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# Tate cycles on some unitary Shimura varieties mod 

David Helm, Yichao Tian and Liang Xiao

Let $F$ be a real quadratic field in which a fixed prime $p$ is inert, and $E_{0}$ be an imaginary quadratic field in which $p$ splits; put $E=E_{0} F$. Let $X$ be the fiber over $\mathbb{F}_{p^{2}}$ of the Shimura variety for $G(U(1, n-1) \times U(n-1,1))$ with hyperspecial level structure at $p$ for some integer $n \geq 2$. We show that under some genericity conditions the middle-dimensional Tate classes of $X$ are generated by the irreducible components of its supersingular locus. We also discuss a general conjecture regarding special cycles on the special fibers of unitary Shimura varieties, and on their relation to Newton stratification.

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## 1. Introduction

The study of the geometry of Shimura varieties lies at the heart of the Langlands program. Arithmetic information of Shimura varieties builds a bridge relating the world of automorphic representations and the world of Galois representations.

One of the interesting topics in this area is to understand the supersingular locus of the special fibers of Shimura varieties, or more generally, any interesting stratifications (e.g., Newton or Ekedahl-Oort stratification) of the special fibers of Shimura varieties. The case of unitary Shimura varieties has been extensively

[^0]studied. Vollaard and Wedhorn [2011] showed that the supersingular locus of the special fiber of the $G U(1, n-1)$-Shimura variety at an inert prime is a union of Deligne-Lusztig varieties. Further, Howard and Pappas [2014] studied the case of $G U(2,2)$ at an inert prime, and Rapoport, Terstiege and Wilson proved similar results for $G U(n-1,1)$ at a ramified prime. Finally, we remark that Görtz and He [2015] studied the basic loci in a slightly more general class of Shimura varieties.

In all the work mentioned above, the authors use the uniformization theorem of Rapoport-Zink to reduce the problem to the study of certain Rapoport-Zink spaces. In this paper, we take a different approach. Instead of using the uniformization theorem, we study the basic locus (or more generally other Newton strata) of certain unitary Shimura varieties by considering correspondences between unitary Shimura varieties of different signatures. This method was introduced by the first author in [Helm 2010; 2012], and applied successfully to quaternionic Shimura varieties by the second and the third authors [Tian and Xiao 2016].

Another new aspect of this work is that we study not only the global geometry of the supersingular locus, but also their relationship with the Tate conjecture for Shimura varieties over finite fields. We show that the basic locus contributes to all "generic" middle-dimensional Tate cycles of the special fiber of the Shimura variety. Similar results have been obtained by the second and the third authors for even-dimensional Hilbert modular varieties at an inert prime [Tian and Xiao 2014]. We believe that, this phenomenon is a general philosophy which holds for more general Shimura varieties. Our slogan is: irreducible components of the basic locus of a Shimura variety should generate all Tate classes under some genericity condition on the automorphic representations.

We explain in more detail the main results of this paper. Let $F$ be a real quadratic field, $E_{0}$ be an imaginary quadratic field, and $E=E_{0} F$. Let $p$ be a prime number inert in $F$, and split in $E_{0}$. Let $\mathfrak{p}, \overline{\mathfrak{p}}$ denote the two places of $E$ above $p$ so that $E_{\mathfrak{p}}$ and $E_{\overline{\mathfrak{p}}}$ are both isomorphic to $\mathbb{Q}_{p^{2}}$, the unique unramified quadratic extension of $\mathbb{Q}_{p}$. For an integer $n \geq 1$, let $G$ be the similitude unitary group associated to a division algebra over $E$ equipped with an involution of second kind. In the notation of Section 3.6, our $G$ is denoted $G_{1, n-1}$. This is an algebraic group over $\mathbb{Q}$ such that $G\left(\mathbb{Q}_{p}\right) \simeq$ $\mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$ and $G(\mathbb{R})$ is the unitary similitude group with signature $(1, n-1)$ and $(n-1,1)$ at the two archimedean places. (For a precise definition, see Section 2.2.)

Let $\mathbb{A}$ denote the ring of finite adeles of $\mathbb{Q}$, and $\mathbb{A}^{\infty}$ be its finite part. Fix a sufficiently small open compact subgroup $K \subseteq G\left(\mathbb{A}^{\infty}\right)$ with $K_{p}=\mathbb{Z}_{p}^{\times} \times \mathrm{GL}_{n}\left(\mathbb{Z}_{p^{2}}\right) \subseteq$ $G\left(\mathbb{Q}_{p}\right)$, where $\mathbb{Z}_{p^{2}}$ is the ring of integers of $\mathbb{Q}_{p^{2}}$. Let $\operatorname{Sh}(G)_{K}$ be the Shimura variety associated to $G$ of level $K .{ }^{1}$

[^1]According to Kottwitz [1992b], when $K^{p}$ is neat, $\operatorname{Sh}(G)_{K}$ admits a proper and smooth integral model over $\mathbb{Z}_{p^{2}}$ which parametrizes certain polarized abelian schemes with $K$-level structure (See Section 2.3). Let $\mathrm{Sh}_{1, n-1}$ denote the special fiber of $\operatorname{Sh}(G)_{K}$ over $\mathbb{F}_{p^{2}}$. This is a proper smooth variety over $\mathbb{F}_{p^{2}}$ of dimension $2(n-1)$. Let $\mathrm{Sh}_{1, n-1}^{\mathrm{ss}}$ denote the supersingular locus of $\mathrm{Sh}_{1, n-1}$, i.e., the reduced closed subvariety of $\mathrm{Sh}_{1, n-1}$ that parametrizes supersingular abelian varieties. We will see in Proposition 4.14 that $\mathrm{Sh}_{1, n-1}^{\mathrm{ss}}$ is equidimensional of dimension $n-1$.

Fix a prime $\ell \neq p$. There is a natural action by $\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{2}}\right) \times \overline{\mathbb{Q}}_{\ell}\left[K \backslash G\left(\mathbb{A}^{\infty}\right) / K\right]$ on the $\ell$-adic étale cohomology group $H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)$. We will take advantage of the Hecke action to consider a variant of the Tate conjecture for $\mathrm{Sh}_{1, n-1}$.

Fix an irreducible admissible representation $\pi$ of $G\left(\mathbb{A}^{\infty}\right)$ (with coefficients in $\overline{\mathbb{Q}}_{\ell}$ ). The $K$-invariant subspace of $\pi$, denoted by $\pi^{K}$, is a finite-dimensional irreducible representation of the Hecke algebra $\overline{\mathbb{Q}}_{\ell}\left[K \backslash G\left(\mathbb{A}^{\infty}\right) / K\right]$. We denote the $\pi^{K}$-isotypic component of $H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \bar{F}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)$ by $H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \bar{F}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}$ and put
$H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{fin}}:=\bigcup_{\mathbb{F}_{q} / \mathbb{F}_{p^{2}}} H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{q}\right)}$,
where $\mathbb{F}_{q}$ runs through all finite extensions of $\mathbb{F}_{p^{2}}$. By projecting to the $\pi^{K}$-isotypic component, we have an $\ell$-adic cycle class map:

$$
\begin{equation*}
\operatorname{cl}_{\pi}^{n-1}: A^{n-1}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}\right) \otimes_{\mathbb{Z}} \overline{\mathbb{Q}}_{\ell} \rightarrow H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{fin}} \tag{1.0.1}
\end{equation*}
$$

where $A^{n-1}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}\right)$ is the abelian group of codimension $n-1$ algebraic cycles on $\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}$. Then the Tate conjecture for $\mathrm{Sh}_{1, n-1}$ predicts that the above map is surjective. Our main result confirms exactly this statement under some "genericity" assumptions on $\pi$.

From now on, we assume that $\pi$ satisfies Hypothesis 2.5 to ensure the nontriviality of the $\pi$-isotypic component of the cohomology groups. In particular, $\pi$ is the finite part of an automorphic cuspidal representation of $G(\mathbb{A})$, and $H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi} \neq 0$. Let $\pi_{p}$ denote the $p$-component of $\pi$, which is an unramified principal series as $K_{p}$ is hyperspecial. Since $G\left(\mathbb{Q}_{p}\right) \simeq \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$, we write $\pi_{p}=\pi_{p, 0} \otimes \pi_{\mathfrak{p}}$, where $\pi_{p, 0}$ is a character of $\mathbb{Q}_{p}^{\times}$and $\pi_{\mathfrak{p}}$ is an irreducible admissible representation of $\mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$.

Our main theorem is the following.
Theorem 1.1. Suppose $\pi$ is the finite part of an automorphic representation of $G(\mathbb{A})$ that admits a cuspidal base change to $\mathrm{GL}_{n}\left(\mathbb{A}_{E}\right) \times \mathbb{A}_{E_{0}}^{\times}$, and the Satake parameters of $\pi_{\mathfrak{p}}$ are distinct modulo roots of unity. Then $H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{fin}}$ is generated by the cohomological classes of the irreducible components of the supersingular locus $\mathrm{Sh}_{1, n-1}^{\mathrm{ss}}$. In particular, the cycle class map (1.0.1) is surjective.

This theorem will be restated in a more precise form in Theorem 4.18. Here, the assumption that the Satake parameters of $\pi_{\mathfrak{p}}$ are distinct modulo roots of unity is crucial for our method. It is closely tied to our geometric description of the irreducible components. This condition will be reformulated in Theorem 4.18 in terms of the Frobenius eigenvalues of certain Galois representation attached to $\pi_{\mathfrak{p}}$ via the unramified local Langlands correspondence. The other automorphic assumption on $\pi$ is of technical nature. It is imposed here to ensure certain equalities on the automorphic multiplicity on $\pi$ (See Remark 4.19). The method of our paper may be extended to more general representations $\pi$ if we have more knowledge of the multiplicity of automorphic forms on unitary groups.

What we will prove is more precise than stated in Theorem 1.1. We need another unitary group $G^{\prime}=G_{0, n}$ over $\mathbb{Q}$ for $E / F$ as in Lemma 2.9 , which is the unique inner form of $G$ such that $G^{\prime}\left(\mathbb{A}^{\infty}\right) \simeq G\left(\mathbb{A}^{\infty}\right)$ and the signatures of $G^{\prime}$ at the two archimedean places are $(0, n)$ and $(n, 0)$. Let $\mathrm{Sh}_{0, n}$ denote the (zero-dimensional) Shimura variety over $\mathbb{F}_{p^{2}}$ associated to $G^{\prime}$. We will show in Proposition 4.14 that the supersingular locus $\mathrm{Sh}_{1, n-1}^{\mathrm{ss}}$ is a union of $n$ closed subvarieties $Y_{j}$ with $1 \leq j \leq n$ such that each of $Y_{j}$ admits a fibration over $\mathrm{Sh}_{0, n}$ of the same level $K \subseteq G\left(\mathbb{A}^{\infty}\right) \simeq G^{\prime}\left(\mathbb{A}^{\infty}\right)$ with fibers isomorphic to a certain proper and smooth closed subvariety in a product of Grassmannians. In other words, each $Y_{j}$ is an algebraic correspondence between $\mathrm{Sh}_{1, n-1}$ and $\mathrm{Sh}_{0, n}$ :

$$
\mathrm{Sh}_{0, n} \leftarrow Y_{j} \rightarrow \mathrm{Sh}_{1, n-1}
$$

This can be viewed as a geometric realization of the Jacquet-Langlands correspondence between $G$ and $G^{\prime}$ in the sense of [Helm 2010]. Alternatively, we may view these $Y_{j}$ as Hecke correspondences between special fibers of unitary Shimura varieties of different signatures. To prove Theorem 1.1, it suffices to show that, when the Satake parameters of $\pi_{\mathfrak{p}}$ are distinct modulo roots of unity, $H_{\mathrm{et}}^{2(n-1)}\left(\mathrm{Sh}_{1, n-1, \overline{\mathbb{F}}_{p}}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{fin}}$ is generated by the cohomology classes of the irreducible components of $Y_{j}$. The key point is to show that the $\pi$-projection of the intersection matrix of $Y_{j}$ is nondegenerate under the assumption above on $\pi_{\mathfrak{p}}$.

We briefly describe the structure of this paper. In Section 2, we consider a more general setup of unitary Shimura varieties, and propose a general conjecture, which roughly predicts the existence of certain algebraic correspondences between the special fibers of Shimura varieties with hyperspecial level at $p$ associated to unitary groups with different signatures at infinity (Conjecture 2.12). Theorem 1.1 is a special case of Conjecture 2.12 . We believe that our conjecture will provide a new perspective to understand the special fibers of Shimura varieties. In Section 3, we review some Dieudonné theory and Grothendieck-Messing deformation theory that will be frequently used in later sections. Section 4 is devoted to the study of the supersingular locus $\mathrm{Sh}_{1, n-1}^{\text {ss }}$, and constructing the subvarieties $Y_{j}$ mentioned
above. In Section 5, we compute certain intersection numbers on products of Grassmannian varieties. These numbers will play a fundamental role in our later computation of the intersection matrix of the $Y_{j}$. In Section 6, we will compute explicitly the intersection matrix of the $Y_{j}$ (Theorem 6.7), and show that its $\pi$ isotypic projection of the intersection matrix is nondegenerate as long as the Satake parameters of $\pi_{\mathfrak{p}}$ are distinct (as opposed to being distinct modulo roots of unity). Then an easy cohomological computation allows us to conclude the proof of our main theorem. In Section 7, we will generalize the construction of the cycles $Y_{j}$ to the Shimura variety associated to unitary group for $E / F$ of signature $(r, s) \times(s, r)$ at infinity. In this case, we only obtain some partial results on these cycles predicted by Conjecture 2.12: the union of these cycles is exactly the supersingular locus of the unitary Shimura variety in question (Theorem 7.8).

## 2. The conjecture on special cycles

We will only discuss certain unitary Shimura varieties so that the description becomes explicit. We will discuss after Conjecture 2.12 on how to possibly extend this conjecture to more general Shimura varieties.
2.1. Notation. We fix a prime number $p$ throughout this paper. We fix an isomorphism $\iota_{p}: \mathbb{C} \xrightarrow{\sim} \overline{\mathbb{Q}}_{p}$. Let $\mathbb{Q}_{p}^{\text {ur }}$ be the maximal unramified extension of $\mathbb{Q}_{p}$ inside $\overline{\mathbb{Q}}_{p}$.

Let $F$ be a totally real field of degree $f$ in which $p$ is inert. We label all real embeddings of $F$, or equivalently (via $\iota_{p}$ ), all $p$-adic embeddings of $F$ (into $\mathbb{Q}_{p}^{\mathrm{ur}}$ ) by $\tau_{1}, \ldots, \tau_{f}$ so that post-composition by the Frobenius map takes $\tau_{i}$ to $\tau_{i+1}$. Here the subindices are taken modulo $f$. Let $E_{0}$ be an imaginary quadratic extension of $\mathbb{Q}$ in which $p$ splits. Put $E=E_{0} F$. Denote by $v$ and $\bar{v}$ the two $p$-adic places of $E_{0}$. Then $p$ splits into two primes $\mathfrak{p}$ and $\overline{\mathfrak{p}}$ in $E$, where $\mathfrak{p}$ (resp. $\overline{\mathfrak{p}}$ ) is the $p$-adic place above $v$ (resp. $\bar{v}$ ). Let $q_{i}$ denote the embedding $E \rightarrow E_{\mathfrak{p}} \cong F_{p} \xrightarrow{\tau_{i}} \overline{\mathbb{Q}}_{p}$ and $\bar{q}_{i}$ the analogous embedding which factors through $E_{\bar{p}}$ instead. Composing with $\iota_{p}^{-1}$, we regard $q_{i}$ and $\bar{q}_{i}$ as complex embeddings of $E$, and we put $\Sigma_{\infty, E}=\left\{q_{1}, \ldots, q_{f}, \bar{q}_{1}, \ldots, \bar{q}_{f}\right\}$.
2.2. Shimura data. Let $D$ be a division algebra of dimension $n^{2}$ over its center $E$, equipped with a positive involution $*$ which restricts to the complex conjugation $c$ on $E$. In particular, $D^{\text {opp }} \cong D \otimes_{E, c} E$. We assume that $D$ splits at $\mathfrak{p}$ and $\overline{\mathfrak{p}}$, and we fix an isomorphism

$$
D \otimes_{\mathbb{Q}} \mathbb{Q}_{p} \simeq \mathrm{M}_{n}\left(E_{\mathfrak{p}}\right) \times \mathrm{M}_{n}\left(E_{\overline{\mathfrak{p}}}\right) \cong \mathrm{M}_{n}\left(\mathbb{Q}_{p^{f}}\right) \times \mathrm{M}_{n}\left(\mathbb{Q}_{p^{f}}\right),
$$

where $*$ switches the two direct factors. We use $\mathfrak{e}$ to denote the element of $D \otimes_{\mathbb{Q}} \mathbb{Q}_{p}$ corresponding to the $(1,1)$-elementary matrix ${ }^{2}$ in the first factor. Let $a_{0}=\left(a_{i}\right)_{1 \leq i \leq f}$

[^2]be a tuple of $f$ numbers with $a_{i} \in\{0, \ldots, n\}$. Assume that there exists an element $\beta_{a_{0}} \in\left(D^{\times}\right)^{*=-1}$ such that the following condition is satisfied: ${ }^{3}$

Let $G_{a_{0}}$ be the algebraic group over $\mathbb{Q}$ such that $G_{a_{\bullet}}(R)$ for a $\mathbb{Q}$-algebra $R$ consists of elements $g \in\left(D^{\mathrm{opp}} \otimes_{\mathbb{Q}} R\right)^{\times}$with $g \beta_{a_{0}} g^{*}=c(g) \beta_{a_{0}}$ for some $c(g) \in R^{\times}$. If $G_{a_{\bullet}}^{1}$ denotes the kernel of the similitude character $c: G_{a_{\bullet}} \rightarrow \mathbb{G}_{m, \mathbb{Q}}$, then there exists an isomorphism

$$
G_{a_{\bullet}}^{1}(\mathbb{R}) \simeq \prod_{i=1}^{f} U\left(a_{i}, n-a_{i}\right),
$$

where the $i$-th factor corresponds to the real embedding $\tau_{i}: F \hookrightarrow \mathbb{R}$.
Note that the assumption on $D$ at $p$ implies that

$$
G_{a \cdot}\left(\mathbb{Q}_{p}\right) \simeq \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right) \cong \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(\mathbb{Q}_{p} f\right)
$$

We put $V_{a_{\bullet}}=D$ and view it as a left $D$-module. Let $\langle-,-\rangle_{a_{0}}: V_{a_{0}} \times V_{a_{0}} \rightarrow \mathbb{Q}$ be the perfect alternating pairing given by

$$
\langle x, y\rangle_{a_{0}}=\operatorname{Tr}_{D / \mathbb{Q}}\left(x \beta_{a_{0}} y^{*}\right) \quad \text { for } x, y \in V_{a_{\bullet}} .
$$

Then $G_{a_{0}}$ is identified with the similitude group associated to $\left(V_{a_{\bullet}},\langle-,-\rangle_{a_{0}}\right)$, i.e., for all $\mathbb{Q}$-algebra $R$, we have
$G_{a_{0}}(R)=\left\{g \in \operatorname{End}_{D \otimes_{\mathbb{Q}} R}\left(V_{a_{0}} \otimes_{\mathbb{Q}} R\right) \mid\langle g x, g y\rangle_{a_{0}}=c(g)\langle x, y\rangle_{a_{0}}\right.$ for some $\left.c(g) \in R^{\times}\right\}$.
Consider the homomorphism of $\mathbb{R}$-algebraic groups $h: \operatorname{Res}_{\mathbb{C} / \mathbb{R}}\left(\mathbb{G}_{m}\right) \rightarrow G_{a_{\mathbf{0}}, \mathbb{R}}$ given by

$$
\begin{equation*}
h(z)=\prod_{i=1}^{f} \operatorname{Diag}(\underbrace{z, \ldots, z}_{a_{i}}, \underbrace{\bar{z}, \ldots, \bar{z}}_{n-a_{i}}), \quad \text { for } z=x+\sqrt{-1} y . \tag{2.2.1}
\end{equation*}
$$

Let $\mu_{h}: \mathbb{G}_{m, \mathbb{C}} \rightarrow G_{a_{\bullet}, \mathbb{C}}$ be the composite of $h_{\mathbb{C}}$ with the map

$$
\mathbb{G}_{m, \mathbb{C}} \rightarrow \operatorname{Res}_{\mathbb{C} / \mathbb{R}}\left(\mathbb{G}_{m}\right)_{\mathbb{C}} \cong \mathbb{C}^{\times} \times \mathbb{C}^{\times}, \quad z \mapsto(z, 1)
$$

Here, the first copy of $\mathbb{C}^{\times}$in $\operatorname{Res}_{\mathbb{C} / \mathbb{R}}\left(\mathbb{G}_{m}\right)$ is the one indexed by the identity element in $\operatorname{Aut}_{\mathbb{R}}(\mathbb{C})$, and the other copy of $\mathbb{C}^{\times}$is indexed by the complex conjugation.

Let $E_{h}$ be the reflex field of $\mu_{h}$, i.e., the minimal subfield of $\mathbb{C}$ where the conjugacy class of $\mu_{h}$ is defined. It has the following explicit description. The group $\operatorname{Aut}_{\mathbb{Q}}(\mathbb{C})$ acts naturally on $\Sigma_{\infty, E}$, and hence on the functions on $\Sigma_{\infty, E}$. Then $E_{h}$ is the subfield of $\mathbb{C}$ fixed by the stabilizer of the $\mathbb{Z}$-valued function $a$ on $\Sigma_{\infty, E}$ defined by $a\left(q_{i}\right)=a_{i}$ and $a\left(\bar{q}_{i}\right)=n-a_{i}$. The isomorphism $\iota_{p}: \mathbb{C} \xrightarrow{\sim} \overline{\mathbb{Q}}_{p}$

[^3]defines a $p$-adic place $\wp$ of $E_{h}$. By our hypothesis on $E$, the local field $E_{h, \wp}$ is an unramified extension of $\mathbb{Q}_{p}$ contained in $\mathbb{Q}_{p^{f}}$, the unique unramified extension over $\mathbb{Q}_{p}$ of degree $f$.
2.3. Unitary Shimura varieties of PEL-type. Let $\mathcal{O}_{D}$ be a $*$-stable order of $D$ and $\Lambda_{a_{0}}$ an $\mathcal{O}_{D}$-lattice of $V_{a_{0}}$ such that $\left\langle\Lambda_{a_{\bullet}}, \Lambda_{a_{\bullet}}\right\rangle_{a_{\bullet}} \subseteq \mathbb{Z}$ and $\Lambda_{a_{\bullet}} \otimes_{\mathbb{Z}} \mathbb{Z}_{p}$ is self-dual under the alternating pairing induced by $\langle-,-\rangle_{a_{\bullet}}$. We put $K_{p}=\mathbb{Z}_{p}^{\times} \times \mathrm{GL}_{n}\left(\mathcal{O}_{E_{\mathfrak{p}}}\right) \subseteq$ $G_{a_{\mathbf{e}}}\left(\mathbb{Q}_{p}\right)$, and fix an open compact subgroup $K^{p} \subseteq G_{a_{\mathbf{0}}}\left(\mathbb{A}^{\infty, p}\right)$ such that $K=K^{p} K_{p}$ is neat, i.e., $G_{a_{0}}(\mathbb{Q}) \cap g K g^{-1}$ is torsion free for any $g \in G_{a_{0}}\left(\mathbb{A}^{\infty}\right)$.

Following [Kottwitz 1992b], we have a unitary Shimura variety $\mathcal{S} h_{a_{0}}$ defined over $\mathbb{Z}_{p^{f}} ;{ }^{4}$ it represents the functor that takes a locally noetherian $\mathbb{Z}_{p^{f-}}$-scheme $S$ to the set of isomorphism classes of tuples $(A, \lambda, \eta)$, where
(1) $A$ is an $f n^{2}$-dimensional abelian variety over $S$ equipped with an action of $\mathcal{O}_{D}$ such that the induced action on $\operatorname{Lie}(A / S)$ satisfies the Kottwitz determinant condition, that is, if we view the reduced relative de Rham homology $H_{1}^{\mathrm{dR}}(A / S)^{\circ}:=\mathfrak{e} H_{1}^{\mathrm{dR}}(A / S)$ and its quotient $\mathrm{Lie}_{A / S}^{\circ}:=\mathfrak{e} \cdot \mathrm{Lie}_{A / S}$ as a module over $F_{p} \otimes_{\mathbb{Z}_{p}} \mathcal{O}_{S} \cong \bigoplus_{i=1}^{f} \mathcal{O}_{S}$, they, respectively, decompose into the direct sums of locally free $\mathcal{O}_{S}$-modules $H_{1}^{\mathrm{dR}}(A / S)_{i}^{\circ}$ of rank $n$ and, their quotients, locally free $\mathcal{O}_{S}$-modules $\mathrm{Lie}_{A / S, i}^{\circ}$ of rank $n-a_{i}$;
(2) $\lambda: A \rightarrow A^{\vee}$ is a prime-to- $p \mathcal{O}_{D}$-equivariant polarization such that the Rosati involution induces the involution $*$ on $\mathcal{O}_{D}$;
(3) $\eta$ is a collection of, for each connected component $S_{j}$ of $S$ with a geometric point $\bar{s}_{j}$, a $\pi_{1}\left(S_{j}, \bar{s}_{j}\right)$-invariant $K^{p}$-orbit of isomorphisms $\eta_{j}: \Lambda_{a} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}^{(p)} \simeq$ $T^{(p)}\left(A_{\bar{s}_{j}}\right)$ such that the following diagram commutes for an isomorphism $v\left(\eta_{j}\right) \in \operatorname{Hom}\left(\widehat{\mathbb{Z}}^{(p)}, \widehat{\mathbb{Z}}^{(p)}(1)\right)$ :

where $\widehat{\mathbb{Z}}^{(p)}=\prod_{\ell \neq p} \mathbb{Z}_{\ell}$ and $T^{(p)}\left(A_{\bar{s}_{j}}\right)$ denotes the product of the $\ell$-adic Tate modules of $A_{\bar{s}_{j}}$ for all $\ell \neq p$.

The Shimura variety $\mathcal{S} h_{a_{0}}$ is smooth and projective over $\mathbb{Z}_{p^{f}}$ of relative dimension $d\left(a_{0}\right):=\sum_{i=1}^{f} a_{i}\left(n-a_{i}\right)$. Note that if $a_{i} \in\{0, n\}$ for all $i$, then $\mathcal{S} h_{a_{0}}$ is of relative dimension zero; we call it a discrete Shimura variety.

[^4]We denote by $\mathcal{S} h_{a_{\bullet}}(\mathbb{C})$ the complex points of $\mathcal{S} h_{a_{\bullet}}$ via the embedding

$$
\mathbb{Z}_{p^{f}} \hookrightarrow \overline{\mathbb{Q}}_{p} \xrightarrow{{l_{p}^{-1}}_{\mathbb{C}} . . . . ~}
$$

Let $K_{\infty} \subseteq G_{a_{0}}(\mathbb{R})$ be the stabilizer of $h(2.2 .1)$ under the conjugation action, and let $X_{\infty}$ denote the $G_{a_{0}}(\mathbb{R})$-conjugacy class of $h$. Then $K_{\infty}$ is a maximal compact-modulo-center subgroup of $G_{a_{\bullet}}(\mathbb{R})$. According to [Kottwitz 1992b, page 400], the complex manifold $S h_{a_{\bullet}}(\mathbb{C})$ is the disjoint union of $\# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{a_{0}}\right)$ copies of

$$
\begin{equation*}
G_{a_{\mathbf{0}}}(\mathbb{Q}) \backslash\left(G_{a_{\mathbf{0}}}\left(\mathbb{A}^{\infty}\right) \times X_{\infty}\right) / K \cong G_{a_{\mathbf{0}}}(\mathbb{Q}) \backslash G_{a_{\mathbf{0}}}(\mathbb{A}) / K \times K_{\infty} \tag{2.3.1}
\end{equation*}
$$

Here, if $n$ is even, then $\operatorname{ker}^{1}\left(\mathbb{Q}, G_{a_{0}}\right)=(0)$, while if $n$ is odd then

$$
\operatorname{ker}^{1}\left(\mathbb{Q}, G_{a_{0}}\right)=\operatorname{Ker}\left(F^{\times} / \mathbb{Q}^{\times} N_{E / F}\left(E^{\times}\right) \rightarrow \mathbb{A}_{F}^{\times} / \mathbb{A}^{\times} N_{E / F}\left(\mathbb{A}_{E}^{\times}\right)\right)
$$

In either case, $\operatorname{ker}^{1}\left(\mathbb{Q}, G_{a_{0}}\right)$ depends only on the CM extension $E / F$ and the parity of $n$ but not on the tuple $a_{\text {. }}$.

Let $\mathrm{Sh}_{a_{0}}:=\mathcal{S} h_{a_{0}} \otimes_{\mathbb{Z}_{p f} f} \mathbb{F}_{p^{f}}$ denote the special fiber of $\mathcal{S} h_{a_{0}}$, and let $\overline{\mathrm{Sh}}_{a_{0}}:=$ $\mathrm{Sh}_{a_{0}} \otimes_{\mathbb{F}_{p f} f} \overline{\mathbb{F}}_{p}$ denote the geometric special fiber.
2.4. $\ell$-adic cohomology. We fix a prime number $\ell \neq p$, and an isomorphism $\iota_{\ell}: \mathbb{C} \simeq \overline{\mathbb{Q}}_{\ell}$. Let $\xi$ be an algebraic representation of $G_{a_{0}}$ over $\overline{\mathbb{Q}}_{\ell}$, and $\xi_{\mathbb{C}}$ be the base change via $\iota_{\ell}^{-1}$. The theory of automorphic sheaves [Milne 1990, Section III] or just reading off from the rational $\ell$-adic Tate modules of the universal abelian variety allows us to attach to $\xi$ a lisse $\overline{\mathbb{Q}}_{\ell}$-sheaf $\mathcal{L}_{\xi}$ over $\mathcal{S} h_{a_{0} .}$. For example, if $\xi$ is the representation of $G_{a_{0}}$ on the vector space $V_{a_{0}}$ (Section 2.2), the corresponding $\ell$-adic local system is given by the rational $\ell$-adic Tate module (tensored with $\overline{\mathbb{Q}}_{\ell}$ ) of the universal abelian scheme over $\mathcal{S} h_{a_{0}}$.

We assume that $\xi$ is irreducible. Let $\mathscr{H}_{K}=\mathscr{H}\left(K, \overline{\mathbb{Q}}_{\ell}\right)$ be the Hecke algebra of compactly supported $K$-bi-invariant $\overline{\mathbb{Q}}_{\ell}$-valued functions on $G_{a_{0}}\left(\mathbb{A}^{\infty}\right)$. The étale cohomology group $H_{\mathrm{et}}^{d\left(a_{\bullet}\right)}\left(\overline{\mathrm{Sh}}_{a_{\bullet}}, \mathcal{L}_{\xi}\right)$ is equipped with a natural action of $\mathscr{H}_{K} \times \operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{f}}\right)$. Since $\mathrm{Sh}_{a_{0}}$ is proper and smooth, there is no continuous spectrum and we have a canonical decomposition of $\mathscr{H}_{K} \times \operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{f}}\right)$-modules (see, e.g., [Harris and Taylor 2001, Proposition III.2.1])

$$
\begin{equation*}
H_{\mathrm{et}}^{d\left(a_{\bullet}\right)}\left(\overline{\mathrm{Sh}}_{a_{\bullet}}, \mathcal{L}_{\xi}\right)=\bigoplus_{\pi \in \operatorname{Irr}\left(G_{a_{\bullet}}\left(\mathrm{A}^{\infty}\right)\right)} \iota_{\ell}\left(\pi^{K}\right) \otimes R_{a_{\bullet}, \ell}(\pi), \tag{2.4.1}
\end{equation*}
$$

where $\operatorname{Irr}\left(G_{a_{0}}\left(\mathbb{A}^{\infty}\right)\right)$ is the set of irreducible admissible representations of $G_{a_{0}}\left(\mathbb{A}^{\infty}\right)$ with coefficients in $\mathbb{C}, \pi^{K}$ is the $K$-invariant subspace of $\pi \in \operatorname{Irr}\left(G_{a_{0}}\left(\mathbb{A}^{\infty}\right)\right)$ and $R_{a_{\bullet} \ell}(\pi)$ is a certain $\ell$-adic representation of $\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{f}}\right)$ which we specify below.

We write $H_{\mathrm{et}}^{d\left(a_{\mathbf{0}}\right)}\left(\overline{\mathrm{Sh}}_{a_{\mathbf{t}}}, \mathcal{L}_{\xi}\right)_{\pi}$ for the $\pi$-isotypic component of the cohomology group, that is, the direct summand of (2.4.1) labeled by $\pi$. We make the following assumptions on $\pi$.

Hypothesis 2.5. (1) We have $\pi^{K} \neq 0$.
(2) There exists an admissible irreducible representation $\pi_{\infty}$ of $G_{a_{\mathbf{0}}}(\mathbb{R})$ such that $\pi \otimes \pi_{\infty}$ is a cuspidal automorphic representation of $G_{a_{0}}(\mathbb{A})$,
(2a) $\pi_{\infty}$ is cohomological in degree $d\left(a_{0}\right)$ for $\xi$ in the sense that

$$
\begin{equation*}
H^{d\left(a_{\bullet}\right)}\left(\operatorname{Lie}\left(G_{a_{\mathbf{0}}}(\mathbb{R})\right), K_{\infty}, \pi_{\infty} \otimes \xi_{\mathbb{C}}\right) \neq 0,{ }^{5} \tag{2.5.1}
\end{equation*}
$$

where $K_{\infty}$ is a maximal compact subgroup of $G_{a_{0}}(\mathbb{R})$,
(2b) and $\pi \otimes \pi_{\infty}$ admits a base change to a cuspidal automorphic representation of $\mathrm{GL}_{n}\left(\mathbb{A}_{E}\right) \times \mathbb{A}_{E_{0}}^{\times}$.
Note that Hypothesis $2.5(1)$ implies that the $p$-component $\pi_{p}$ is unramified. Hypothesis 2.5 (2a) ensures that $R_{a_{\bullet}, \ell}(\pi)$ is nontrivial. Moreover, by [Caraiani 2012, Theorem 1.2], this hypothesis implies that the base change of $\pi \otimes \pi_{\infty}$ to $\mathrm{GL}_{n, E}$ is tempered at all finite places, and hence $\pi_{p}$ is tempered.

We recall now an explicit description, due to Kottwitz [1992a], of the Galois module $R_{a_{\bullet}, \ell}(\pi)$. As $G_{a_{\bullet}}\left(\mathbb{Q}_{p}\right)=\mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$, we may write $\pi_{p}=\pi_{p, 0} \otimes \pi_{\mathfrak{p}}$, where $\pi_{p, 0}$ is a character of $\mathbb{Q}_{p}^{\times}$trivial on $\mathbb{Z}_{p}^{\times}$, and $\pi_{\mathfrak{p}}$ is an irreducible admissible representation of $\mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$ such that $\pi_{\mathfrak{p}}^{\mathrm{GL}_{n}\left(\mathcal{O}_{E_{\mathfrak{p}}}\right)} \neq 0$. Choose a square root $\sqrt{p}$ of $p$ in $\overline{\mathbb{Q}}$. Depending on this choice of $\sqrt{p}$, one has an (unramified) local Langlands parameter attached to $\pi_{p}$ :

$$
\varphi_{\pi_{p}}=\left(\varphi_{\pi_{p, 0}}, \varphi_{\pi_{\mathfrak{p}}}\right): W_{\mathbb{Q}_{p}} \rightarrow^{L}\left(G_{a, \mathbb{Q}_{p}}\right) \simeq \mathbb{C}^{\times} \times\left(\mathrm{GL}_{n}(\mathbb{C})^{\mathbb{Z} / f \mathbb{Z}} \rtimes \operatorname{Gal}\left(\overline{\mathbb{Q}}_{p} / \mathbb{Q}_{p}\right)\right)
$$

Here, $W_{\mathbb{Q}_{p}}$ is the Weil group of $\mathbb{Q}_{p}$, and $\operatorname{Gal}\left(\overline{\mathbb{Q}}_{p} / \mathbb{Q}_{p}\right)$ permutes cyclically the $f$ copies of $\mathrm{GL}_{n}(\mathbb{C})$ though the quotient $\operatorname{Gal}\left(\mathbb{Q}_{p^{f}} / \mathbb{Q}_{p}\right) \cong \mathbb{Z} / f \mathbb{Z}$. The image of $\varphi_{\pi_{p}} \mid W_{Q_{p}}$ lies in $\left({ }^{L} G_{a_{0}}\right)^{\circ} \simeq \mathbb{C}^{\times} \times \mathrm{GL}_{n}(\mathbb{C})^{\mathbb{Z}} / f \mathbb{Z}$. The cocharacter $\mu_{h}: \mathbb{G}_{m, E_{h}} \rightarrow G_{a_{\mathbf{\bullet}}, E_{h}}$ induces a character $\check{\mu}_{h}$ of $\left({ }^{L} G_{a_{0}}\right)^{\circ}$ over $E_{h}$. Let $r_{\mu_{h}}$ denote the algebraic representation of $\left({ }^{L} G_{a_{0}}\right)^{\circ}$ with extreme weight $\check{\mu}_{h}$. Denote by $\operatorname{Frob}_{p^{f}}$ a geometric Frobenius element in $W_{\mathbb{Q}_{p} f}$. Let $\overline{\mathbb{Q}}_{\ell}(1 / 2)$ denote the unramified representation of $W_{\mathbb{Q}_{p} f}$ which sends $\operatorname{Frob}_{p^{f}}$ to $(\sqrt{p})^{-f}$. Then $R_{a_{\bullet} \ell \ell}(\pi)$ can be described in terms of $\varphi_{\pi_{p}}$ as follows.
Theorem 2.6 [Kottwitz 1992a, Theorem 1]. Under the hypothesis and notation above, we have an equality in the Grothendieck group of $W_{\mathbb{Q}_{p} f}$-modules:

$$
\left[R_{a_{\bullet}, \ell}(\pi)\right]=\# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{a_{\bullet}}\right) m_{a_{\mathbf{\bullet}}}(\pi)\left[\iota_{\ell}\left(r_{\mu_{h}} \circ \varphi_{\pi_{p}}\right) \otimes \overline{\mathbb{Q}}_{\ell}\left(-\frac{1}{2} d\left(a_{\mathbf{\bullet}}\right)\right)\right]
$$

where $m_{a_{\bullet}}(\pi)$ is a certain integer related to the automorphic multiplicities of automorphic representations of $G_{a_{0}}$ with finite part $\pi .^{6}$

[^5]In our case, one can make Kottwitz's theorem more transparent. Define an $\ell$-adic representation

$$
\begin{equation*}
\rho_{\pi_{\mathfrak{p}}}=\iota_{\ell}\left(\varphi_{\pi_{\mathfrak{p}}}^{(1), \vee}\right) \otimes \overline{\mathbb{Q}}_{\ell}\left(\frac{1}{2}(1-n)\right): W_{\mathbb{Q}_{p} f} \rightarrow \mathrm{GL}_{n}\left(\overline{\mathbb{Q}}_{\ell}\right), \tag{2.6.1}
\end{equation*}
$$

where $\varphi_{\pi_{\mathfrak{p}}}^{(1), \vee}: W_{\mathbb{Q}_{p} f} \rightarrow \mathrm{GL}_{n}(\mathbb{C})$ denotes the contragredient of the projection to the first (or any) copy of $\mathrm{GL}_{n}(\mathbb{C})$. Both $\varphi_{\pi_{\mathfrak{p}}}$ and $\overline{\mathbb{Q}}_{\ell}\left(\frac{1}{2}\right)$ depend on the choice of $\sqrt{p}$, but $\rho_{\pi_{\mathfrak{p}}}$ does not. Explicitly, $\rho_{\pi_{\mathfrak{p}}}\left(\operatorname{Frob}_{p^{f}}\right)$ is semisimple with the characteristic polynomial given by [Gross 1998, (6.7)]:

$$
\begin{equation*}
X^{n}+\sum_{i=1}^{n}(-1)^{i}(N \mathfrak{p})^{i(i-1) / 2} a_{\mathfrak{p}}^{(i)} X^{n-i} \tag{2.6.2}
\end{equation*}
$$

where $a_{\mathfrak{p}}^{(i)}$ is the eigenvalue on $\pi_{\mathfrak{p}} \mathrm{GL}_{n}\left(\mathcal{O}_{E_{\mathfrak{p}}}\right)$ of the Hecke operator

$$
T_{\mathfrak{p}}^{(i)}=\mathrm{GL}_{n}\left(\mathcal{O}_{E_{\mathfrak{p}}}\right) \cdot \operatorname{Diag}(\underbrace{p, \ldots, p}_{i}, \underbrace{1, \ldots, 1}_{n-i}) \cdot \mathrm{GL}_{n}\left(\mathcal{O}_{E_{\mathfrak{p}}}\right)
$$

An easy computation shows that $r_{\mu_{h}}=\operatorname{Std}_{\mathbb{Q}_{p}^{\times}}^{-1} \otimes \bigotimes_{i=1}^{f}\left(\wedge^{a_{i}} \operatorname{Std}^{\vee}\right)$. Since the projection of $\left.\varphi_{\pi_{p}}\right|_{W_{Q_{p} f}}$ to each copy of $\mathrm{GL}_{n}(\mathbb{C})$ is conjugate to all others, Theorem 2.6 is equivalent to

$$
\begin{align*}
& {\left[R_{a_{\mathbf{\bullet}}, \ell}(\pi)\right]} \\
& \quad=\# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{a_{\mathbf{\bullet}}}\right) \cdot m_{a_{\mathbf{\bullet}}}(\pi)\left[\rho_{a_{\mathbf{\bullet}}}\left(\pi_{\mathfrak{p}}\right) \otimes \chi_{\pi_{p, 0}}^{-1} \otimes \overline{\mathbb{Q}}_{\ell}\left(\sum_{i} \frac{1}{2} a_{i}\left(a_{i}-1\right)\right)\right] \tag{2.6.3}
\end{align*}
$$

where $\rho_{a_{\bullet}}\left(\pi_{\mathfrak{p}}\right)=r_{a_{\bullet}} \circ \rho_{\pi_{\mathfrak{p}}}$ with $r_{a_{\mathbf{\bullet}}}=\bigotimes_{i=1}^{f} \wedge^{a_{i}}$ Std, and $\chi_{\pi_{p, 0}}$ denotes the character of $\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{f}}\right)$ sending $\operatorname{Frob}_{p^{f}}$ to $\iota_{\ell}\left(\pi_{p, 0}\left(p^{f}\right)\right)$.
Remark 2.7. The reason why we normalize the Galois representation as above is the following: By Hypothesis $2.5, \pi$ is the finite part of an automorphic representation of $G_{a_{\bullet}}$ (A) which admits a base change to a cuspidal automorphic representation $\Pi \otimes \chi$ of $\mathrm{GL}_{n}\left(\mathbb{A}_{E}\right) \times \mathbb{A}_{E_{0}}^{\times}$. If $\rho_{\Pi}$ denotes the Galois representation of $\mathrm{Gal}(\overline{\mathbb{Q}} / E)$ attached to $\Pi$, then $\rho_{\pi_{\mathfrak{p}}}$ is the semisimplification of the restriction of $\rho_{\Pi}$ to $W_{E_{\mathfrak{p}}}$ (See [Caraiani 2012, Theorem 1.1]).
2.8. Tate conjecture. We recall first the Tate conjecture [1966] over finite fields. Let $X$ be a projective smooth variety over a finite field $\mathbb{F}_{q}$ of characteristic $p$. Put $\bar{X}=X_{\overline{\mathbb{F}}_{p}}$. For each prime $\ell \neq p$ and integer $r \leq \operatorname{dim}(X)$, we have a cycle class map

$$
\mathrm{cl}_{X}^{r}: A^{r}(X) \otimes_{\mathbb{Z}} \overline{\mathbb{Q}}_{\ell} \rightarrow H_{\mathrm{et}}^{2 r}\left(\bar{X}, \overline{\mathbb{Q}}_{\ell}(r)\right)^{\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{q}\right)}
$$

where $A^{r}(X)$ denotes the abelian group of codimension $r$ algebraic cycles in $X$ defined over $\mathbb{F}_{q}$. Then the Tate conjecture predicts that this map is surjective. One
has a geometric variant of the Tate conjecture, which claims that the geometric cycle class map:

$$
\operatorname{cl}_{\bar{X}}^{r}: A^{r}(\bar{X}) \otimes_{\mathbb{Z}} \overline{\mathbb{Q}}_{\ell} \rightarrow H_{\mathrm{et}}^{2 r}\left(\bar{X}, \overline{\mathbb{Q}}_{\ell}(r)\right)^{\mathrm{fin}}:=\bigcup_{m \geq 1} H_{\mathrm{et}}^{2 r}\left(\bar{X}, \overline{\mathbb{Q}}_{\ell}(r)\right)^{\operatorname{Gal}\left(\bar{F}_{p} / \mathbb{F}_{q^{\prime}} m\right)}
$$

is surjective. Here, the superscript "fin" means the subspace on which $\operatorname{Gal}\left(\mathbb{F}_{p} / \mathbb{F}_{q}\right)$ acts through a finite quotient. Note that the surjectivity of $\mathrm{cl}_{\bar{X}}^{r}$ implies that of $\mathrm{cl}_{X}^{r}$ by taking the $\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{q}\right)$-invariant subspace.

Consider the case $X=\mathrm{Sh}_{a_{0}}$ with $d\left(a_{0}\right)$ even. Let $\pi$ be an irreducible admissible representation of $G_{a_{0}}\left(\mathbb{A}^{\infty}\right)$ as in Theorem 2.6. By Theorem 2.6, the $\pi$-isotypic component of $H_{\mathrm{et}}^{d\left(a_{\bullet}\right.}{ }^{\circ}\left(\overline{\operatorname{Sh}}_{a_{0}}, \overline{\mathbb{Q}}_{\ell}\left(\frac{1}{2} d\left(a_{\bullet}\right)\right)\right)^{\text {fin }}$ is, up to Frobenius semisimplification ${ }^{7}$, isomorphic to $\operatorname{dim}\left(\pi^{K}\right) \cdot \# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{a_{\bullet}}\right) \cdot m_{a_{\bullet}}(\pi)$ copies of

$$
\begin{equation*}
\left(\rho_{a_{\mathbf{0}}}\left(\pi_{\mathfrak{p}}\right) \otimes \chi_{\pi_{p, 0}}^{-1} \otimes \overline{\mathbb{Q}}_{\ell}\left(\frac{(n-1)}{2} \sum_{i=1}^{f} a_{i}\right)\right)^{\mathrm{fin}} \tag{2.8.1}
\end{equation*}
$$

Note that $\chi_{\pi_{p, 0}}\left(\operatorname{Frob}_{p^{f}}\right)=\pi_{p, 0}\left(p^{f}\right)$ is a root of unity. Hence, the dimension of (2.8.1) is equal to the sum of the dimensions of the Frob $p^{f}$-eigenspaces of $\rho_{a_{\mathbf{0}}}\left(\pi_{\mathfrak{p}}\right)$ with eigenvalues $\left(p^{f}\right)^{(n-1) / 2 \sum_{i} a_{i}} \zeta$ for some root of unity $\zeta$. In many examples, this space is known to be nonzero.

For instance, when $f=2, a_{1}=r$ and $a_{2}=n-r$ for some $1 \leq r \leq n-1$, we have $d\left(a_{0}\right)=2 r(n-r)$ and

$$
\rho_{a_{\mathbf{0}}}\left(\pi_{\mathfrak{p}}\right)=\wedge^{r} \rho_{\pi_{\mathfrak{p}}} \otimes \wedge^{n-r} \rho_{\pi_{\mathfrak{p}}}
$$

Let $V_{\pi_{\mathfrak{p}}, a_{\mathbf{\bullet}}}$ denote the space of representation $\rho_{a_{\mathbf{0}}}\left(\pi_{\mathfrak{p}}\right)$. If $\rho_{\pi_{\mathfrak{p}}}\left(\right.$ Frob $\left._{p^{f}}\right)$ has distinct eigenvalues $\alpha_{1}, \ldots, \alpha_{n}$, then the eigenvalues of $\mathrm{Frob}_{p^{f}}$ on $V_{\pi_{\mathrm{p}}, a_{0}}$ are given by $\alpha_{i_{1}} \cdots \alpha_{i_{r}} \cdot \alpha_{j_{1}} \cdots \alpha_{j_{n-r}}$, for distinct subscripts $i_{1}, \ldots, i_{r}$ and distinct subscripts $j_{1}, \ldots, j_{n-r}$. This product is exactly $\left(p^{f}\right)^{n(n-1) / 2} a_{\mathfrak{p}}^{(n)}$ (note that $a_{\mathfrak{p}}^{(n)}$ is a root of unity) if the set $\left\{i_{1}, \ldots, i_{r}\right\}$ and the set $\left\{j_{1}, \ldots, j_{n-r}\right\}$ are the complement of each other as subsets of $\{1, \ldots, n\}$. On the other hand, if the subsets $\left\{i_{1}, \ldots, i_{r}\right\}$ and $\left\{j_{1}, \ldots, j_{n-r}\right\}$ are not the complement of each other and if the $\alpha_{i}$ are "sufficiently generic" ${ }^{8}$, the eigenvalue $\alpha_{i_{1}} \cdots \alpha_{i_{r}} \cdot \alpha_{j_{1}} \cdots \alpha_{j_{n-r}}$ is not a root of unity. In other words, the dimension of (2.8.1) is "generically" equal to $\binom{n}{r}$. As predicted by the Tate conjecture, these cohomology classes should come from algebraic cycles. Our main conjecture addresses exactly this, and it predicts that those desired "generic"

[^6]algebraic cycles can be given by the irreducible components of the basic locus, and are birationally equivalent to certain fiber bundles over the special fiber of some other Shimura varieties associated to inner forms of $G_{a_{0}}$. To make this precise, we need the following lemma.
Lemma 2.9. Let $b_{\bullet}=\left(b_{i}\right)_{1 \leq i \leq f}$ be a tuple with $b_{i} \in\{0, \ldots, n\}$ such that $\sum_{i=1}^{f} b_{i} \equiv$ $\sum_{i=1}^{f} a_{i}(\bmod 2)$ if $n$ is even. Then there exists $\beta_{b_{0}} \in\left(D^{\times}\right)^{*=-1}$ such that

- the alternating $D$-Hermitian space $\left(V_{b_{\bullet}},\langle-,-\rangle_{b_{0}}\right)$ defined using $\beta_{b_{0}}$ in place of $\beta_{a_{0}}$ is isomorphic to $\left(V_{a_{\bullet}},\langle-,-\rangle_{a_{0}}\right)$ when tensored with $\mathbb{A}^{\infty}$, and
- if $G_{b_{0}}$ denotes the corresponding algebraic group over $\mathbb{Q}$ defined in the similar way with $\beta_{a_{0}}$ replaced by $\beta_{b_{\bullet}}$, then

$$
G_{b_{\mathbf{\bullet}}}^{1}(\mathbb{R}) \simeq \prod_{i=1}^{f} U\left(b_{i}, n-b_{i}\right)
$$

Proof. We reduce the problem to the existence of a certain cohomology class. Note that $G_{a_{0}}^{1}=\operatorname{Aut}\left(V_{a_{0}},\langle-,-\rangle_{a_{0}}\right)$ is the Weil restriction to $\mathbb{Q}$ of a unitary group $U_{a_{0}}$ over $F$. The cohomology set $H^{1}\left(\mathbb{Q}, G_{a_{0}}^{1}\right) \cong H^{1}\left(F, U_{a_{\bullet}}\right)$ is in bijection with the isomorphism classes of one-dimensional skew-Hermitian $D$-modules $V$. As $U_{a_{0}} \times_{F} E \simeq \mathrm{GL}_{n, E}$, Hilbert's Theorem 90 for $\mathrm{GL}_{n}$ implies that the inflation map induces an isomorphism

$$
H^{1}\left(E / F, U_{a_{\bullet}}\right) \xrightarrow{\sim} H^{1}\left(F, U_{a_{0}}\right)
$$

Denote by $g \mapsto g^{\sharp \beta_{a_{\bullet}}}=\beta_{a_{0}} g^{*} \beta_{a_{\bullet}}^{-1}$ the involution on $D$ induced by the alternating pairing $\langle-,-\rangle_{a_{0}}$. Then a 1-cocycle of $\operatorname{Gal}(E / F)$ with values in $U_{a_{0}}$ is given by an element $\alpha \in D^{\times}$such that $\alpha=\alpha^{\sharp \beta_{a_{\bullet}}}$, and $\alpha_{1}, \alpha_{2} \in D^{\times}$define the same cohomology class in $H^{1}\left(F, U_{a_{0}}\right)$ if and only if there exists $g \in D^{\times}$such that $g \alpha_{1} g^{\sharp \beta_{a_{\bullet}}}=\alpha_{2}$. Explicitly, given such an $\alpha$, the corresponding skew-Hermitian $D$-module is given by $V=D$ equipped with the alternating pairing

$$
\langle-,-\rangle_{\alpha}: V \times V \rightarrow \mathbb{Q}, \quad(x, y) \mapsto \operatorname{Tr}_{D / \mathbb{Q}}\left(x \alpha \beta_{a_{0}} y^{*}\right)
$$

For a place $v$ of $F$, we denote by

$$
\operatorname{loc}_{v}: H^{1}\left(F, U_{a_{0}}\right) \rightarrow H^{1}\left(F_{v}, U_{a_{0}}\right)
$$

the canonical localization map. By [Helm 2012, Proposition 8.1], if $\sum_{i=1}^{f} b_{i} \equiv$ $\sum_{i=1}^{f} a_{i} \bmod 2$, there exists a cohomology class $[\alpha] \in H^{1}\left(F, U_{a_{0}}\right)$ such that

- $\operatorname{loc}_{v}([\alpha])$ is trivial for every finite place $v$ of $F$, and
- if $v=\tau_{i}$ with $1 \leq i \leq n$ is an archimedean place, one has an isomorphism of unitary groups over $\mathbb{R}: \operatorname{Aut}\left(V \otimes_{F, \tau_{i}} \mathbb{R},\langle-,-\rangle_{\alpha}\right) \simeq U\left(b_{i}, n-b_{i}\right)$.
Then the element $\beta_{b_{\bullet}}=\alpha \beta_{a_{0}}$ meets the requirements of Lemma 2.9.

In the sequel, we always fix a choice of $\beta_{b_{\bullet}}$, and as well as an isomorphism $\gamma_{a_{0}, b_{\mathbf{0}}}: V_{a_{\bullet}} \otimes_{\mathbb{Q}} \mathbb{A}^{\infty} \xrightarrow{\sim} V_{b_{\mathbf{0}}} \otimes_{\mathbb{Q}} \mathbb{A}^{\infty}$, which induces an isomorphism $G_{a_{\mathbf{0}}}\left(\mathbb{A}^{\infty}\right) \simeq$ $G_{b_{0}}\left(\mathbb{A}^{\infty}\right)$. Recall that we have chosen a lattice $\Lambda_{a_{0}} \subseteq V_{a_{0}}$ to define the moduli problem for $\mathcal{S} h_{a_{\bullet}}$. We put $\Lambda_{b_{\bullet}}:=V_{b_{\bullet}} \cap \gamma_{a_{\bullet}, b_{\bullet}}\left(\Lambda_{a_{\bullet}} \otimes_{\mathbb{Z}} \widehat{\mathbb{Z}}\right)$. Then applying the construction of Section 2.3 to the lattice $\Lambda_{b_{\mathbf{\bullet}}} \subseteq V_{b_{\boldsymbol{\bullet}}}$ and the open compact subgroup $K^{p} \subseteq G_{a_{\mathbf{0}}}\left(\mathbb{A}^{\infty, p}\right) \simeq G_{b_{\bullet}}\left(\mathbb{A}^{\infty, p}\right)$, we get a Shimura variety $S h_{b_{\bullet}}$ over $\mathbb{Z}_{p^{f}}$ of level $K^{p}$ as well as its special fiber $\mathrm{Sh}_{b_{0}}$. Moreover, an algebraic representation $\xi$ of $G_{a_{0}}$ over $\overline{\mathbb{Q}}_{\ell}$ corresponds, via the fixed isomorphism $G_{a_{0}}\left(\mathbb{A}^{\infty}\right) \simeq G_{b_{0}}\left(\mathbb{A}^{\infty}\right)$, to an algebraic representation of $G_{b_{\bullet}}$ over $\overline{\mathbb{Q}}_{\ell}$. We use the same notation $\mathcal{L}_{\xi}$ to denote the étale sheaf on $\mathcal{S} h_{a_{0}}$ and $\mathcal{S} h_{b_{\bullet}}$ defined by $\xi$.
2.10. Gysin/trace maps. Before stating the main conjecture of this paper, we recall the general definition of Gysin maps. Let $f: Y \rightarrow X$ be a proper morphism of smooth varieties over an algebraically closed field $k$. Let $d_{X}$ and $d_{Y}$ be the dimensions of $X$ and $Y$ respectively. Recall that the derived direct image $R f_{*}$ on the derived category of constructible $\ell$-adic étale sheaves has a left adjoint $f^{!}$. Since both $X$ and $Y$ are smooth, the $\ell$-adic dualizing complex of $X$ (resp. $Y$ ) is $\overline{\mathbb{Q}}_{\ell}\left(d_{X}\right)\left[2 d_{X}\right]$ (resp. $\overline{\mathbb{Q}}_{\ell}\left(d_{Y}\right)\left[2 d_{Y}\right]$ ). Therefore, one has

$$
f^{!}\left(\overline{\mathbb{Q}}_{\ell}\left(d_{X}\right)\left[2 d_{X}\right]\right)=\overline{\mathbb{Q}}_{\ell}\left(d_{Y}\right)\left[2 d_{Y}\right]
$$

The adjunction map $R f_{*} f^{!} \overline{\mathbb{Q}}_{\ell} \rightarrow \overline{\mathbb{Q}}_{\ell}$ induces a canonical morphism

$$
\operatorname{Tr}_{f}: R f_{*} \overline{\mathbb{Q}}_{\ell} \rightarrow \overline{\mathbb{Q}}_{\ell}\left(d_{X}-d_{Y}\right)\left[2\left(d_{X}-d_{Y}\right)\right]
$$

More generally, if $\mathcal{L}$ is a lisse $\overline{\mathbb{Q}}_{\ell}$-sheaf on $X$, it induces a Gysin/trace map

$$
R f_{*}\left(f^{*} \mathcal{L}\right) \cong \mathcal{L} \otimes R f_{*}\left(\overline{\mathbb{Q}}_{\ell}\right) \xrightarrow{1 \otimes \operatorname{Tr}_{f}} \mathcal{L}\left(d_{X}-d_{Y}\right)\left[2\left(d_{X}-d_{Y}\right)\right],
$$

where the first isomorphism is the projection formula [SGA $4_{2}$ 1972, XVII 5.2.9]. When $f$ is flat with equidimensional fibers of dimension $d_{Y}-d_{X}$, this is the trace map as defined in [SGA $4_{2}$ 1972, XVIII 2.9]. When $f$ is a closed immersion of codimension $r=d_{X}-d_{Y}$, it is the usual Gysin map. For any integer $q$, the Gysin/trace map induces a morphism on cohomology groups:

$$
\begin{equation*}
f_{!}: H_{\mathrm{et}}^{q}\left(Y, f^{*} \mathcal{L}\right) \rightarrow H_{\mathrm{et}}^{q+2\left(d_{X}-d_{Y}\right)}\left(X, \mathcal{L}\left(d_{X}-d_{Y}\right)\right) \tag{2.10.1}
\end{equation*}
$$

2.11. Representation theory of $\mathbf{G L}_{\boldsymbol{n}}$. As suggested by the description of Galois representations appearing in the middle cohomology group of Shimura varieties in Theorem 2.6, as well as by the Tate conjecture, we need to understand the representation theory of $\mathrm{GL}_{n}$ embedded diagonally into the Langlands dual group

$$
\left({ }^{L} G_{a_{0}}\right)^{\circ} \simeq \mathbb{C}^{\times} \times \mathrm{GL}_{n}(\mathbb{C})^{\mathbb{Z} / f \mathbb{Z}}
$$

The Hodge cocharacter $\mu$ of $G_{a_{0}}$ gives rise to the representation $r_{a_{\bullet}}=\bigotimes_{i=1}^{f}\left(\wedge^{a_{i}} \operatorname{Std}\right)$ of the diagonal $\mathrm{GL}_{n}$. If $\lambda$ is a dominant weight of $\mathrm{GL}_{n}$ (with respect to the usual diagonal torus and upper triangular Borel subgroup) appearing in $r_{a_{\bullet}}$, we can write this weight $\lambda$ as the sum of $f$ dominant minuscule weights $\omega_{b_{1}}+\cdots+\omega_{b_{f}}$, where $\omega_{i}$ for $0 \leq i \leq n$ is the weight of $\mathrm{GL}_{n}$ that takes $\operatorname{Diag}\left(t_{1}, \ldots, t_{n}\right)$ to $t_{1} \cdots t_{i}$. The set $\left\{b_{1}, \ldots, b_{f}\right\}$ (counted with multiplicity) is unique, which we denote by $B_{\lambda}$. Explicitly, if $\lambda$ takes $\operatorname{Diag}\left(t_{1}, \ldots, t_{n}\right)$ to $t_{1}^{\beta_{1}} \cdots t_{n}^{\beta_{n}}$ (necessarily $\beta_{1} \leq f$ ), then

$$
B_{\lambda}=\{\underbrace{n, \ldots, n}_{\beta_{n}}, \underbrace{n-1, \ldots, n-1}_{\beta_{n-1}-\beta_{n}}, \ldots, \underbrace{1, \ldots, 1}_{\beta_{1}-\beta_{0}}, \underbrace{0, \ldots, 0}_{f-\beta_{1}}\} .
$$

Moreover, we always have $\sum a_{i}=\sum b_{i}$. In particular, this implies by Lemma 2.9 that the Shimura variety $\mathrm{Sh}_{b_{\bullet}}$ makes sense, and the étale sheaf $\mathcal{L}_{\xi}$ is well defined on $\mathrm{Sh}_{b}$.

We write $m_{\lambda}\left(a_{0}\right)$ for the multiplicity of the weight $\lambda$ in $r_{a_{0}}$.
Conjecture 2.12. Let $\mathrm{Sh}_{a_{\bullet}}$ and $\mathcal{L}_{\xi}$ be as in Section 2.4. Let $\lambda$ be a dominant weight that appears in the representation $r_{a_{0}}$ as in Section 2.11. Define $B_{\lambda}$ and $m_{\lambda}\left(a_{0}\right)$ as in Section 2.11.

Then there exist varieties $Y_{1}, \ldots, Y_{m_{\lambda}\left(a_{\mathbf{0}}\right)}$ of dimension $\frac{1}{2}\left(d\left(a_{\mathbf{0}}\right)+d\left(b_{\mathbf{0}}\right)\right)$ over $\mathbb{F}_{p^{f}}$, equipped with natural action of prime-to-p Hecke correspondences, such that each $Y_{j}$ fits into a diagram

satisfying the following properties.
(1) For each $j, b_{\cdot}^{(j)}=\left(b_{1}^{(j)}, \ldots, b_{f}^{(j)}\right)$ is a reordering of the elements of the set $B_{\lambda}$, and both $\mathrm{pr}_{a_{0}}^{(j)}$ and $\mathrm{pr}_{b_{\mathbf{0}}^{(j)}}$ are equivariant for the prime-to-p Hecke correspondences.
(2) The morphism $\operatorname{pr}_{a_{0}}^{(j)}$ is a proper morphism and is birational onto the image. The morphism $\mathrm{pr}_{b_{(j)}^{(j)}}$ is proper and generically smooth of relative dimension $\frac{1}{2}\left(d\left(a_{\bullet}\right)-d\left(b_{\bullet}\right)\right)\left(\right.$ note that $d\left(b_{\bullet}\right) \equiv d\left(a_{\bullet}\right)(\bmod 2)$ since $\left.\sum_{i} a_{i}=\sum_{i} b_{i}\right)$.
(3) There exists a p-isogeny of abelian schemes over $Y_{j}$

$$
\phi_{b_{\bullet}^{(j)}, a_{\bullet}}: \operatorname{pr}_{b_{\bullet}(j)}^{*}\left(\mathcal{A}_{b_{\bullet}^{(j)}}\right) \rightarrow \operatorname{pr}_{a_{\bullet}}^{(j), *}\left(\mathcal{A}_{a_{\bullet}}\right)
$$

where $\mathcal{A}_{a_{0}}$ and $\mathcal{A}_{b_{0}^{(j)}}$ denote respectively the universal abelian scheme on $\mathrm{Sh}_{a_{0}}$ and $\mathrm{Sh}_{b_{\bullet}^{(j)}}$. Let

$$
\phi_{b_{\bullet}^{(j)}, a_{\bullet}, *}: \operatorname{pr}_{b_{\bullet}^{(j)}}^{*} \mathcal{L}_{\xi} \xrightarrow{\sim} \operatorname{pr}_{a_{\bullet}}^{(j), *} \mathcal{L}_{\xi}
$$

be the isomorphism of the $\ell$-adic sheaves induced by $\phi_{b_{\mathbf{*}}^{(j)}, a_{0}}$ via the construction in Section 2.4. ${ }^{9}$
(4) Let $\pi$ be an irreducible admissible representation of $G_{a_{0}}\left(\mathbb{A}^{\infty}\right) \simeq G_{\left.b_{0}\right)}\left(\mathbb{A}^{\infty}\right)$ satisfying Hypothesis 2.5 for both $a_{\bullet}$ and $b_{\bullet}$, and assume that $m_{a_{\bullet}}(\pi)=m_{b_{\bullet}^{(j)}}(\pi)$ for all $j^{10}$. Suppose that the $n$ eigenvalues $\alpha_{1}, \ldots, \alpha_{n}$ of $\rho_{\pi_{\mathfrak{p}}}\left(\operatorname{Frob}_{p^{f}}\right)$ are "sufficiently generic" in the sense that the generalized eigenspace decomposition of $\rho_{a_{\bullet}}\left(\operatorname{Frob}_{p^{N}}\right)$ for any large $N$ is the same as the weight space decomposition of the algebraic representation $r_{a_{0}}$. Then the natural homomorphism of $\pi$ isotypic components ${ }^{11}$ of the cohomology groups

$$
\begin{aligned}
& \bigoplus_{j=1}^{m_{\lambda}\left(a_{\bullet}\right)} H_{\mathrm{et}}^{d\left(b_{\bullet}\right)}\left(\overline{\operatorname{Sh}}_{b_{\bullet}(j)}, \mathcal{L}_{\xi}\left(\frac{1}{2} d\left(b_{\bullet}\right)\right)\right)_{\pi}^{\operatorname{Frob}_{p} f=\lambda} \\
& \xrightarrow{\oplus \mathrm{pr}_{b_{\mathbf{\bullet}}^{*}(j)}} \bigoplus_{j=1}^{m_{\lambda}\left(a_{\bullet}\right)} H_{\mathrm{et}}^{d\left(b_{\mathbf{\bullet}}\right)}\left(\bar{Y}_{j}, \mathrm{pr}_{b_{\mathbf{\bullet}}}^{*} \mathcal{L}_{\xi}\left(\frac{1}{2} d\left(b_{\mathbf{\bullet}}\right)\right)\right)_{\pi}^{\mathrm{Frob}_{p} f=\lambda} \\
& \xrightarrow{\oplus \phi_{b_{\bullet}(j)},_{\mathbf{\bullet}}, *} \bigoplus_{j=1}^{m_{\lambda}\left(a_{\bullet}\right)} H_{\mathrm{et}}^{d\left(b_{\bullet}\right)}\left(\bar{Y}_{j}, \operatorname{pr}_{a_{\bullet}}^{*} \mathcal{L}_{\xi}\left(\frac{1}{2} d\left(b_{\mathbf{\bullet}}\right)\right)\right)_{\pi}^{\operatorname{Frob}_{p} f=\lambda} \\
& \xrightarrow{\sum \mathrm{pr}_{a_{\bullet},}^{(j)}} H_{\mathrm{et}}^{d\left(a_{\bullet}\right)}\left(\overline{\operatorname{Sh}}_{a_{\mathbf{\bullet}}}, \mathcal{L}_{\xi}\left(\frac{1}{2} d\left(a_{\mathbf{\bullet}}\right)\right)\right)_{\pi}^{\mathrm{Frob}_{p} f=\lambda}
\end{aligned}
$$

is an isomorphism, where $\mathrm{pr}_{a_{0},!}^{(j)}$ is the Gysin map (2.10.1) and the superscript $\mathrm{Frob}_{p^{f}}=\lambda$ means taking the (direct sum of) generalized Frob $_{p^{f}}$-eigenspace with eigenvalues in the Weyl group orbit

$$
\lambda \circ \rho_{\pi_{\mathfrak{p}}}\left(\operatorname{Frob}_{p^{f}}\right) \cdot \chi_{\pi_{p, 0}}^{-1}\left(p^{f}\right)(\sqrt{p})^{-f(n-1) \sum_{i} b_{i}}
$$

Here, since the semisimple conjugacy classes of $\mathrm{GL}_{n}\left(\overline{\mathbb{Q}}_{\ell}\right)$ is in natural bijection with the orbits of $T\left(\overline{\mathbb{Q}}_{\ell}\right)$ under the Weyl group of $\mathrm{GL}_{n}$, it makes sense to evaluate a dominant weight of $T$ on $\rho_{\pi_{\mathfrak{p}}}\left(\mathrm{Frob}_{p^{f}}\right)$ to get an orbit under the action of the Weyl group of $\mathrm{GL}_{n}$; hence the notation $\lambda \circ \rho_{\pi_{\mathfrak{p}}}\left(\operatorname{Frob}_{p^{f}}\right)$.

In particular, when $\xi$ is the trivial representation and the weight $\lambda$ is a power of the determinant (so automatically, $\sum_{i} a_{i}$ is divisible by $n$, and $d\left(a_{0}\right)$ is even), the cycles given by the images of $Y_{1}, \ldots, Y_{m_{\lambda}\left(a_{0}\right)}$ parameterized by the discrete Shimura varieties $\mathrm{Sh}_{b_{\bullet}^{(j)}}$, generate the Tate classes of $H_{\mathrm{et}}^{d\left(a_{\bullet}\right)}\left(\overline{\mathrm{S}}_{a_{\bullet}}, \overline{\mathbb{Q}}_{\ell}\left(\frac{1}{2} d\left(a_{0}\right)\right)\right)_{\pi}$ when $\rho_{\pi_{\mathfrak{p}}}\left(\mathrm{Frob}_{p^{f}}\right)$ is "sufficiently generic".

[^7]Remark 2.13. (1) A key feature of this conjecture is that the codimension of the cycle map $\mathrm{pr}_{a_{\bullet}}: Y_{j} \rightarrow \mathrm{Sh}_{a_{\bullet}}$ is the same as the fiber dimension of $\mathrm{pr}_{b_{\mathbf{0}}^{(j)}}: Y_{j} \rightarrow \mathrm{Sh}_{b_{\mathbf{0}}^{(j)}}$.
(2) It seems that the fiber of $\mathrm{pr}_{b_{\bullet}^{(j)}}: Y_{j} \rightarrow \mathrm{Sh}_{b_{\bullet}^{(j)}}$ over a generic point $\eta \in \mathrm{Sh}_{b_{\bullet}(j)}$ is likely to be isomorphic to a certain "iterated Deligne-Lusztig variety," that is, a tower of maps $Y_{j, \eta}=Z_{\alpha} \rightarrow \cdots \rightarrow Z_{0}=\eta$ such that each $Z_{i} \rightarrow Z_{i-1}$ is a fiber bundle with certain Deligne-Lusztig varieties as fibers.
(3) Xinwen Zhu pointed out to us that since the universal abelian varieties $\mathcal{A}_{a}$ and $\mathcal{A}_{b_{\bullet}}$ are isogenous over each $Y_{j}$, the union of the images of $Y_{1}, \ldots, Y_{m_{\lambda}\left(a_{\bullet}\right)}$ on $\mathrm{Sh}_{a_{0}}$ is contained in the closure of the Newton strata, where the slope is the same as the $\mu$-ordinary slope of the universal abelian varieties on $\mathrm{Sh}_{b_{\bullet}^{(j)}}$ (for different $j$, they have the same $\mu$-ordinary slopes). In fact, one should expect the union of images to be the same as the closure of this Newton stratum.

When $\lambda$ is central (i.e., a power of the determinant), Conjecture 2.12 says: irreducible components of the basic locus of the special fiber of a Shimura variety, generically, contribute to all Tate cycles in the cohomology. Implicitly, this means that the dimension of the basic locus is half of the dimension of the Shimura variety if and only if the Galois representations of the Shimura variety has generically nontrivial Tate classes. Here two appearances of "generic" both mean that we only consider those $\pi$-isotypic components where the Satake parameter for $\pi_{p}$ is sufficiently generic as in Conjecture 2.12(4). For example, the supersingular locus of Hilbert modular surface at a split prime or the supersingular locus of a Siegel modular variety (over $\mathbb{Q}$ ) is not half the dimension. This is related to the fact that the $\pi$-isotypic component of the cohomology of the Shimura varieties are not expected to have Tate classes, at least when the Satake parameter of $\pi_{p}$ is sufficiently general. ${ }^{12}$
(4) These varieties $Y_{j}$ may be viewed as Hecke correspondences at $p$ between the special fibers of two different Shimura varieties $\mathrm{Sh}_{a_{0}}$ and $\mathrm{Sh}_{b_{0}^{(j)}}$. These correspondences certainly cannot be lifted to characteristic zero. We hope that the conjecture will bring interests into the study of such Hecke correspondences.
Remark 2.14. (1) The assumption on the decomposition of the place $p$ in $E / \mathbb{Q}$ and working with unitary Shimura varieties is to simplify our presentation and to get to a situation where most terms can be defined. We certainly expect the validity of analogous conjectures for the special fibers of Shimura varieties of PEL-type or more generally of abelian type (using the integral model of M. Kisin [2010]). This would be a more precise version of the Tate conjecture in the context of special fibers of Shimura varieties: if $\mathrm{Sh}_{G}$ and $\mathrm{Sh}_{G^{\prime}}$ are the

[^8]special fibers of two unitary Shimura varieties associated to the groups $G$ and $G^{\prime}$ such that $G\left(\mathbb{A}_{f}\right) \simeq G^{\prime}\left(\mathbb{A}_{f}\right)$, then, generically, the cycles on the product $\mathrm{Sh}_{G} \times \mathrm{Sh}_{G^{\prime}}$ predicted by the Tate conjectures are likely to be constructed by understanding the "isogenies" between the corresponding universal abelian varieties, and are closely related to the Newton stratifications of $\mathrm{Sh}_{G}$ and $\mathrm{Sh}_{G^{\prime}}$. In the case of Shimura varieties of abelian type, we expect some technical difficulties in reinterpreting the meaning of isogenies of abelian varieties in terms of certain " $G$-crystals".

For example, consider a real quadratic field $F / \mathbb{Q}$ in which a prime $p$ is inert. Let $\mathrm{Sh}_{G}$ denote the special fiber of the Hilbert-Siegel modular variety for $G:=$ $\operatorname{Res}_{F / \mathbb{Q}} \mathrm{GSp}_{2 g}$, with hyperspecial level structure at $p$. Then by Langlands's prediction of the cohomology of $\mathrm{Sh}_{G}$, we should look at the representation $r_{\text {spin }}^{\otimes 2}$ of the "essential part" $\operatorname{Spin}_{2 g+1}$ of the Langlands dual group, where $r_{\text {spin }}$ is the $2^{g}$-dimensional spin representation. ${ }^{13}$ The central weight space of $r_{\text {spin }}^{\otimes 2}$ has dimension $2^{g}$. So we expect that the supersingular locus of $\mathrm{Sh}_{G}$ is the union of $2^{g}$ collection of varieties parameterized by the discrete Shimura variety $\mathrm{Sh}_{G^{\prime}}$ where $G^{\prime}$ is the inner form of $G$ which is split at all finite places and is compact modulo center at both archimedean places. Unfortunately, the moduli problem that describes $G^{\prime}$ uses a different division algebra from that describing $G$. We do not know how to interpret the meaning of isogenies of universal abelian varieties in this case, and the method of our paper does not apply directly to this case.
(2) Xinwen Zhu pointed out to us that even if $p$ is ramified, we should expect Conjecture 2.12 continue to hold for (the special fiber of) the "splitting models" of Pappas and Rapoport [2005]. Some evidences of this have already appeared in the case of Hilbert modular varieties; see [Rapoport et al. 2014; Reduzzi and Xiao 2017].
(3) In our setup, we took advantage of many coincidences that ensures that for example the Shimura variety is compact and there is no endoscopy. It would be certainly an interesting future question to study the case involving Eisenstein series, as well as the case when the representations come from endoscopy transfers.
(4) As explained in Remark 2.13(3), the images of $Y_{j}$ are expected to form the closure of a certain Newton polygon where the slopes are related to $\lambda$. Conjecture 2.12(1)-(3) may have a degenerate situation: when $\sum_{i} a_{i}$ is not divisible by $n$, the representation $V_{a}$ does not contain a weight corresponding to a power of the determinant (which corresponds to the basic locus). So our

[^9]conjecture does not describe the basic locus of $\mathrm{Sh}_{a_{0}}$, and it is indeed not of half dimension of $\mathrm{Sh}_{a_{\bullet}}$.

Yet, this basic locus may still have a good description as the union of some fiber bundles over the special fibers of some other Shimura varieties for reductive groups which are not quasisplit at $p$. For example, the supersingular locus of modular curve is related to the Shimura variety associated to the definite quaternion algebra which is ramified at $p$, by a theorem of Serre and Deuring [Serre 1996]. More such examples are given in [Tian and Xiao 2016] and [Vollaard and Wedhorn 2011].
2.15. Known cases of Conjecture 2.12. Conjecture 2.12 is largely inspired by the work of Tian and Xiao [2014; 2016], where we proved the analogous conjecture for the special fibers of the Hilbert modular varieties assuming that $p$ is inert in the totally real field.

Another strong evidence of Conjecture 2.12 is the work of Vollaard and Wedhorn [2011], where they considered certain stratification of the supersingular locus of the Shimura variety for $G U(1, n-1)$ with $s \in \mathbb{N}$ at an inert prime $p$. What concerns us is the case when $n-1$ is even. In this case, it is hidden in the writing of their Section 6 that one gets a correspondence (in the notation of loc. cit.)


Note that $I\left(\mathbb{A}_{f}\right) \simeq \mathbb{G}\left(\mathbb{A}_{f}\right)$. Here $N_{n}$ is a certain Deligne-Lusztig variety. In [Vollaard and Wedhorn 2011], the parameterizing space, namely the first term in (2.15.1), is interpreted very differently, in terms of Bruhat-Tits building. The method of this paper should be applicable to their situation to verify the analogous Conjecture 2.12. In fact, in their case, there will be only one collection of cycles as given by (2.15.1), but the computation of the intersection matrix (only essentially one entry in this case) of them requires some nontrivial Schubert calculus similar to Section 5.

When $n-1$ is odd, the result of [Vollaard and Wedhorn 2011] is related to the degenerate version of the Conjecture 2.12 in the sense of Remark 2.14(4).

The aim of the rest of the paper is to provide evidence for Conjecture 2.12 for some large rank groups. In particular, we will construct cycles in the case of the unitary group $G(U(r, s) \times U(s, r))$ with $s, r \in \mathbb{N}$ (Section 7). While we expect these cycles to verify Conjecture 2.12 , we do not know how to compute the "intersection matrix" in general. Nonetheless, when $r=1$, we are able to
make the computation and prove Conjecture 2.12 (with trivial coefficients for the sake of a simple presentation) in this case; see Section 4-6. We point out that our method should be applicable to many other examples, and even in general reduce Conjecture 2.12 to a question of a combinatorial nature. This combinatorics problem is the heart of the question. In the Hilbert case [Tian and Xiao 2014], we model the combinatorics question by the so-called periodic semimeander (for $\mathrm{GL}_{2}$ ). The generalization of the usual (as opposed to periodic) semimeander to other groups has been introduced; see [Fontaine et al. 2013] for the corresponding references. The straightforward generalization to the periodic case does seem to agree with some of our computations with small groups. Nonetheless, the corresponding Gram determinant formula seems to be extremely difficult. Even in the nonperiodic case, we only know it for a special case; see [Di Francesco 1997].

We also mention that in a very recent work [Xiao and Zhu 2017] of Zhu and the last author, we relate Conjecture 2.12 with the geometric Satake theory of Zhu [2017] in mixed characteristic, and we proved many new cases of Conjecture 2.12.

## 3. Preliminaries on Dieudonné modules and deformation theory

We first introduce the basic tools that we will use in this paper.
3.1. Notation. Recall that we have an isomorphism
$\mathcal{O}_{D} \otimes_{\mathbb{Z}} \mathbb{Z}_{p^{f}} \cong \bigoplus_{i=1}^{f}\left(\mathcal{O}_{D} \otimes_{\mathcal{O}_{E}, q_{i}} \mathbb{Z}_{p^{f}} \oplus \mathcal{O}_{D} \otimes_{\mathcal{O}_{E}, \bar{q}_{i}} \mathbb{Z}_{p^{f}}\right) \simeq \bigoplus_{i=1}^{f}\left(\mathrm{M}_{n}\left(\mathbb{Z}_{p^{f}}\right) \oplus \mathrm{M}_{n}\left(\mathbb{Z}_{p^{f}}\right)\right)$.
Let $S$ be a locally noetherian $\mathbb{Z}_{p^{f}}$-scheme. An $\mathcal{O}_{D} \otimes_{\mathbb{Z}} \mathcal{O}_{S}$-module $M$ admits a canonical decomposition

$$
M=\bigoplus_{i=1}^{f}\left(M_{q_{i}} \oplus M_{\bar{q}_{i}}\right),
$$

where $M_{q_{i}}\left(\right.$ resp. $M_{\bar{q}_{i}}$ ) is the direct summand of $M$ on which $\mathcal{O}_{E}$ acts via $q_{i}$ (resp. via $\left.\bar{q}_{i}\right)$. Then each $M_{q_{i}}$ has a natural action by $\mathrm{M}_{n}\left(\mathcal{O}_{S}\right)$. Let $\mathfrak{e}$ denote the element of $\mathrm{M}_{n}\left(\mathcal{O}_{S}\right)$ whose (1,1)-entry is 1 and whose other entries are 0 . We put $M_{i}^{\circ}:=\mathfrak{e} M_{q_{i}}$, and call it the reduced part of $M_{q_{i}}$.

Let $A$ be an $f n^{2}$-dimensional abelian variety over an $\mathbb{F}_{p^{f}}$-scheme $S$, equipped with an $\mathcal{O}_{D}$-action. The de Rham homology $H_{1}^{\mathrm{dR}}(A / S)$ has a Hodge filtration

$$
0 \rightarrow \omega_{A^{\vee} / S} \rightarrow H_{1}^{\mathrm{dR}}(A / S) \rightarrow \operatorname{Lie}_{A / S} \rightarrow 0
$$

compatible with the natural action of $\mathcal{O}_{D} \otimes_{\mathbb{Z}} \mathcal{O}_{S}$ on $H_{1}^{\mathrm{dR}}(A / S)$. We call $H_{1}^{\mathrm{dR}}(A / S)_{i}^{\circ}$ (resp. $\omega_{A^{\vee} / S, i}^{\circ}, \mathrm{Lie}_{A / S, i}^{\circ}$ ) the reduced de Rham homology of $A / S$ (resp. the reduced invariant 1 -forms of $A^{\vee} / S$, the reduced Lie algebra of $A / S$ ) at $q_{i}$. In particular, the
former is a locally free $\mathcal{O}_{S}$-module of rank $n$ and the latter is a subbundle ${ }^{14}$ of the former; when $A \rightarrow S$ satisfies the moduli problem in Section 2.3, $\omega_{A^{\vee} / S, i}^{\circ}$ is locally free of rank $a_{i}$.

The Frobenius morphism $A \rightarrow A^{(p)}$ induces a natural homomorphism

$$
V: H_{1}^{\mathrm{dR}}(A / S)_{i}^{\circ} \rightarrow H_{1}^{\mathrm{dR}}(A / S)_{i-1}^{\circ,(p)}
$$

where the index $i$ is considered as an element of $\mathbb{Z} / f \mathbb{Z}$, and the superscript " $(p)$ " means the pullback via the absolute Frobenius of $S$. The image of $V$ is exactly $\omega_{A^{\vee} / S, i-1}^{0,(p)}$. Similarly, the Verschiebung morphism $A^{(p)} \rightarrow A$ induces a natural homomorphism ${ }^{15}$

$$
F: H_{1}^{\mathrm{dR}}(A / S)_{i-1}^{\circ,(p)} \rightarrow H_{1}^{\mathrm{dR}}(A / S)_{i}^{\circ}
$$

We have $\operatorname{Ker}(F)=\operatorname{Im}(V)$ and $\operatorname{Ker}(V)=\operatorname{Im}(F)$.
When $S=\operatorname{Spec}(k)$ with $k$ a perfect field containing $\mathbb{F}_{p^{f}}$, let $W(k)$ denote the ring of Witt vectors in $k$. Let $\tilde{\mathcal{D}}(A)$ denote the (covariant) Dieudonné module associated to the $p$-divisible group of $A$. This is a free $W(k)$-module of rank $2 f n^{2}$ equipped with a Frob-linear action of $F$ and a $\mathrm{Frob}^{-1}$-linear action of $V$ such that $F V=V F=p$. The $\mathcal{O}_{D}$-action on $A$ induces a natural action of $\mathcal{O}_{D}$ on $\tilde{\mathcal{D}}(A)$ that commutes with $F$ and $V$. Moreover, there is a canonical isomorphism $\tilde{\mathcal{D}}(A) / p \tilde{\mathcal{D}}(A) \cong H_{1}^{\mathrm{dR}}(A / k)$ compatible with all structures on both sides. For each $i \in \mathbb{Z} / f \mathbb{Z}$, we have the reduced part $\tilde{\mathcal{D}}(A)_{i}^{\circ}:=\mathfrak{e} \tilde{\mathcal{D}}(A)_{q_{i}}$. The Verschiebung and the Frobenius induce natural maps

$$
V: \tilde{\mathcal{D}}(A)_{i}^{\circ} \rightarrow \tilde{\mathcal{D}}(A)_{i-1}^{\circ}, \quad F: \tilde{\mathcal{D}}(A)_{i}^{\circ} \rightarrow \tilde{\mathcal{D}}(A)_{i+1}^{\circ}
$$

Note that $\tilde{\mathcal{D}}(A)_{q_{i}}=\left(\tilde{\mathcal{D}}(A)_{i}^{\circ}\right)^{\oplus n}$, and $\bigoplus_{i \in \mathbb{Z} / f \mathbb{Z}} \tilde{\mathcal{D}}(A)_{q_{i}}$ is the covariant Dieudonné module of the $p$-divisible group $A\left[\mathfrak{p}^{\infty}\right]$.

For any $f n^{2}$-dimensional abelian variety $A^{\prime}$ over $k$ equipped with an $\mathcal{O}_{D}$-action, an $\mathcal{O}_{D}$-equivariant isogeny $A^{\prime} \rightarrow A$ induces a morphism $\tilde{\mathcal{D}}\left(A^{\prime}\right)_{i}^{\circ} \rightarrow \tilde{\mathcal{D}}(A)_{i}^{\circ}$ compatible with the actions of $F$ and $V$. Conversely, we have the following.

Proposition 3.2. Let A be an abelian variety of dimension $f n^{2}$ over prefect field $k$ which contains $\mathbb{F}_{p^{f}}$, equipped with an $\mathcal{O}_{D^{\prime}}$-action and an $\mathcal{O}_{D^{\prime}}$-compatible prime-to- $p$ polarization $\lambda$. Suppose given an integer $m \geq 1$ and a $W(k)$-submodule $\tilde{\mathcal{E}}_{i} \subseteq \tilde{\mathcal{D}}(A)_{i}^{\circ}$ for each $i \in \mathbb{Z} / f \mathbb{Z}$ such that

$$
\begin{equation*}
p^{m} \tilde{\mathcal{D}}(A)_{i}^{\circ} \subseteq \tilde{\mathcal{E}}_{i}, \quad F\left(\tilde{\mathcal{E}}_{i}\right) \subseteq \tilde{\mathcal{E}}_{i+1}, \quad \text { and } \quad V\left(\tilde{\mathcal{E}}_{i}\right) \subseteq \tilde{\mathcal{E}}_{i-1} \tag{3.2.1}
\end{equation*}
$$

[^10]Then there exists a unique abelian variety $A^{\prime}$ over $k$ (depending on $m$ ) equipped with an $\mathcal{O}_{D}$-action, a prime-to-p polarization $\lambda^{\prime}$, and an $\mathcal{O}_{D^{-}}$equivariant p-isogeny $\phi: A^{\prime} \rightarrow A$ such that the natural inclusion $\tilde{\mathcal{E}}_{i} \subseteq \tilde{\mathcal{D}}(A)_{i}^{\circ}$ is naturally identified with the map $\phi_{*, i}: \tilde{\mathcal{D}}\left(A^{\prime}\right)_{i}^{\circ} \rightarrow \tilde{\mathcal{D}}(A)_{i}^{\circ}$ induced by $\phi$ and such that $\phi^{\vee} \circ \lambda \circ \phi=p^{m} \lambda^{\prime}$. Moreover, we have
(1) If $\operatorname{dim} \omega_{A^{\vee} / k, i}^{\circ}=a_{i}$ and length ${ }_{W(k)}\left(\tilde{\mathcal{D}}(A)_{i}^{\circ} / \tilde{\mathcal{E}}_{i}\right)=\ell_{i}$ for $i \in \mathbb{Z} / f \mathbb{Z}$, then

$$
\begin{equation*}
\operatorname{dim} \omega_{A^{\prime} / k, i}^{\circ}=a_{i}+\ell_{i}-\ell_{i+1} \tag{3.2.2}
\end{equation*}
$$

(2) If $A$ is equipped with a prime-to- $p$ level structure $\eta$ in the sense of Section 2.3(1), then there exists a unique prime-to-p level structure $\eta^{\prime}$ on $A^{\prime}$ such that $\eta=\phi \circ \eta^{\prime}$.
Proof. By Dieudonné theory, the Dieudonné submodule

$$
\bigoplus_{i \in \mathbb{Z} / f \mathbb{Z}}\left(\tilde{\mathcal{E}}_{i} / p^{m} \tilde{\mathcal{D}}(A)_{i}^{\circ}\right)^{\oplus n} \subseteq \bigoplus_{i \in \mathbb{Z} / f \mathbb{Z}}\left(\tilde{\mathcal{D}}(A)_{i}^{\circ} / p^{m} \tilde{\mathcal{D}}(A)_{i}^{\circ}\right)^{\oplus n}
$$

corresponds to a closed subgroup scheme $H_{\mathfrak{p}} \subseteq A\left[p^{m}\right]$. The prime-to- $p$ polarization $\lambda$ induces a perfect pairing

$$
\langle-,-\rangle_{\lambda}: A\left[\mathfrak{p}^{m}\right] \times A\left[\overline{\mathfrak{p}}^{m}\right] \rightarrow \mu_{p^{m}} .
$$

Let $H_{\overline{\mathfrak{p}}}=H_{\mathfrak{p}}^{\perp} \subseteq A\left[\overline{\mathfrak{p}}^{m}\right]$ denote the orthogonal complement of $H_{\mathfrak{p}}$. Put $H_{p}=H_{\mathfrak{p}} \oplus H_{\overline{\mathfrak{p}}}$. Let $\psi: A \rightarrow A^{\prime}$ be the canonical quotient with kernel $H_{p}$, and $\phi: A^{\prime} \rightarrow A$ be the quotient with kernel $\psi\left(A\left[\mathfrak{p}^{m}\right]\right)$ so that $\psi \circ \phi=p^{m} \mathrm{id}_{A^{\prime}}$ and $\phi \circ \psi=p^{m} \mathrm{id}_{A}$. By construction, $H_{p} \subseteq A\left[p^{m}\right]$ is a maximal totally isotropic subgroup. By [Mumford 2008, $\S 23$, Theorem 2], there is a prime-to- $p$ polarization $\lambda^{\prime}$ on $A^{\prime}$ such that $p^{m} \lambda=$ $\psi^{\vee} \circ \lambda^{\prime} \circ \psi$. It follows also that $p^{m} \lambda^{\prime}=\phi^{\vee} \circ \lambda \circ \phi$. The fact that $\phi_{*, i}: \tilde{\mathcal{D}}\left(A^{\prime}\right)_{i}^{\circ} \rightarrow \tilde{\mathcal{D}}(A)_{i}^{\circ}$ is identified with the natural inclusion $\tilde{\mathcal{E}}_{i} \subseteq \tilde{\mathcal{D}}(A)_{i}^{\circ}$ follows from the construction. The existence and uniqueness of the tame level structure is clear. The dimension of the differential forms can be computed as follows:

$$
\begin{aligned}
\operatorname{dim}_{k} \omega_{A^{\prime} / k, i}^{\circ} & =\operatorname{dim}_{k} \frac{V\left(\tilde{\mathcal{D}}\left(A^{\prime}\right)_{i+1}^{\circ}\right)}{p \tilde{\mathcal{D}}\left(A^{\prime}\right)_{i}^{\circ}}=\operatorname{dim}_{k} \frac{V\left(\tilde{\mathcal{E}}_{i+1}\right)}{p \tilde{\mathcal{E}}_{i}} \\
& =\operatorname{dim}_{k} \frac{V\left(\tilde{\mathcal{D}}(A)_{i+1}^{\circ}\right)}{p \tilde{\mathcal{D}}(A)_{i}^{\circ}}-\operatorname{length}_{W(k)} \frac{V\left(\tilde{\mathcal{D}}(A)_{i+1}^{\circ}\right)}{V\left(\tilde{\mathcal{E}}_{i+1}\right)}+\text { length }_{W(k)} \frac{p \tilde{\mathcal{D}}(A)_{i}^{\circ}}{p \tilde{\mathcal{E}}_{i}} \\
& =a_{i}-\ell_{i+1}+\ell_{i} .
\end{aligned}
$$

3.3. Deformation theory. We shall frequently use Grothendieck-Messing deformation theory to compare the tangent spaces of moduli spaces. We make this explicit in our setup.

Let $\hat{R}$ be a noetherian $\mathbb{F}_{p^{f} \text {-algebra }}$ and $\hat{I} \subset \hat{R}$ an ideal such that $\hat{I}^{2}=0$. Put $R=\hat{R} / \hat{I}$. Let $\mathscr{C}_{\hat{R}}$ denote the category of tuples $(\hat{A}, \hat{\lambda}, \hat{\eta})$, where $\hat{A}$ is an $f n^{2}$ dimensional abelian variety over $\hat{R}$ equipped with an $\mathcal{O}_{D}$-action, $\hat{\lambda}$ is a polarization
on $\hat{A}$ such that the Rosati involution induces the $*$-involution on $\mathcal{O}_{D}$, and $\hat{\eta}$ is a level structure as in Section 2.3(3). We define $\mathscr{C}_{R}$ in the same way. For an object $(A, \lambda, \eta)$ in the category $\mathscr{C}_{R}$, let $H_{1}^{\text {cris }}(A / \hat{R})$ be the evaluation of the first relative crystalline homology (i.e., dual crystal of the first crystalline cohomology) of $A / R$ at the divided power thickening $\hat{R} \rightarrow R$, and $H_{1}^{\text {cris }}(A / \hat{R})_{i}^{\circ}:=\mathfrak{e} H_{1}^{\text {cris }}(A / \hat{R})_{q_{i}}$ be the $i$-th reduced part. We denote by $\mathscr{D e f}(R, \hat{R})$ the category of tuples $\left(A, \lambda, \eta,\left(\hat{\omega}_{i}^{\circ}\right)_{i=1, \ldots, f}\right)$, where $(A, \lambda, \eta)$ is an object in $\mathscr{C}_{R}$, and $\hat{\omega}_{i}^{\circ} \subseteq H_{1}^{\text {cris }}(A / \hat{R})_{i}^{\circ}$ for each $i \in \mathbb{Z} / f \mathbb{Z}$ is a subbundle that lifts $\omega_{A^{\vee} / R, i}^{\circ} \subseteq H_{1}^{\mathrm{dR}}(A / R)_{i}^{\circ}$. The following is a combination of Serre-Tate and Grothendieck-Messing deformation theory.
Theorem 3.4 (Serre-Tate, Grothendieck-Messing). The functor

$$
(\hat{A}, \hat{\lambda}, \hat{\eta}) \mapsto\left(\hat{A} \otimes_{\hat{R}} R, \lambda, \eta, \omega_{\hat{A}^{\vee} / \hat{R}, i}^{\circ}\right),
$$

where $\lambda$ and $\eta$ are the natural induced polarization and level structure on $\hat{A} \otimes_{\hat{R}} R$, is an equivalence of categories between $\mathscr{C}_{\hat{R}}$ and $\mathscr{D} \operatorname{ef}(R, \hat{R})$.
Proof. The main theorem of the crystalline deformation theory (cf., [Grothendieck 1974, pp. 116-118], [Mazur and Messing 1974, Chapter II §1]) says that the category $\mathscr{C}_{\hat{R}}$ is equivalent to the category of objects $(A, \lambda, \eta)$ in $\mathscr{C}_{R}$ together with a lift of $\omega_{A^{\vee} / R} \subseteq H_{1}^{\text {cris }}(A / R)$ to a subbundle $\hat{\omega}$ of $H_{1}^{\text {cris }}(A / \hat{R})$, such that $\hat{\omega}$ is stable under the induced $\mathcal{O}_{D}$-action and is isotropic for the pairing on $H_{1}^{\text {cris }}(A / \hat{R})$ induced by the polarization $\lambda$. But the additional information $\hat{\omega}$ is clearly equivalent to the subbundles $\hat{\omega}_{i}^{\circ} \subseteq H_{1}^{\text {cris }}(A / \hat{R})_{i}^{\circ}$ lifting $\omega_{A^{\vee} / R, i}^{\circ}$.
Corollary 3.5. If $\mathcal{A}_{a_{0}}$ denotes the universal abelian variety over $\mathrm{Sh}_{a_{\bullet}}$, then the tangent space $T_{\mathrm{Sh}_{a_{0}}}$ of $\mathrm{Sh}_{a_{0}}$ is

$$
\bigoplus_{i=1}^{f} \operatorname{Lie}_{\mathcal{A}_{a_{\bullet}}^{\vee} / \mathrm{Sh}_{a_{\bullet}}, i}^{\circ} \otimes \mathrm{Lie}_{\mathcal{A}_{a_{\bullet}} / \mathrm{Sh}_{a_{\bullet}}, i}^{\circ}
$$

Proof. Even though this is a well-known statement often referred to as the KodairaSpencer isomorphism (e.g., [Lan 2013, Proposition 2.3.4.2]), we include a short proof, as the proof serves as a toy model of many arguments later. Let $\hat{R}$ be
 By Theorem 3.4, to lift an $R$-point $(A, \lambda, \eta)$ of $\mathrm{Sh}_{a_{0}}$ to an $\hat{R}$-point, it suffices to lift, for $i=1, \ldots, f$, the differentials $\omega_{A^{\vee}, i}^{\circ} \subseteq H_{1}^{\text {cris }}(A / R)_{i}^{\circ}$ to a subbundle $\hat{\omega}_{i} \subseteq H_{1}^{\text {cris }}(A / \hat{R})_{i}^{\circ}$. Such lifts form a torsor for the group

$$
\operatorname{Hom}_{R}\left(\omega_{A^{\vee} / R, i}^{\circ}, \operatorname{Lie}_{A / R, i}^{\circ}\right) \otimes_{R} \hat{I} .
$$

It follows from this

$$
T_{\mathrm{Sh}_{a_{\bullet}}} \cong \bigoplus_{i=1}^{f} \mathcal{H o m}\left(\omega_{\mathcal{A}_{a_{\bullet}}^{\vee} / \mathrm{Sh}_{a_{0}}, i}^{\circ}, \mathrm{Lie}_{\mathcal{A}_{a_{0}} / \mathrm{Sh}_{a_{\bullet}}, i}^{\circ} \cong \bigoplus_{i=1}^{f} \mathrm{Lie}_{\mathcal{A}_{a_{0}}^{\vee} / \mathrm{Sh}_{a_{0}}, i} \otimes \mathrm{Lie}_{\mathcal{A}_{a_{0}} / \mathrm{Sh}_{a_{\bullet}}, i}^{\circ}\right.
$$

Note that this proof also shows that $\mathrm{Sh}_{a_{0}}$ is smooth.
3.6. Notation in the real quadratic case. For the rest of the paper, we assume $f=2$ so that $F$ is a real quadratic field in which $p$ is inert. For nonnegative integers $r \leq s$ such that $n=r+s$, we denote by $G_{r, s}$ the algebraic group previously denoted by $G_{a_{0}}$ with $a_{1}=r$ and $a_{2}=s$; in particular, $G_{r, s}(\mathbb{R})=G(U(r, s) \times U(s, r))$. If $r^{\prime}, s^{\prime}$ is another pair of nonnegative integers such that $n=r^{\prime}+s^{\prime}$ and $r^{\prime} \leq s^{\prime}$, Lemma 2.9 gives an isomorphism $G_{r, s}\left(\mathbb{A}^{\infty}\right) \simeq G_{r^{\prime}, s^{\prime}}\left(\mathbb{A}^{\infty}\right)$.

Let $\mathcal{S} h_{r, s}$ be the Shimura variety over $\mathbb{Z}_{p^{2}}$ attached to $G_{r, s}$ defined in Section 2.3 of some fixed sufficiently small prime-to- $p$ level $K^{p} \subseteq G_{r, s}\left(\mathbb{A}^{\infty, p}\right)$. Let $\mathrm{Sh}_{r, s}$ denote its special fiber over $\mathbb{F}_{p^{2}}$. Let $\mathcal{A}=\mathcal{A}_{r, s}$ denote the universal abelian variety over $\mathrm{Sh}_{r, s}$. It is a $2 n^{2}$-dimensional abelian variety, equipped with an action of $\mathcal{O}_{D}$ and a prime-to- $p$ polarization $\lambda_{\mathcal{A}}$. Moreover, $\omega_{\mathcal{A}^{\vee} / \mathrm{Sh}_{r, s}, 1}^{\circ}\left(\right.$ resp. $\left.\omega_{\mathcal{A}^{\vee} / \mathrm{Sh}_{r, s}, 2}^{\circ}\right)$ is a locally free module over $\mathrm{Sh}_{r, s}$ of rank $r$ (resp. rank $s$ ).

Remark 3.7. When $r=0$ and $s=n$, the universal abelian variety $\mathcal{A}=\mathcal{A}_{0, n}$ over $\mathrm{Sh}_{0, n}$ is supersingular. Indeed, for each $\overline{\mathbb{F}}_{p}$-point $z$ of $\mathrm{Sh}_{0, n}$, the Kottwitz condition implies that the Frobenius induces isomorphisms

$$
\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \xrightarrow{F} \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \xrightarrow{F} p \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} .
$$

In particular, $(1 / p) F^{2}$ induces a $\sigma^{2}$-linear automorphism of $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}$. By Hilbert's Theorem 90 , there exists a $\mathbb{Z}_{p^{2}}$-lattice $\mathbb{L}$ of $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}$ that is invariant under the action of $(1 / p) F^{2}$; in other words, $F^{2}$ acts by multiplication by $p$ for a basis chosen from this lattice. It follows that all slopes of the Frobenius on $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)$ are $\frac{1}{2}$, and hence $\mathcal{A}_{z}$ is supersingular.

## 4. The case of $G(U(1, n-1) \times U(n-1,1))$

We will verify Conjecture 2.12 for $\mathrm{Sh}_{1, n-1}$, namely the existence of some cycles $Y_{j}$ having morphisms to both $\mathrm{Sh}_{0, n}$ and $\mathrm{Sh}_{1, n-1}$ and generating Tate classes of $\mathrm{Sh}_{1, n-1}$ under a certain genericity hypothesis on the Satake parameters. We always fix an isomorphism $G_{1, n-1}\left(\mathbb{A}^{\infty}\right) \simeq G_{0, n}\left(\mathbb{A}^{\infty}\right)$, and write $G\left(\mathbb{A}^{\infty}\right)$ for either group.

Notation 4.1. For a smooth variety $X$ over $\mathbb{F}_{p^{2}}$, we denote by $T_{X}$ the tangent bundle of $X$, and for a locally free $\mathcal{O}_{X}$-module $M$, we put $M^{*}=\mathcal{H o m}_{\mathcal{O}_{X}}\left(M, \mathcal{O}_{X}\right)$.
4.2. Cycles on $\mathbf{S h}_{1, n-1}$. For each integer $j$ with $1 \leq j \leq n$, we first define the variety $Y_{j}$ we briefly mentioned in the introduction. Let $Y_{j}$ be the moduli space over $\mathbb{F}_{p^{2}}$ that associates to each locally noetherian $\mathbb{F}_{p^{2}}$-scheme $S$, the set of isomorphism classes of tuples $\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$, where

- $(A, \lambda, \eta)$ is an $S$-point of $\mathrm{Sh}_{1, n-1}$,
- $\left(B, \lambda^{\prime}, \eta^{\prime}\right)$ is an $S$-point of $\mathrm{Sh}_{0, n}$ and
- $\phi: B \rightarrow A$ is an $\mathcal{O}_{D}$-equivariant isogeny whose kernel is contained in $B[p]$,
such that
- $p \lambda^{\prime}=\phi^{\vee} \circ \lambda \circ \phi$,
- $\phi \circ \eta^{\prime}=\eta$ and
- the cokernels of the maps

$$
\phi_{*, 1}: H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ} \rightarrow H_{1}^{\mathrm{dR}}(A / S)_{1}^{\circ} \quad \text { and } \quad \phi_{*, 2}: H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ} \rightarrow H_{1}^{\mathrm{dR}}(A / S)_{2}^{\circ}
$$

are locally free $\mathcal{O}_{S}$-modules of rank $j-1$ and $j$, respectively.
There is a unique isogeny $\psi: A \rightarrow B$ such that $\psi \circ \phi=p \cdot \mathrm{id}_{B}$ and $\phi \circ \psi=p \cdot \mathrm{id}_{A}$. We have

$$
\operatorname{Ker}\left(\phi_{*, i}\right)=\operatorname{Im}\left(\psi_{*, i}\right) \quad \text { and } \quad \operatorname{Ker}\left(\phi_{*, i}\right)=\operatorname{Im}\left(\psi_{*, i}\right),
$$

where $\psi_{*, i}$ for $i=1,2$ is the induced homomorphism on the reduced de Rham homology in the evident sense. This moduli space $Y_{j}$ is represented by a scheme of finite type over $\mathbb{F}_{p^{2}}$. We have a natural diagram of morphisms:

where $\mathrm{pr}_{j}$ and $\operatorname{pr}_{j}^{\prime}$ send a tuple $\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$ to $(A, \lambda, \eta)$ and to $\left(B, \lambda^{\prime}, \eta^{\prime}\right)$, respectively. Letting $K^{p}$ vary, we see easily that both $\mathrm{pr}_{j}$ and $\mathrm{pr}_{j}^{\prime}$ are equivariant under prime-to- $p$ Hecke actions given by the double cosets $K^{p} \backslash G\left(\mathbb{A}^{\infty}, p\right) / K^{p}$.
4.3. Some auxiliary moduli spaces. The moduli problem for $Y_{j}$ is slightly complicated. We will introduce a more explicit moduli space $Y_{j}^{\prime}$ below and then show they are isomorphic.

Consider the functor $\underline{Y}_{j}^{\prime}$ which associates to each locally noetherian $\mathbb{F}_{p^{2}}$-scheme $S$ the set of isomorphism classes of tuples ( $B, \lambda^{\prime}, \eta^{\prime}, H_{1}, H_{2}$ ), where

- $\left(B, \lambda^{\prime}, \eta^{\prime}\right)$ is an $S$-valued point of $\mathrm{Sh}_{0, n}$;
- $H_{1} \subset H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ}$ and $H_{2} \subset H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ}$ are $\mathcal{O}_{S}$-subbundles of rank $j$ and $j-1$ respectively such that

$$
\begin{equation*}
V^{-1}\left(H_{2}^{(p)}\right) \subseteq H_{1}, \quad H_{2} \subseteq F\left(H_{1}^{(p)}\right) \tag{4.3.1}
\end{equation*}
$$

Here,

$$
F: H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ,(p)} \xrightarrow{\sim} H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ} \quad \text { and } \quad V: H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ} \xrightarrow{\sim} H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ,(p)}
$$

are respectively the Frobenius and Verschiebung homomorphisms, which are actually isomorphisms because of the signature condition on $\mathrm{Sh}_{0, n}$.

It follows from the moduli problem that the quotients $H_{1} / V^{-1}\left(H_{2}^{(p)}\right), F\left(H_{1}^{(p)}\right) / H_{2}$ are both locally free $\mathcal{O}_{S}$-modules of rank one.

There is a natural projection $\pi_{j}^{\prime}: \underline{Y}_{j}^{\prime} \rightarrow \mathrm{Sh}_{0, n}$ given by $\left(B, \lambda^{\prime}, \eta^{\prime}, H_{1}, H_{2}\right) \mapsto$ ( $B, \lambda^{\prime}, \eta^{\prime}$ ).

Proposition 4.4. The functor $\underline{Y}_{j}^{\prime}$ is representable by a scheme $Y_{j}^{\prime}$ smooth and projective over $\mathrm{Sh}_{0, n}$ of dimension $n-1$. Moreover, if $\left(\mathcal{B}, \lambda^{\prime}, \eta^{\prime}, \mathcal{H}_{1}, \mathcal{H}_{2}\right)$ denotes the universal object over $Y_{j}^{\prime}$, then the tangent bundle of $Y_{j}^{\prime}$ is

$$
T_{Y_{j}^{\prime}} \cong\left(\left(\mathcal{H}_{1} / V^{-1}\left(\mathcal{H}_{2}^{(p)}\right)\right)^{*} \otimes\left(H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \mathrm{Sh}_{0, n}\right)_{1}^{\circ} / \mathcal{H}_{1}\right) \oplus\left(\mathcal{H}_{2}^{*} \otimes F\left(\mathcal{H}_{1}^{(p)}\right) / \mathcal{H}_{2}\right)\right.
$$

Proof. For each integer $m$ with $0 \leq m \leq n$ and $i=1$, 2, let $\boldsymbol{G r}\left(H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \mathrm{Sh}_{0, n}\right)_{i}^{\circ}, m\right)$ be the Grassmannian scheme over $\mathrm{Sh}_{0, n}$ that parametrizes subbundles of the universal de Rham homology $H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \mathrm{Sh}_{0, n}\right)_{i}^{\circ}$ of rank $m$. Then $\underline{Y}_{j}^{\prime}$ is a closed subfunctor of the product of the Grassmannian schemes

$$
\boldsymbol{G r}\left(H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \mathrm{Sh}_{0, n}\right)_{1}^{\circ}, j\right) \times \boldsymbol{G} \boldsymbol{r}\left(H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \mathrm{Sh}_{0, n}\right)_{2}^{\circ}, j-1\right)
$$

The representability of $\underline{Y}_{j}^{\prime}$ follows. Moreover, $Y_{j}^{\prime}$ is projective.
We show now that the structural map $\pi_{j}^{\prime}: Y_{j}^{\prime} \rightarrow \mathrm{Sh}_{0, n}$ is smooth of relative dimension $n-1$. Let $S_{0} \hookrightarrow S$ be an immersion of locally noetherian $\mathbb{F}_{p^{2}}$-schemes with ideal sheaf $I$ satisfying $I^{2}=0$. Suppose we are given a commutative diagram

with solid arrows. We have to show that, locally for the Zariski topology on $S_{0}$, there is a morphism $g: S \rightarrow Y_{j}^{\prime}$ making the diagram commute. Let $B$ be the abelian scheme over $S$ given by $h$, and $B_{0}$ be the base change to $S_{0}$. The morphism $g_{0}$ gives rises to subbundles $\bar{H}_{1} \subset H_{1}^{\mathrm{dR}}\left(B_{0} / S_{0}\right)_{1}^{\circ}$ and $\bar{H}_{2} \subset H_{1}^{\mathrm{dR}}\left(B_{0} / S_{0}\right)_{2}^{\circ}$ with

$$
F\left(\bar{H}_{1}^{(p)}\right) \supset \bar{H}_{2}, \quad V^{-1}\left(\bar{H}_{2}^{(p)}\right) \subset \bar{H}_{1}
$$

Finding $g$ is equivalent to finding a subbundle $H_{i} \subset H_{1}^{\mathrm{dR}}(B / S)_{i}^{\circ}$ which lifts each $\bar{H}_{i}$ for $i=1,2$ and satisfies (4.3.1); this is certainly possible when passing to small enough affine open subsets of $S_{0}$. Thus $\pi_{j}^{\prime}: Y_{j}^{\prime} \rightarrow \mathrm{Sh}_{0, n}$ is formally smooth, and hence smooth. We note that $F_{S}^{*}: \mathcal{O}_{S} \rightarrow \mathcal{O}_{S}$ factors through $\mathcal{O}_{S_{0}}$. Hence $V^{-1}\left(H_{2}^{(p)}\right)$ and $F\left(H_{1}^{(p)}\right)$ actually depend only on $\bar{H}_{1}, \bar{H}_{2}$, but not on the lifts $H_{1}$ and $H_{2}$. Therefore, the possible lifts $H_{2}$ form a torsor under the group

$$
\mathcal{H o m}_{\mathcal{O}_{S_{0}}}\left(\bar{H}_{2}, F\left(\bar{H}_{1}^{(p)}\right) / \bar{H}_{2}\right) \otimes_{\mathcal{O}_{S_{0}}} I
$$

and similarly the possible lifts $H_{1}$ form a torsor under the group

$$
\mathcal{H o m}_{\mathcal{O}_{S_{0}}}\left(\bar{H}_{1} / V^{-1}\left(\bar{H}_{2}^{(p)}\right), H_{1}^{\mathrm{dR}}\left(B_{0} / S_{0}\right)_{1}^{\circ} / \bar{H}_{1}\right) \otimes_{\mathcal{O}_{S_{0}}} I
$$

To compute the tangent bundle $T_{Y_{j}^{\prime}}$, we take $S=\operatorname{Spec}\left(\mathcal{O}_{S_{0}}[\epsilon] / \epsilon^{2}\right)$ and $I=\epsilon \mathcal{O}_{S}$. The morphism $g_{0}: S_{0} \rightarrow Y_{j}^{\prime}$ corresponds to an $S_{0}$-valued point of $Y_{j}^{\prime}$, say $y_{0}$. Then the possible liftings $g$ form the tangent space $T_{Y_{j}^{\prime}}$ at $y_{0}$, denote by $T_{Y_{j}^{\prime}, y_{0}}$. The discussion above shows that
$T_{Y_{j}^{\prime}, y_{0}} \cong \mathcal{H o m}_{\mathcal{O}_{S_{0}}}\left(\bar{H}_{2}, F\left(\bar{H}_{1}^{(p)}\right) / \bar{H}_{2}\right) \oplus \mathcal{H o m}_{\mathcal{O}_{S_{0}}}\left(\bar{H}_{1} / V^{-1}\left(\bar{H}_{2}^{(p)}\right), H_{1}^{\mathrm{dR}}\left(B_{0} / S_{0}\right)_{1}^{\circ} / \bar{H}_{1}\right)$, which is certainly a vector bundle over $S_{0}$ of rank $j-1+(n-j)=n-1$. Applying this to the universal case when $g_{0}: S_{0} \rightarrow Y_{j}^{\prime}$ is the identity morphism, the formula of the tangent bundle follows.
Remark 4.5. Let $\left(B, \lambda^{\prime}, \eta^{\prime}, H_{1}, H_{2}\right)$ be an $S$-point of $Y_{j}^{\prime}$.
(a) If $j=n, H_{1}$ has to be $H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ}$, and $H_{2}$ is a hyperplane of $H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ}$. Condition (4.3.1) is trivial. In this case, $Y_{n}^{\prime}$ is the projective space over $\mathrm{Sh}_{0, n}$ associated to $H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \mathrm{Sh}_{0, n}\right)_{2}^{\circ}$, where $\mathcal{B}$ is the universal abelian scheme over $\mathrm{Sh}_{0, n}$. So it is geometrically a union of copies of $\mathbb{P}_{\mathbb{F}_{p}}^{n-1}$.
(b) If $j=1$, then $H_{1}$ is a line in $H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ}$ and $H_{2}=0$. So $Y_{1}^{\prime}$ is the projective space over $\mathrm{Sh}_{0, n}$ associated to $\left(H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \mathrm{Sh}_{0, n}\right)_{1}^{\circ}\right)^{*}$.
(c) If $j=2, H_{2} \subseteq H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ}$ is a line, and $H_{1} \subseteq H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ}$ is a subbundle of rank 2 such that $F\left(H_{1}^{(p)}\right)$ contains both $H_{2}$ and $F\left(V^{-1}\left(H_{2}^{(p)}\right)^{(p)}\right)$. Therefore, if $H_{2} \neq F\left(V^{-1}\left(H_{2}^{(p)}\right)\right)^{(p)}, H_{1}$ is determined up to Frobenius pullback. If $H_{2}=$ $F\left(V^{-1}\left(H_{2}^{(p)}\right)^{(p)}\right)$, then $H_{1}$ could be any rank 2 subbundle containing $V^{-1}\left(H_{2}^{(p)}\right)$.

We fix a geometric point $z=\left(B, \lambda^{\prime}, \eta^{\prime}\right) \in \mathrm{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$. It is possible to find good bases for $H_{1}^{\mathrm{dR}}\left(B / \overline{\mathbb{F}}_{p}\right)_{1}^{\circ}, H_{1}^{\mathrm{dR}}\left(B / \overline{\mathbb{F}}_{p}\right)_{2}^{\circ}$ such that $F, V: H_{1}^{\mathrm{dR}}\left(B / \overline{\mathbb{F}}_{p}\right)_{1}^{\circ} \rightarrow H_{1}^{\mathrm{dR}}\left(B / \overline{\mathbb{F}}_{p}\right)_{2}^{\circ}$ are both given by the identity matrix. With these choices, we may identify the fiber $Y_{2, z}^{\prime}=\pi_{2}^{\prime-1}(z)$ with a closed subvariety of

$$
\boldsymbol{G r}\left(\overline{\mathbb{F}}_{p}^{n}, 2\right) \times \boldsymbol{G r}\left(\overline{\mathbb{F}}_{p}^{n}, 1\right)
$$

Moreover, one may equip $\boldsymbol{\operatorname { G r }}\left(\overline{\mathbb{F}}_{p}^{n}, 1\right) \cong \mathbb{P}_{\overline{\mathbb{F}}_{p}}^{n-1}$ with an $\mathbb{F}_{p^{2}}$-rational structure such that $H_{2}=F\left(V^{-1}\left(H_{2}^{(p)}\right)^{(p)}\right)$ if and only if $\left[H_{2}\right] \in \mathbb{P}_{\overline{\mathbb{F}}_{p}}^{n-1}$ is an $\mathbb{F}_{p^{2}}$-rational point. So $Y_{2, z}^{\prime}$ is isomorphic to a "Frobenius twisted" blow-up of $\mathbb{P}_{\mathbb{F}_{p}}^{n-1}$ at all of its $\mathbb{F}_{p^{2}}$-rational points. Here, "Frobenius twisted" means that each irreducible component of the exceptional divisor has multiplicity $p$. For instance, when $n=3$, each $Y_{2, z}$ is isomorphic to the closed subscheme of $\mathbb{P}_{\mathbb{F}_{p}}^{2} \times \mathbb{P}_{\mathbb{F}_{p}}^{2}$ defined by

$$
a_{1} b_{1}^{p}+a_{2} b_{2}^{p}+a_{3} b_{3}^{p}=0, \quad a_{1}^{p} b_{1}+a_{2}^{p} b_{2}+a_{3}^{p} b_{3}=0
$$

where $\left(a_{1}: a_{2}: a_{3}\right)$ and $\left(b_{1}: b_{2}: b_{3}\right)$ are the homogeneous coordinates on the two copies of $\mathbb{P}^{2}$.

Lemma 4.6. Let $\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$ be an $S$-point of $Y_{j}$. Then the image of $\phi_{*, 1}$ contains both $\omega_{A^{\vee} / S, 1}^{\circ}$ and $F\left(H_{1}^{\mathrm{dR}}(A / S)_{2}^{\circ,(p)}\right)$, and the image of $\phi_{*, 2}$ is contained in both $\omega_{A^{\vee} / S, 2}^{\circ}$ and $F\left(H_{1}^{\mathrm{dR}}(A / S)_{1}^{\circ,(p)}\right)$.
Proof. By the functoriality, $\phi_{2, *}$ sends $\omega_{B^{\vee} / S, 2}^{\circ}$ to $\omega_{A^{\vee} / S, 2}^{\circ}$. Since $\omega_{B^{\vee} / S, 2}^{\circ}=$ $H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ}$ by the Kottwitz determinant condition, it follows that $\operatorname{Im}\left(\phi_{*, 2}\right)$ in contained in $\omega_{A^{\vee} / S, 2}^{\circ}$. Similar arguments by considering $\psi_{*, 1}$ shows that $\omega_{A^{\vee} / S, 1}^{\circ} \subseteq$ $\operatorname{Ker}\left(\psi_{*, 1}\right)=\operatorname{Im}\left(\phi_{*, 1}\right)$. The fact that $\operatorname{Im}\left(\phi_{*, 2}\right)$ is contained in $F\left(H_{1}^{\mathrm{dR}}(A / S)_{1}^{\mathrm{o},(p)}\right)$ follows from the commutative diagram

$$
\begin{align*}
& H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ,(p)} \xrightarrow{\phi_{*, 1}^{(p)}} H_{1}^{\mathrm{dR}}(A / S)_{1}^{\circ,(p)}  \tag{4.6.1}\\
& F \downarrow \\
& H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ} \xrightarrow{\cong} \xrightarrow{\phi_{*, 2}} \underset{ }{\square} H_{1}^{\mathrm{dR}}(A / S)_{2}^{\circ}
\end{align*}
$$

and the fact that the left vertical arrow is an isomorphism. Similarly, the inclusion $F\left(H_{1}^{\mathrm{dR}}(A / S)_{2}^{\circ,(p)}\right) \subseteq \operatorname{Im}\left(\phi_{*, 1}\right)=\operatorname{Ker}\left(\psi_{*, 1}\right)$ can be proved using the functoriality of Verschiebung homomorphisms.
4.7. A morphism from $\boldsymbol{Y}_{\boldsymbol{j}}$ to $\boldsymbol{Y}_{\boldsymbol{j}}^{\prime}$. There is a natural morphism $\alpha: Y_{j} \rightarrow Y_{j}^{\prime}$ for $1 \leq j \leq n$ defined as follows. For a locally noetherian $\mathbb{F}_{p^{2}}$-scheme $S$ and an $S$-point $\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$ of $Y_{j}$, we define

$$
\begin{equation*}
H_{1}:=\phi_{*, 1}^{-1}\left(\omega_{A^{\vee} / S, 1}^{\circ}\right) \subseteq H_{1}^{\mathrm{dR}}(B / S)_{1}^{\circ}, \quad H_{2}:=\psi_{*, 2}\left(\omega_{A^{\vee} / S, 2}^{\circ}\right) \subseteq H_{1}^{\mathrm{dR}}(B / S)_{2}^{\circ} \tag{4.7.1}
\end{equation*}
$$

In particular, $H_{1}$ and $H_{2}$ are $\mathcal{O}_{S}$-subbundles of rank $j$ and $j-1$, respectively. Also, there is a canonical isomorphism $\omega_{A^{\vee} / S, 2}^{\circ} / \operatorname{Im}\left(\phi_{*, 2}\right) \xrightarrow{\sim} H_{2}$. From the commutative diagram (4.6.1), it is easy to see that $F\left(H_{1}^{(p)}\right) \subseteq \operatorname{Ker}\left(\phi_{*, 2}\right)=\operatorname{Im}\left(\psi_{*, 2}\right)$, but comparing the rank forces this to be an equality. It follows that $H_{2} \subseteq F\left(H_{1}^{(p)}\right)$. Similarly, $V^{-1}\left(H_{2}^{(p)}\right)$ is identified with $\operatorname{Im}\left(\psi_{*, 1}\right)=\operatorname{Ker}\left(\phi_{*, 1}\right)$, hence $V^{-1}\left(H_{2}^{(p)}\right) \subseteq H_{1}$. From these, we deduce two canonical isomorphisms:

$$
\begin{align*}
H_{1} / V^{-1}\left(H_{2}^{(p)}\right) & \xrightarrow{\hookrightarrow} \omega_{A^{\vee} / S, 1}^{\circ}  \tag{4.7.2}\\
F\left(H_{1}^{(p)}\right) / H_{2} & \xrightarrow{\longrightarrow} H_{1}^{\mathrm{dR}}(A / S)_{2}^{\circ} / \omega_{A^{\vee} / S, 2}^{\circ} \cong \mathrm{Lie}_{A / S, 2}^{\circ} .
\end{align*}
$$

Therefore, we have a well-defined map $\alpha: Y_{j} \rightarrow Y_{j}^{\prime}$ given by

$$
\alpha:\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right) \mapsto\left(B, \lambda^{\prime}, \eta^{\prime}, H_{1}, H_{2}\right)
$$

Moreover, it is clear from the definition that $\pi_{j}^{\prime} \circ \alpha=\mathrm{pr}_{j}^{\prime}$.
Proposition 4.8. The morphism $\alpha$ is an isomorphism.
Proof. Let $k$ be a perfect field containing $\mathbb{F}_{p^{2}}$. We first prove that $\alpha$ induces a bijection of points $\alpha: Y_{j}(k) \xrightarrow{\sim} Y_{j}^{\prime}(k)$. It suffices to show that there exists a
morphism of sets $\beta: Y_{j}^{\prime}(k) \rightarrow Y_{j}(k)$ inverse to $\alpha$. Let $y=\left(B, \lambda^{\prime}, \eta^{\prime}, H_{1}, H_{2}\right) \in$ $Y_{j}^{\prime}(k)$. We define $\beta(y)=\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$ as follows. Let $\tilde{\mathcal{E}}_{1} \subseteq \tilde{\mathcal{D}}(B)_{1}^{\circ}$ and $\tilde{\mathcal{E}}_{2} \subseteq \tilde{\mathcal{D}}(B)_{2}^{\circ}$ be respectively the inverse images of $V^{-1}\left(H_{2}^{(p)}\right) \subseteq H_{1}^{\mathrm{dR}}(B / k)_{1}^{\circ}$ and $F\left(H_{1}^{(p)}\right) \subseteq H_{1}^{\mathrm{dR}}(B / k)_{2}^{\circ}$ under the natural reduction maps

$$
\tilde{\mathcal{D}}(B)_{i}^{\circ} \rightarrow \tilde{\mathcal{D}}(B)_{i}^{\circ} / p \tilde{\mathcal{D}}(B)_{i}^{\circ} \cong H_{1}^{\mathrm{dR}}(B / k)_{i}^{\circ} \quad \text { for } i=1,2
$$

The condition (4.3.1) ensures that $F\left(\tilde{\mathcal{E}}_{i}\right) \subseteq \tilde{\mathcal{E}}_{3-i}$ and $V\left(\tilde{\mathcal{E}}_{i}\right) \subseteq \tilde{\mathcal{E}}_{3-i}$ for $i=1,2$. Applying Proposition 3.2 with $m=1$, we get a triple $(A, \lambda, \eta)$ and an $\mathcal{O}_{D}$-equivariant isogeny $\psi: A \rightarrow B$, where $A$ is an abelian variety over $k$ with an action of $\mathcal{O}_{D}, \lambda$ is a prime-to- $p$ polarization on $A$, and $\eta$ is a prime-to- $p$ level structure on $A$, such that $\psi^{\vee} \circ \lambda^{\prime} \circ \psi=p \lambda, p \eta^{\prime}=\psi \circ \eta$ and such that $\psi_{*, i}: \tilde{\mathcal{D}}(A)_{i}^{\circ} \rightarrow \tilde{\mathcal{D}}(B)_{i}^{\circ}$ is naturally identified with the inclusion $\tilde{\mathcal{E}}_{i} \hookrightarrow \tilde{\mathcal{D}}(B)_{i}^{\circ}$ for $i=1$, 2. Moreover, the dimension formula (3.2.2) implies that $\omega_{A^{\vee} / k, 1}^{\circ}$ has dimension 1 , and $\omega_{A^{\vee} / k, 2}^{\circ}$ has dimension $n-1$. Therefore, $(A, \lambda, \eta)$ is a point of $\mathrm{Sh}_{1, n-1}$. Finally, we take $\phi: B \rightarrow A$ to be the unique isogeny such that $\phi \circ \psi=p \cdot \mathrm{id}_{A}$ and $\psi \circ \phi=p \cdot \mathrm{id}_{B}$. Thus we have $\phi \circ \eta^{\prime}=\eta$. This finishes the construction of $\beta(y)$. It is direct to check that $\beta$ is the set theoretic inverse to $\alpha: Y_{j}(k) \rightarrow Y_{j}^{\prime}(k)$.

We show now that $\alpha$ induces an isomorphism on the tangent spaces at each closed point; as we have already shown that $Y_{j}^{\prime}$ is smooth, it will then follow that $\alpha$ is an isomorphism. Let $x=\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right) \in Y_{j}(k)$ be a closed point. Consider the infinitesimal deformation over $k[\epsilon]=k[t] / t^{2}$. Note that ( $B, \lambda^{\prime}, \eta^{\prime}$ ) has a unique deformation $\left(\hat{B}, \hat{\lambda}^{\prime}, \hat{\eta}^{\prime}\right)$ to $k[\epsilon]$, namely the trivial deformation. By the Serre-Tate and Grothendieck-Messing deformation theory (cf., Theorem 3.4), giving a deformation $(\hat{A}, \hat{\lambda}, \hat{\eta})$ of $(A, \lambda, \eta)$ to $k[\epsilon]$ is equivalent to giving free $k[\epsilon]-$ submodules $\hat{\omega}_{A^{\vee}, i}^{\circ} \subseteq H_{1}^{\text {cris }}(A / k[\epsilon])_{i}^{\circ}$ for $i=1,2$ which lift $\omega_{A^{\vee} / k, i}^{\circ}$. The isogeny $\phi$ and the polarization $\lambda$ deform to an isogeny $\hat{\phi}: \hat{B} \rightarrow \hat{A}$ and a polarization $\hat{\lambda}: \hat{A}^{\vee} \rightarrow \hat{A}$ (satisfying $p \hat{\lambda}^{\prime}=\hat{\phi}^{\vee} \circ \hat{\lambda} \circ \hat{\phi}$ ), necessarily unique if they exist, if and only if

$$
\hat{\omega}_{A^{\vee}, 2}^{\circ} \supseteq \phi_{*, 2}^{\text {cris }}\left(H_{1}^{\text {cris }}(B / k[\epsilon])_{2}^{\circ}\right) \quad \text { and } \quad\left(\phi_{*, 1}^{\text {cris }}\left(H_{1}^{\text {cris }}(B / k[\epsilon])_{1}^{\circ}\right)\right)^{\vee} \subseteq\left(\hat{\omega}_{A^{\vee}, 1}^{\circ}\right)^{\vee}
$$

where the second inclusion comes from the consideration at the embedding $\bar{q}_{2}$ by taking duality using the polarization $\lambda$ and is equivalent to $\hat{\omega}_{A^{\vee}, 1}^{\circ} \subseteq \phi_{*, 1}^{\text {cris }}\left(H_{1}^{\text {cris }}(B / k[\epsilon])_{1}^{\circ}\right)$.

As discussed before Proposition 4.8, we have $\operatorname{Ker}\left(\phi_{*, 1}\right)=V^{-1}\left(H_{2}^{(p)}\right)$ and $F\left(H_{1}^{(p)}\right)=\operatorname{Ker}\left(\phi_{*, 2}\right)=\operatorname{Im}\left(\psi_{*, 2}\right)$. Then according to the relation between $\omega_{A^{\vee} / k, i}^{\circ}$ and $H_{1}$ in (4.7.1), giving such $\hat{\omega}_{A^{\vee}, i}^{\circ}$ for $i=1,2$ is equivalent to lifting each $H_{i}$ to a free $k[\epsilon]$-submodule $\hat{H}_{i} \subseteq H_{1}^{\mathrm{dR}^{\prime}}(B / k)_{i}^{\circ} \otimes_{k} k[\epsilon] \cong H_{1}^{\text {cris }}(B / k[\epsilon])_{i}^{\circ}$ for $i=1,2$ such that $\hat{H}_{1} \supseteq V^{-1}\left(H_{2}^{(p)}\right) \otimes_{k} k[\epsilon]$ and $\hat{H}_{2} \subseteq F\left(H_{1}^{(p)}\right) \otimes_{k} k[\epsilon]$. This is exactly the description of the tangent space of $Y_{j}^{\prime}$ at $\alpha(x)$. This concludes the proof.

In the sequel, we will always identify $Y_{j}$ with $Y_{j}^{\prime}$ and $\mathrm{pr}_{j}^{\prime}$ with $\pi_{j}^{\prime}$. Before proceeding, we prove some results on the structure of $\operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$.

We turn to the study of the Shimura variety $\mathcal{S} h_{0, n}$. The following proposition was suggested by an anonymous referee of this article.
Proposition 4.9. (1) The Shimura variety $\mathcal{S} h_{0, n}$ is finite and étale over $\mathbb{Z}_{p^{2}}$. In particular, the reduction map induces a bijection of geometric points

$$
\mathcal{S} h_{0, n}\left(\overline{\mathbb{Q}}_{p}\right) \xrightarrow{\sim} \mathrm{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right) .
$$

(2) Let $\tilde{x}_{i}=\left(\tilde{B}_{i}, \tilde{\lambda}_{i}, \tilde{\eta}_{i}\right) \in \mathcal{S} h_{0, n}\left(\overline{\mathbb{Q}}_{p}\right)$ for $i=1,2$ be two geometric points in characteristic 0 , and $x_{i}=\left(B_{i}, \lambda_{i}, \eta_{i}\right) \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ be their reductions. Then the reduction map on

$$
\operatorname{Hom}_{\mathcal{O}_{D}}\left(\tilde{B}_{1}, \tilde{B}_{2}\right) \xrightarrow{\sim} \operatorname{Hom}_{\mathcal{O}_{D}}\left(B_{1}, B_{2}\right)
$$

is an isomorphism.
Proof. (1) Let $\tilde{z} \in(\tilde{B}, \tilde{\lambda}, \tilde{\eta}) \in \mathcal{S} h_{0, n}(\mathbb{C})$. Put $H=H_{1}(\tilde{B}(\mathbb{C}), \mathbb{Q})$. It is a left $D$-module of rank 1 equipped with an alternating $D$-Hermitian pairing $\langle-,-\rangle_{\tilde{\lambda}}$ induced by the polarization $\tilde{\lambda}$. Let $\left(V_{0, n}=D,\langle-,-\rangle_{0, n}\right)$ be the left $D$-module together with its alternating $D$-Hermitian pairing as in the definition of $\mathcal{S} h_{0, n}$. By results of Kottwitz [1992b, §8], for every place $v$ of $\mathbb{Q}$, the skew-Hermitian $D_{\mathbb{Q}_{v}}$-modules $H_{\mathbb{Q}_{v}}$ and $V_{0, n, \mathbb{Q}_{v}}$ are isomorphic. ${ }^{16}$ Then $\operatorname{End}_{\mathcal{O}_{D}}\left(\tilde{B}_{\mathbb{C}}\right)_{\mathbb{Q}}$ consists of the elements of $D^{\mathrm{opp}}=\operatorname{End}_{D}(H)$ that preserves the complex structure on $H_{1, \mathbb{R}} \simeq V_{0, n, \mathbb{R}}$ induced the Deligne homomorphism by $h: \mathbb{C}^{\times} \rightarrow G_{0, n}(\mathbb{R})$. Since $h(i)$ is necessarily central (because $G_{0, n}^{1}$ is compact), it follows that $\operatorname{End}_{\mathcal{O}_{D}}\left(\tilde{B}_{\mathbb{C}}\right)_{\mathbb{Q}}=D^{\text {opp }}$, and

$$
D \otimes_{E} D^{\mathrm{opp}} \simeq \mathrm{M}_{n^{2}}(E) \subseteq \operatorname{End}(\tilde{B})_{\mathbb{Q}}
$$

For dimension reasons, the inclusion above is an equality, and $\tilde{B}$ is isogenous to the product of $n^{2}$-copies of abelian varieties with complex multiplication by $E$. Therefore, $\tilde{B}$ is defined over a number field and has potentially good reduction everywhere. This implies that $\mathcal{S} h_{0, n}$ is proper over $\mathbb{Z}_{p^{2}}$.

To see that $\mathcal{S} h_{0, n}$ is finite and étale over $\mathbb{Z}_{p^{2}}$, it remains to show its étaleness over $\mathbb{Z}_{p^{2}}$. But this is clear from the description of its relative differential sheaf in Corollary 3.5, which is trivial as $\operatorname{Lie}_{\mathcal{A}^{\vee} / \mathcal{S} h_{0, n}, 1}=\operatorname{Lie}_{\mathcal{A} / \mathcal{S} h_{0, n}, 2}=0$ by Kottwitz's determinant condition.
(2) In general, the reduction map

$$
\operatorname{Hom}_{\mathcal{O}_{D}}\left(\tilde{B}_{1}, \tilde{B}_{2}\right) \hookrightarrow \operatorname{Hom}_{\mathcal{O}_{D}}\left(B_{1}, B_{2}\right)
$$

is injective. It remains to see that every element $f \in \operatorname{Hom}_{\mathcal{O}_{D}}\left(B_{1}, B_{2}\right)$ lifts to a homomorphism $\tilde{f} \in \operatorname{Hom}_{\mathcal{O}_{D}}\left(\tilde{B}_{1}, \tilde{B}_{2}\right)$. Note that points $\tilde{x}_{1}, \tilde{x}_{2}$ can be viewed over

[^11]$W\left(\overline{\mathbb{F}}_{p}\right)$. As recalled in Section 3, to show that $f$ lifts to a map $\tilde{f}: \tilde{B}_{1} \rightarrow \tilde{B}_{2}$, it suffices to see that the induced map on crystalline homology
$$
f^{*}: H_{1}^{\text {cris }}\left(B_{2} / W\left(\overline{\mathbb{F}}_{p}\right)\right) \rightarrow H_{1}^{\text {cris }}\left(B_{1} / W\left(\overline{\mathbb{F}}_{p}\right)\right)
$$
preserves the Hodge filtrations
$$
\omega_{\tilde{B}_{i}^{\vee}} \subseteq H_{1}^{\mathrm{dR}}\left(\tilde{B}_{i} / W\left(\overline{\mathbb{F}}_{p}\right)\right) \cong H_{1}^{\text {cris }}\left(B_{i} / W\left(\overline{\mathbb{F}}_{p}\right)\right)
$$

It is clear that $f^{*}$ preserves the decomposition

$$
H_{1}^{\mathrm{dR}}\left(\tilde{B}_{i} / W\left(\overline{\mathbb{F}}_{p}\right)\right)=H_{1}^{\mathrm{dR}}\left(\tilde{B}_{i} / W\left(\overline{\mathbb{F}}_{p}\right)\right)_{1} \oplus H_{1}^{\mathrm{dR}}\left(\tilde{B}_{i} / W\left(\overline{\mathbb{F}}_{p}\right)\right)_{2}
$$

according to the two embeddings of $\mathcal{O}_{E}$ into $W\left(\overline{\mathbb{F}}_{p}\right)$. By the Kottwitz's determinant condition for $\mathcal{S} h_{0, n}$, the Hodge filtrations on $H_{1}^{\mathrm{dR}}\left(\tilde{B}_{i} / W\left(\overline{\mathbb{F}}_{p}\right)\right)$ are trivial, namely,

$$
\omega_{\tilde{B}_{i}^{V} / W\left(\overline{\mathbb{F}}_{p}\right), 1}^{\circ}=0, \quad \text { and } \quad \omega_{\tilde{B}_{i}^{V} / W\left(\overline{\mathbb{F}}_{p}\right), 2}^{\circ}=H_{1}^{\mathrm{dR}}\left(\tilde{B}_{i} / W\left(\overline{\mathbb{F}}_{p}\right)\right)_{2}^{\circ} \quad \text { for } i=1,2
$$

It is clear now $f^{*}$ preserves this trivial Hodge filtration, since it does so when tensoring with $\overline{\mathbb{F}}_{p}$.

Fix a geometric point $z=(B, \lambda, \eta) \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$. Put $C=\operatorname{End}_{\mathcal{O}_{D}}(B)_{\mathbb{Q}}$, and denote by $\dagger$ the Rosati involution on $C$ induced by $\lambda$. Let $I$ be the algebraic group over $\mathbb{Q}$ such that

$$
\begin{equation*}
I(R)=\left\{x \in C \otimes_{\mathbb{Q}} R \mid x x^{\dagger} \in R^{\times}\right\}, \quad \text { for all } \mathbb{Q} \text {-algebras } R \tag{4.9.1}
\end{equation*}
$$

Corollary 4.10. We have an isomorphism of algebraic groups over $\mathbb{Q}: I \simeq G_{0, n}$.
Proof. Let $\tilde{z}=(\tilde{B}, \tilde{\lambda}, \tilde{\eta}) \in \mathcal{S} h_{0, n}\left(\overline{\mathbb{Q}}_{p}\right)$ denote the unique lift of $z$ according to Proposition 4.9 (1). By 4.9 (2), we have a canonical isomorphism

$$
\operatorname{End}_{\mathcal{O}_{D}}(\tilde{B})_{\mathbb{Q}} \xrightarrow{\sim} \operatorname{End}_{\mathcal{O}_{D}}(B)_{\mathbb{Q}}=C
$$

In the proof of 4.9, we have seen that $C=D^{\text {opp }}$. Moreover, the Rosati involution on $C$ corresponds to the involution $b \mapsto b^{\sharp \beta_{0, n}}=\beta_{0, n} b^{*} \beta_{0, n}^{-1}$ on $D^{\text {opp }}$, where $\beta_{0, n}$ is the element in the definition of $\langle-,-\rangle_{0, n}$. It follows immediately that $I \simeq G_{0, n}$.

Let $\operatorname{Isog}(z) \subseteq \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ denote the subset of points $z^{\prime}=\left(B^{\prime}, \lambda^{\prime}, \eta^{\prime}\right)$ such that there exists an $\mathcal{O}_{D^{\prime}}$-equivariant quasi-isogeny $\phi: B^{\prime} \rightarrow B$ such that $\phi^{\vee} \circ \lambda \circ \phi=c_{0} \lambda^{\prime}$ for some $c_{0} \in \mathbb{Q}_{>0}$. We denote such a quasi-isogeny by $\phi: z^{\prime} \rightarrow z$ for simplicity.
Corollary 4.11. There exists a natural bijection of sets

$$
\Theta_{z}: \operatorname{Isog}(z) \xrightarrow{\sim} G_{0, n}(\mathbb{Q}) \backslash G_{0, n}\left(\mathbb{A}^{\infty}\right) / K
$$

Proof. First, we give the construction of $\Theta_{z}$. Put $V^{(p)}(B)=T^{(p)}(B) \otimes_{\hat{\mathbb{Z}}}(p), A^{\infty, p}$. Then $\eta$ determines an isomorphism

$$
\tilde{\eta}: V_{0, n}^{(p)} \otimes_{\mathbb{Q}} A^{\infty, p} \xrightarrow{\sim} V^{(p)}(B),
$$

modulo right translation by $K^{p}$. For any $z^{\prime}=\left(B^{\prime}, \lambda^{\prime}, \eta^{\prime}\right) \in \operatorname{Isog}(z)$ and a choice of $\phi: B^{\prime} \rightarrow B$ as above. The quasi-isogeny $\phi$ induces an isomorphism $\phi_{*}: V^{(p)}\left(B^{\prime}\right) \xrightarrow{\sim}$ $V^{(p)}(B)$. Then there exists a $g^{p} \in G_{0, n}\left(\mathbb{A}^{\infty, p}\right)$, unique up to right multiplication by elements of $K^{p}$, such that the $K^{p}$-orbit of $\phi_{*}^{-1} \circ \tilde{\eta} \circ g$ gives $\eta^{\prime}$.

We put

$$
\begin{equation*}
\mathbb{L}_{z}=\tilde{\mathcal{D}}(B)_{1}^{\circ, F^{2}=p}=\left\{v \in \tilde{\mathcal{D}}(B)_{1}^{\circ}: F^{2}(v)=p v\right\} \tag{4.11.1}
\end{equation*}
$$

Since $B$ is supersingular (See Remark 3.7), this is a free $\mathbb{Z}_{p^{2}}$-module of rank $n$, and we have $\tilde{\mathcal{D}}(B)_{1}^{\circ}=\mathbb{L}_{z} \otimes_{\mathbb{Z}_{p^{2}}} W\left(\overline{\mathbb{F}}_{p}\right)$. Put $\mathbb{L}_{z}[1 / p]=\mathbb{L}_{z} \otimes_{\mathbb{Z}_{p^{2}}} \mathbb{Q}_{p^{2}}$. Then $\phi$ induces an isomorphism $\phi_{*}: \mathbb{L}_{z^{\prime}}[1 / p] \xrightarrow{\sim} \mathbb{L}_{z}[1 / p]$. Fix a $\mathbb{Z}_{p^{2}}$-basis for $\mathbb{L}_{z}$. Then there exists a $g_{\mathbb{\unrhd}} \in \mathrm{GL}_{n}\left(\mathbb{Q}_{p^{2}}\right)$ such that $\phi_{*}\left(\mathbb{L}_{z^{\prime}}\right)=g_{\mathbb{L}}\left(\mathbb{L}_{z}\right)$, and the right coset $g_{\square} \mathrm{GL}_{n}\left(\mathbb{Z}_{p^{2}}\right)$ is independent of the choice of such a basis. We put $g_{p}=\left(c_{0}, g_{\square}\right) \in$ $\mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(\mathbb{Q}_{p^{2}}\right) \simeq G_{0, n}\left(\mathbb{Q}_{p}\right)$, which is well defined up to right multiplication by elements of $K_{p}=\mathbb{Z}_{p^{2}}^{\times} \times \mathrm{GL}_{n}\left(\mathbb{Z}_{p^{2}}\right)$.

Finally, note that the quasi-isogeny $\phi^{\prime}: B^{\prime} \rightarrow B$ is well determined by $z^{\prime}$ up to left composition with an element $\gamma \in I(\mathbb{Q})=G_{0, n}(\mathbb{Q})$. If we replace $\phi$ by $\gamma \circ \phi$, then $g:=\left(g^{p}, g_{p}\right) \in G_{0, n}\left(\mathbb{A}^{\infty}\right)$ is replaced by $\gamma g=\left(\gamma g^{p}, \gamma g_{p}\right)$. Therefore, the map

$$
\Theta_{z}: \operatorname{Isog}(z) \rightarrow G_{0, n}(\mathbb{Q}) \backslash G_{0, n}\left(\mathbb{A}^{\infty}\right) / K, \quad z^{\prime} \mapsto G_{0, n}(\mathbb{Q}) g K
$$

is well defined. The fact that $\Theta_{z}$ is a bijection follows from the similar classical statement in characteristic 0 .

Remark 4.12. It follows from Proposition 4.9 and the description of $\mathcal{S} h_{0, n}(\mathbb{C})$ in Section 2.3 that $\operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ consists of $\# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{0, n}\right)$ isogeny classes of abelian varieties equipped with additional structures.

Lemma 4.13. Let $N$ be a fixed nonnegative integer. Up to replacing $K^{p}$ by an open compact subgroup of itself, the following properties are satisfied: if $(B, \lambda, \eta)$ is an $\overline{\mathbb{F}}_{p}$-point of $\mathrm{Sh}_{0, n}$ and $f: B \rightarrow B$ is an $\mathcal{O}_{D}$-quasi-isogeny such that $p^{N} f \in$ $\operatorname{End}_{\mathcal{O}_{D}}(B), f^{\vee} \circ \lambda \circ f=\lambda$ and $f \circ \eta=\eta$, then $f=\mathrm{id}$.

Proof. It suffices to prove the lemma for $(B, \lambda, \eta)$ in a fixed isogeny class $\operatorname{Isog}(z)$ of $\operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$. We write $G_{0, n}\left(\mathbb{A}^{\infty}\right)=\coprod_{i \in I} G_{0, n}(\mathbb{Q}) g_{i} K$ with $K=K^{p} K_{p}$, where $g_{i}=g_{i}^{p} g_{i, p}$, with $g_{i}^{p} \in G_{0, n}\left(\mathbb{A}^{\infty, p}\right)$ and $g_{i, p} \in G_{0, n}\left(\mathbb{Q}_{p}\right)$, runs through a finite set of representatives of the double coset

$$
G_{0, n}(\mathbb{Q}) \backslash G_{0, n}\left(\mathbb{A}^{\infty}\right) / K
$$

Let $(B, \lambda, \eta)$ be a point of $\mathrm{Sh}_{0, n}$ corresponding to $G_{0, n}(\mathbb{Q}) g_{i} K$ for some $i \in I$, and $f$ be an $\mathcal{O}_{D}$-quasi-isogeny of $B$ as in the statement. Then $f$ is given by an element of $G_{0, n}^{1}(\mathbb{Q})$. The condition that $f \circ \eta=\eta$ is equivalent to saying that the image of $f$ in $G_{0, n}\left(\mathbb{A}^{\infty, p}\right)$ lies in $g_{i}^{p} K^{p} g_{i}^{p,-1}$. Moreover, $p^{N} f \in \operatorname{End}_{\mathcal{O}_{D}}(B)$ implies that
the image of $f$ in $G_{0, n}\left(\mathbb{Q}_{p}\right)$ belongs to $\coprod_{\delta} g_{i, p}\left(K_{p} \delta K_{p}\right) g_{i, p}^{-1}$, where $\delta$ runs through the set

$$
\begin{aligned}
& \left\{\left(1, \operatorname{Diag}\left(p^{a_{1}}, p^{a_{2}}, \ldots, p^{a_{n}}\right)\right) \in G_{0, n}\left(\mathbb{Q}_{p}\right) \simeq \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(\mathbb{Q}_{p^{2}}\right)\right. \\
& \left.\quad 0 \geq a_{1} \geq a_{2} \geq \cdots \geq a_{n} \geq-N\right\}
\end{aligned}
$$

Write $\coprod_{\delta} K_{p} \delta K_{p}=\coprod_{j \in J} h_{j} K_{p}$ for some finite set $J$. Hence, it suffices to show that there exists an open compact subgroup $K^{p} \subseteq K^{p}$ such that for all $g_{i}$,

$$
G_{0, n}^{1}(\mathbb{Q}) \cap g_{i}\left(K^{\prime p} \cdot h_{j} K_{p}\right) g_{i}^{-1}=\{1\}
$$

if $h_{j} K_{p}=K_{p}$, and empty otherwise. Since $K$ is neat, we have

$$
G_{0, n}^{1}(\mathbb{Q}) \cap g_{i}\left(K^{\prime p} K_{p}\right) g_{i}^{-1}=\{1\} \quad \text { for any } g_{i} \text { and any } K^{\prime p} \subseteq K^{p}
$$

Note that this implies that, for each $i \in I, G_{0, n}^{1}(\mathbb{Q}) \cap g_{i}\left(K^{p} \cdot h_{j} K_{p}\right) g_{i}^{-1}$ contains at most one element (because if it contains both $x$ and $y$, then $x^{-1} y$ is contained in $\left.G_{0, n}^{1}(\mathbb{Q}) \cap g_{i} K g_{i}^{-1}=\{1\}\right)$. Let $S \subset I \times J$ be the subset consisting of $(i, j)$ such that $h_{j} K_{p} \neq K_{p}$ and $G_{0, n}^{1}(\mathbb{Q}) \cap g_{i}\left(K^{p} \cdot h_{j} K_{p}\right) g_{i}^{-1}$ indeed contains one element, say $x_{i, j}$. Then $x_{i, j} \neq 1$ for all $(i, j) \in S$. Hence, one can choose a normal open compact subgroup $K^{\prime p} \subseteq K^{p}$ so that $x_{i, j} \notin g_{i}^{p} K^{\prime p} g_{i}^{p,-1}$ for all $i$. We claim that this choice of $K^{\prime p}$ will satisfy the desired property. Indeed, if $K^{p}=\coprod_{l} b_{l} K^{\prime p}$, then the double coset $G_{0, n}(\mathbb{Q}) \backslash G_{0, n}\left(\mathbb{A}^{\infty}\right) / K^{\prime p} K_{p}$ has a set of representatives of the form $g_{i} b_{l}$. Here, by abuse of notation, we consider $b_{l}$ as an element of $K$ with $p$-component equal to 1 . Then one has, for $h_{j} K_{p} \neq K_{p}$,

$$
G_{0, n}^{1}(\mathbb{Q}) \cap g_{i} b_{l}\left(K^{\prime p} h_{j} K_{p}\right) b_{l}^{-1} g_{i}^{-1}=G_{0, n}^{1}(\mathbb{Q}) \cap g_{i}\left(K^{\prime p} h_{j} K_{p}\right) g_{i}^{-1}=\varnothing .
$$

The first equality uses the fact that $K^{\prime p}$ is normal in $K^{p}$. This finishes the proof.
We come back to the discussion on the cycles $Y_{j} \subseteq \mathrm{Sh}_{1, n-1}$ for $1 \leq j \leq n$.
Proposition 4.14. Let $\left(\mathcal{A}, \lambda, \eta, \mathcal{B}, \lambda^{\prime}, \eta^{\prime}, \phi^{\text {univ }}\right)$ denote the universal object on $Y_{j}$ for $1 \leq j \leq n$, and $\mathcal{H}_{i} \subset H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \mathrm{Sh}_{0, n}\right)$ for $i=1,2$ be the universal subbundles on $Y_{j}^{\prime} \cong Y_{j}$.
(1) The induced map $T_{Y_{j}} \rightarrow \operatorname{pr}_{j}^{*} T_{\mathrm{Sh}_{1, n-1}}$ is universally injective, and we have canonical isomorphisms

$$
\begin{aligned}
N_{Y_{j}}\left(\mathrm{Sh}_{1, n-1}\right): & =\operatorname{pr}_{j}^{*} T_{\mathrm{Sh}_{1, n-1}} / T_{Y_{j}} \\
& \cong\left(\mathcal{H}_{1} / V^{-1}\left(\mathcal{H}_{2}^{(p)}\right)\right)^{*} \otimes V^{-1}\left(\mathcal{H}_{2}^{(p)}\right) \\
& \oplus\left(F\left(\mathcal{H}_{1}^{(p)}\right) / \mathcal{H}_{2}\right) \otimes\left(H_{1}^{\mathrm{dR}}\left(\mathcal{B} / \operatorname{Sh}_{0, n}\right)_{2}^{\circ} / F\left(\mathcal{H}_{1}^{(p)}\right)\right)^{*} \\
& \cong \operatorname{Lie}_{\mathcal{A}^{\vee}, 1}^{\circ} \otimes \operatorname{Coker}\left(\phi_{*, 1}^{\text {univ }}\right) \oplus \operatorname{Lie}_{\mathcal{A}, 2}^{\circ} \otimes \operatorname{Im}\left(\phi_{*, 2}^{\text {univ }}\right)^{*} .
\end{aligned}
$$

(2) Assume that $K^{p}$ is sufficiently small so that the consequences of Lemma 4.13 hold for $N=1$. For each fixed closed point $z \in \mathrm{Sh}_{0, n}$, the map $\mathrm{pr}_{j, z}:=$ $\left.\mathrm{pr}_{j}\right|_{Y_{j, z}}: Y_{j, z} \rightarrow \mathrm{Sh}_{1, n-1}$ is a closed immersion, or equivalently, the morphism $\left(\mathrm{pr}_{j}, \mathrm{pr}_{j}^{\prime}\right): Y_{j} \rightarrow \operatorname{Sh}_{1, n-1} \times \operatorname{Spec}\left(\mathbb{F}_{p^{2}}\right) \mathrm{Sh}_{0, n}$ is a closed immersion.
(3) The union of the images of $\mathrm{pr}_{j}$ for all $1 \leq j \leq n$ is the supersingular locus of $\mathrm{Sh}_{1, n-1}$, i.e., the reduced closed subscheme of $\mathrm{Sh}_{1, n-1}$ where all the slopes of the Newton polygon of the $p$-divisible group $\mathcal{A}\left[p^{\infty}\right]$ are $1 / 2$.

Proof. (1) Let $S$ be an affine noetherian $\mathbb{F}_{p^{2}}$-scheme and let $y=\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$ be an $S$-point of $Y_{j}$. Put $\hat{S}=S \times{ }_{\operatorname{Spec}\left(\mathbb{F}_{p^{2}}\right)} \operatorname{Spec}\left(\mathbb{F}_{p^{2}}[t] / t^{2}\right)$. Then we have a natural bijection

$$
\mathscr{D e f}(y, \hat{S}) \cong \Gamma\left(S, y^{*} T_{Y_{j}}\right),
$$

where $\mathscr{D e f}(y, \hat{S})$ is the set of deformations of $y$ to $\hat{S}$. Similarly, $\mathscr{D} \operatorname{ef}\left(\operatorname{pr}_{j} \circ y, \hat{S}\right) \cong$ $\Gamma\left(S, y^{*} \operatorname{pr}_{j}^{*} T_{\mathrm{Sh}_{1, n-1}}\right)$. To prove the universal injectivity of $T_{Y_{j}} \rightarrow \operatorname{pr}_{j}^{*} T_{\mathrm{Sh}_{1, n-1}}$, it suffices to show that the natural map $\mathscr{D e f}(y, \hat{S}) \rightarrow \mathscr{D} \operatorname{ef}\left(\operatorname{pr}_{j} \circ y, \hat{S}\right)$ is injective. By crystalline deformation theory (Theorem 3.4), giving a point of $\mathscr{D e f}(y, \hat{S})$ is equivalent to giving $\mathcal{O}_{\hat{S}^{-s u b b u n d l e s}} \hat{\omega}_{A^{\vee}, i}^{\circ} \subseteq H_{1}^{\text {cris }}(A / \hat{S})_{i}^{\circ}$ over $\hat{S}$ for $i=1,2$ such that

- $\hat{\omega}_{A^{\vee}, i}^{\circ}$ lifts $\omega_{A^{\vee} / S, i}^{\circ}$;
- $\hat{\omega}_{A^{\vee}, 1}^{\circ} \subseteq \operatorname{Im}\left(\phi_{*, 1}\right) \otimes \mathbb{F}_{p^{2}}[t] / t^{2}$ and $\operatorname{Im}\left(\phi_{*, 2}\right) \otimes \mathbb{F}_{p^{2}}[t] / t^{2} \subseteq \hat{\omega}_{A^{\vee}, 2}^{\circ}$ are locally direct factors.

Hence, one sees easily that

$$
\begin{aligned}
& \mathscr{D e f}(y, \hat{S}) \cong \operatorname{Hom}_{\mathcal{O}_{S}}\left(\omega_{A^{\vee} / S, 1}^{\circ},\right. \\
&\left.\quad \operatorname{Im}\left(\phi_{*, 1}\right) / \omega_{A^{\vee} / S, 1}^{\circ}\right) \\
& \oplus \operatorname{Hom}_{\mathcal{O}_{S}}\left(\omega_{A^{\vee} / S, 2}^{\circ} / \operatorname{Im}\left(\phi_{*, 2}\right), H_{1}^{\mathrm{dR}}(A / S)_{2}^{\circ} / \omega_{A^{\vee} / S, 2}^{\circ}\right) \\
& \cong \operatorname{Lie}_{A^{\vee} / S, 1}^{\circ} \otimes\left(\operatorname{Im}\left(\phi_{*, 1}\right) / \omega_{A^{\vee} / S, 1}^{\circ}\right) \oplus\left(\omega_{A^{\vee} / S, 2}^{\circ} / \operatorname{Im}\left(\phi_{*, 2}\right)\right)^{*} \otimes \operatorname{Lie}_{A / S, 2}^{\circ}
\end{aligned}
$$

Similarly, $\mathscr{D e f}\left(\operatorname{pr}_{j} \circ y, \hat{S}\right)$ is given by the lifts of $\omega_{A^{\vee} / S, i}^{\circ}$ to $\hat{S}$ for $i=1,2$. These lifts are classified by $\operatorname{Hom}_{\mathcal{O}_{S}}\left(\omega_{A^{\vee} / S, i}^{\circ}, H_{1}^{\mathrm{dR}}(A / S)_{i}^{\circ} / \omega_{A^{\vee} / S, i}^{\circ}\right) \cong \operatorname{Lie}_{A^{\vee} / S, i}^{\circ} \otimes_{k} \operatorname{Lie}_{A / S, i}^{\circ}$. Hence, $\mathscr{D e f}\left(\operatorname{pr}_{j} \circ y, \hat{S}\right)$ is canonically isomorphic to

$$
\mathrm{Lie}_{A^{\vee} / S, 1}^{\circ} \otimes \mathcal{O}_{S} \mathrm{Lie}_{A / S, 1}^{\circ} \oplus \mathrm{Lie}_{A^{\vee} / S, 2}^{\circ} \otimes \mathcal{O}_{S} \mathrm{Lie}_{A / S, 2}^{\circ}
$$

The natural map $\mathscr{D} \operatorname{ef}(y, \hat{S}) \rightarrow \mathscr{D} \operatorname{ef}\left(\operatorname{pr}_{j} \circ y, \hat{S}\right)$ is induced by the natural maps

$$
\begin{array}{rlrl}
\operatorname{Im}\left(\phi_{*, 1}\right) / \omega_{A^{\vee} / S, 1}^{\circ} \hookrightarrow H_{1}^{\mathrm{dR}}(A / S)_{1}^{\circ} / \omega_{A^{\vee} / S, 1}^{\circ} & \cong \mathrm{Lie}_{A / S, 1}^{\circ}, \\
\left(\omega_{A^{\vee} / S, 2}^{\circ} / \operatorname{Im}\left(\phi_{*, 2}\right)\right)^{*} \hookrightarrow & \omega_{A^{\vee} / S, 2}^{\circ} \leftrightarrows & \cong \mathrm{Lie}_{A^{\vee} / S, 2}^{\circ}
\end{array}
$$

It follows that $\mathscr{D e f}(y, \hat{S}) \rightarrow \mathscr{D} \operatorname{ef}\left(\operatorname{pr}_{j} \circ y, \hat{S}\right)$ is injective. To prove the formula for $N_{Y_{j}}\left(\mathrm{Sh}_{1, n-1}\right)$, we apply the arguments above to affine open subsets of $Y_{j}$. We see
easily that

$$
\begin{aligned}
& N_{Y_{j}}\left(\mathrm{Sh}_{1, n-1}\right) \cong \operatorname{Lie}_{\mathcal{A}^{\vee}, 1}^{\circ} \otimes_{\mathcal{O}_{Y_{j}}} \operatorname{Coker}\left(\phi_{*, 1}^{\mathrm{univ}}\right) \oplus \operatorname{Lie}_{\mathcal{A}, 2}^{\circ} \otimes_{\mathcal{O}_{Y_{j}}} \operatorname{Im}\left(\phi_{*, 2}^{\mathrm{univ}}\right)^{*} \\
& \cong\left(\mathcal{H}_{1} / V^{-1}\left(\mathcal{H}_{2}^{(p)}\right)\right)^{*} \otimes V^{-1}\left(\mathcal{H}_{2}^{(p)}\right) \\
& \oplus\left(F\left(\mathcal{H}_{1}^{(p)}\right) / \mathcal{H}_{2}\right) \otimes\left(H_{1}^{\mathrm{dR}}\left(\mathcal{B} / Y_{j}\right)_{2}^{\circ} / F\left(\mathcal{H}_{1}^{(p)}\right)\right)^{*}
\end{aligned}
$$

Here, the last step uses (4.7.2) and the isomorphism

$$
\operatorname{Im}\left(\phi_{*, 2}^{\mathrm{univ}}\right) \cong H_{1}^{\mathrm{dR}}\left(\mathcal{B} / Y_{j}\right)_{2}^{\circ} / \operatorname{Ker}\left(\phi_{*, 2}^{\mathrm{univ}}\right) \cong H_{1}^{\mathrm{dR}}\left(\mathcal{B} / Y_{j}\right)_{2}^{\circ} / F\left(\mathcal{H}_{1}^{(p)}\right)
$$

(2) By statement (1), $\mathrm{pr}_{j, z}$ induces an injection of tangent spaces at each closed points of $Y_{j, z}$. To complete the proof, it suffices to prove that $\pi_{j, z}$ induces injections on the closed points. Write $z=\left(B, \lambda^{\prime}, \eta^{\prime}\right) \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$. Assume $y_{1}$ and $y_{2}$ are two closed points of $Y_{j, z}$ with $\pi_{j}\left(y_{1}\right)=\pi_{j}\left(y_{2}\right)=(A, \lambda, \eta)$. Let $\phi_{1}, \phi_{2}: B \rightarrow A$ be the isogenies given by $y_{1}$ and $y_{2}$. Then the quasi-isogeny $\phi_{1,2}=\phi_{2}^{-1} \phi_{1} \in \operatorname{End}_{\mathcal{O}_{D}}(B)_{\mathbb{Q}}$ satisfies the conditions of Lemma 4.13 for $N=1$. Hence, we get $\phi_{1,2}=\mathrm{id}_{B}$, which is equivalent to $y_{1}=y_{2}$. This proves that $\pi_{j, z}$ is injective on closed points.
(3) The proof resembles the work of Vollaard and Wedhorn [2011]. Since all the points of $\mathrm{Sh}_{0, n}\left(\mathbb{F}_{p}\right)$ are supersingular by Remark 3.7, it is clear that the image of each $\mathrm{pr}_{j}$ lies in the supersingular locus of $\mathrm{Sh}_{1, n-1}$. Suppose now we are given a supersingular point $x=(A, \lambda, \eta) \in \operatorname{Sh}_{1, n-1}\left(\mathbb{F}_{p}\right)$. We have to show that there exists $\left(B, \lambda^{\prime}, \eta^{\prime}\right) \in \mathrm{Sh}_{0, n}$ and an isogeny $\phi: B \rightarrow A$ such that $\left(A, \lambda, \eta, \lambda^{\prime}, \eta^{\prime} ; \phi\right)$ lies in $Y_{j}$ for some $1 \leq j \leq n$.

Consider

$$
\mathbb{Z}_{\mathbb{Q}}=\left(\tilde{\mathcal{D}}(A)_{1}^{\circ}[1 / p]\right)^{F^{2}=p}=\left\{a \in \tilde{\mathcal{D}}(A)_{1}^{\circ}[1 / p] \mid F^{2}(a)=p a\right\}
$$

Since $x$ is supersingular, $\mathbb{Q}_{\mathbb{Q}}$ is a $\mathbb{Q}_{p^{2}}$-vector space of dimension $n$ by the DieudonnéManin classification, and $\tilde{\mathcal{D}}(A)_{1}^{\circ}[1 / p]=\mathbb{Q}_{\mathbb{Q}} \otimes_{\mathbb{Q}_{p^{2}}} W\left(\overline{\mathbb{F}}_{p}\right)[1 / p]$. We put $\tilde{\mathcal{E}}_{1}^{\circ}=$ $\left(\mathbb{Q}_{\mathbb{Q}} \cap \tilde{\mathcal{D}}(A)_{1}^{\circ}\right) \otimes_{\mathbb{Z}_{p^{2}}} W\left(\overline{\mathbb{F}}_{p}\right)$, and $\tilde{\mathcal{E}}_{2}^{\circ}=F\left(\tilde{\mathcal{E}}_{1}^{\circ}\right) \subseteq \tilde{\mathcal{D}}(A)_{2}^{\circ}$. Thus $\tilde{\mathcal{E}}^{\circ}=\tilde{\mathcal{E}}_{1}^{\circ} \oplus \tilde{\mathcal{E}}_{2}^{\circ}$ is a Dieudonné submodule of $\tilde{\mathcal{D}}(A)^{\circ}$. We claim that $\tilde{\mathcal{E}}^{\circ}$ contains $p \tilde{\mathcal{D}}(A)^{\circ}$ as a submodule. Then applying Proposition 3.2 with $m=1$, we get an $\mathcal{O}_{D}$-abelian variety ( $B, \lambda^{\prime}, \eta^{\prime}$ ) together with an $\mathcal{O}_{D}$-isogeny $\phi: B \rightarrow A$ with $\phi^{\vee} \circ \lambda \circ \phi=p \lambda$. It is easy to see in this case that $\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$ defines a point in $Y_{j}$ with $j=\operatorname{dim}_{\overline{\mathbb{F}}_{p}}\left(\tilde{\mathcal{D}}(A)_{2}^{\circ} / \tilde{\mathcal{E}}_{2}^{\circ}\right)$.

It then suffices to prove the claim that $p \tilde{\mathcal{D}}(A)^{\circ} \subseteq \tilde{\mathcal{E}}^{\circ}$. Suppose not, then $\tilde{\mathcal{D}}(A)^{\circ} \nsubseteq(1 / p) \tilde{\mathcal{E}}^{\circ}$. Consider $M_{i}:=\tilde{\mathcal{D}}(A)_{i}^{\circ} / \tilde{\mathcal{E}}_{i}^{\circ}$ for $i=1,2$. For any integer $\alpha \geq 0$, its $p^{\alpha}$-torsion submodule is

$$
M_{i}\left[p^{\alpha}\right]=\left(\tilde{\mathcal{D}}(A)_{i}^{\circ} \cap \frac{1}{p^{\alpha}} \tilde{\mathcal{E}}_{i}^{\circ}\right) / \tilde{\mathcal{E}}_{i}^{\circ}
$$

It follows easily that

$$
M_{i}\left[p^{\alpha+1}\right] / M_{i}\left[p^{\alpha}\right] \cong\left(\frac{1}{p^{\alpha+1}} \tilde{\mathcal{E}}_{i}^{\circ} \cap\left(\tilde{\mathcal{D}}(A)_{i}^{\circ}+\frac{1}{p^{\alpha}} \tilde{\mathcal{E}}_{i}^{\circ}\right)\right) / \frac{1}{p^{\alpha}} \tilde{\mathcal{E}}_{i}^{\circ}
$$

On the other hand, the Kottwitz's signature condition implies that both $F$ and $V: \tilde{\mathcal{D}}(A)_{1}^{\circ} \rightarrow \tilde{\mathcal{D}}(A)_{2}^{\circ}$ have cokernel isomorphic to $\overline{\mathbb{F}}_{p}$, and both $F$ and $V: \tilde{\mathcal{E}}_{1}^{\circ} \rightarrow \tilde{\mathcal{E}}_{2}^{\circ}$ are isomorphism. Therefore, the two induced morphisms

$$
F \text { and } V: M_{1} \rightarrow M_{2}
$$

are injective and both have cokernel isomorphic to $\overline{\mathbb{F}}_{p}$. It follows that the induced maps on the graded pieces

$$
\begin{align*}
F \text { and } V:\left(\frac { 1 } { p ^ { \alpha + 1 } } \tilde { \mathcal { E } } _ { 1 } ^ { \circ } \cap \left(\tilde{\mathcal{D}}(A)_{1}^{\circ}\right.\right. & \left.\left.+\frac{1}{p^{\alpha}} \tilde{\mathcal{E}}_{1}^{\circ}\right)\right) / \frac{1}{p^{\alpha}} \tilde{\mathcal{E}}_{1}^{\circ} \\
& \rightarrow\left(\frac{1}{p^{\alpha+1}} \tilde{\mathcal{E}}_{2}^{\circ} \cap\left(\tilde{\mathcal{D}}(A)_{2}^{\circ}+\frac{1}{p^{\alpha}} \tilde{\mathcal{E}}_{2}^{\circ}\right)\right) / \frac{1}{p^{\alpha}} \tilde{\mathcal{E}}_{2}^{\circ} \tag{4.14.1}
\end{align*}
$$

are injective maps, and are isomorphisms for all $\alpha \geq 0$ except for exactly one $\alpha .^{17}$ The assumption $\tilde{\mathcal{D}}(A)^{\circ} \nsubseteq(1 / p) \tilde{\mathcal{E}}^{\circ}$ implies that there are at least two $\alpha \geq 0$ for which the right hand side of (4.14.1) is nonzero. So there exists $\alpha \geq 0$ such that (4.14.1) are isomorphisms of nonzero $\overline{\mathbb{F}}_{p}$-vector spaces. Multiplication by $p^{\alpha}$ gives isomorphisms:

$$
\begin{equation*}
F \text { and } V:\left(\frac{1}{p} \tilde{\mathcal{E}}_{1}^{\circ} \cap\left(p^{\alpha} \tilde{\mathcal{D}}(A)_{1}^{\circ}+\tilde{\mathcal{E}}_{1}^{\circ}\right)\right) \rightarrow\left(\frac{1}{p} \tilde{\mathcal{E}}_{2}^{\circ} \cap\left(p^{\alpha} \tilde{\mathcal{D}}(A)_{2}^{\circ}+\tilde{\mathcal{E}}_{2}^{\circ}\right)\right) \tag{4.14.2}
\end{equation*}
$$

Now, Hilbert 90 theorem implies that $\mathbb{L}^{\prime}=\left((1 / p) \tilde{\mathcal{E}}_{1}^{\circ} \cap\left(p^{\alpha} \tilde{\mathcal{D}}(A)_{1}^{\circ}+\tilde{\mathcal{E}}_{1}^{\circ}\right)\right)^{F^{2}=p}$ in fact generates the source of (4.14.2). On the other hand, it is obvious that $\mathbb{L}^{\prime} \subset \mathbb{L}_{\mathbb{Q}}$ and $\mathbb{L}^{\prime} \subseteq p^{\alpha} \tilde{\mathcal{D}}(A)_{1}^{\circ}+\tilde{\mathcal{E}}_{1}^{\circ} \subseteq \tilde{\mathcal{D}}(A)_{1}^{\circ}$. This means that $\mathbb{L}^{\prime}$, and hence $\mathbb{L}_{\mathbb{Q}} \cap \tilde{\mathcal{D}}(A)_{1}^{\circ}$, generates the entire $(1 / p) \tilde{\mathcal{E}}_{1}^{\circ} \cap\left(p^{\alpha} \tilde{\mathcal{D}}(A)_{1}^{\circ}+\tilde{\mathcal{E}}_{1}^{\circ}\right)$, i.e., one has $(1 / p) \tilde{\mathcal{E}}_{1}^{\circ} \cap\left(p^{\alpha} \tilde{\mathcal{D}}(A)_{1}^{\circ}+\tilde{\mathcal{E}}_{1}^{\circ}\right)=\tilde{\mathcal{E}}_{1}^{\circ}$. But this contradicts with the nontriviality of the vector spaces in (4.14.1) by our choice of $\alpha$. Now the proposition is proved.

Corollary 4.15. The morphism $\mathrm{pr}_{1}$ (resp. $\mathrm{pr}_{n}$ ) is a closed immersion, and it identifies $Y_{1}$ (resp. $Y_{n}$ ) with the closed subscheme of $\mathrm{Sh}_{1, n-1}$ defined by the vanishing of $V: \omega_{\mathcal{A}^{\vee}, 2}^{\circ} \rightarrow \omega_{\mathcal{A}^{\vee}, 1}^{\circ,(p)}\left(\right.$ resp. $\left.V: \omega_{\mathcal{A}^{\vee}, 1}^{\circ} \rightarrow \omega_{\mathcal{A}^{\vee}, 2}^{\circ,(p)}\right)$.

Proof. We just prove the statement for $\mathrm{pr}_{1}$, and the case of $\mathrm{pr}_{n}$ is similar. Let $Z_{1}$ be the closed subscheme of $\mathrm{Sh}_{1, n-1}$ defined by the condition that $V: \omega_{\mathcal{A}^{\vee}, 2}^{\circ} \rightarrow \omega_{\mathcal{A}^{v}, 1}^{\mathrm{o}(p)}$ vanishes. We show first that $\mathrm{pr}_{1}: Y_{1} \rightarrow \mathrm{Sh}_{1, n-1}$ factors through the natural inclusion $Z_{1} \hookrightarrow \operatorname{Sh}_{1, n-1}$. Let $y=\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$ be an $S$-valued point of $Y_{1}$. By Lemma 4.6, $\operatorname{Im}\left(\phi_{2, *}\right)$ has rank $n-1$ and contains both $\omega_{A^{\vee} / S, 2}^{\circ}$ and $F\left(H_{1}^{\mathrm{dR}}(A / S)_{1}^{\mathrm{o},(p)}\right)$, which are both $\mathcal{O}_{S}$-subbundles of rank $n-1$. This forces $\omega_{A^{\vee} / S, 2}^{\circ}=F\left(H_{1}^{\mathrm{dR}}(A / S)_{1}^{\circ,(p)}\right)$, and therefore $V: \omega_{A^{\vee} / S, 2}^{\circ} \rightarrow \omega_{A^{\vee} / S, 1}^{\circ,(p)}$ vanishes. This shows that $\operatorname{pr}_{1}(y) \in Z_{1}$.

[^12]To prove that $\mathrm{pr}_{1}: Y_{1} \rightarrow Z_{1}$ is an isomorphism, as $Y_{1}$ is smooth, it suffices to show that it induces a bijection between closed points and tangent spaces of $Y_{1}$ and $Z_{1}$. For any perfect field $k$ containing $\mathbb{F}_{p^{2}}$, one constructs a map $\theta: Z_{1}(k) \rightarrow Y_{1}(k)$ inverse to $\mathrm{pr}_{1}: Y_{1}(k) \rightarrow Z_{1}(k)$ as follows. Given $x=(A, \lambda, \eta) \in Z_{1}(k)$. Let $\tilde{\mathcal{E}}_{1}^{\circ}=\tilde{\mathcal{D}}(A)_{1}^{\circ}$ and $\tilde{\mathcal{E}}_{2}^{\circ} \subseteq \tilde{\mathcal{D}}(A)_{2}^{\circ}$ be the inverse image of $\omega_{A^{\vee} / k, 2}^{\circ} \subseteq \tilde{\mathcal{D}}(A)_{2}^{\circ} / p \tilde{\mathcal{D}}(A)_{2}^{\circ}$. Then the condition that $y \in Z_{1}$ implies that $\tilde{\mathcal{E}}_{1}^{\circ} \oplus \tilde{\mathcal{E}}_{2}^{\circ}$ is stable under $F$ and $V$. Applying Proposition 3.2 with $m=1$, we get a tuple $\left(B, \lambda^{\prime}, \eta^{\prime}, \phi\right)$ such that $y=$ $\left(A, \lambda, \eta, B, \lambda^{\prime}, \eta^{\prime}, \phi\right) \in Y_{1}(k)$. It is immediate to check that $x \mapsto y$ and $\mathrm{pr}_{1}$ are the set theoretic inverse of each other. It remains to show that $\mathrm{pr}_{1}$ induces a bijection between $T_{Y_{1}, y}$ and $T_{Z_{1}, x}$. Proposition 4.14 already implies that we have an inclusion $T_{Y_{1}, y} \hookrightarrow T_{Z_{1}, x} \hookrightarrow T_{\mathrm{Sh}_{1, n-1}, x}$. It suffices to check that $\operatorname{dim} T_{Z_{1}, x}=n-1$. The tangent space $T_{Z_{1}, x}$ is the space of deformations $(\hat{A}, \hat{\lambda}, \hat{\eta})$ over $\hat{k}=k[\epsilon] /\left(\epsilon^{2}\right)$ of $(A, \lambda, \eta)$ such that $V: \omega_{\hat{A}^{\vee} / \hat{k}, 2}^{\circ} \rightarrow \omega_{\hat{A}^{\vee} / \hat{k}, 1}^{\circ,(p)}=\omega_{A^{\vee} / k, 1}^{\circ,(p)} \otimes_{k} \hat{k}$ vanishes. This uniquely determines the lift $\hat{\omega}_{A^{\vee}, 2}^{\circ}=\omega_{\hat{A}}{ }^{\vee} / \hat{k}, 2$. So by deformation theory (Theorem 3.4), the tangent space $T_{Z_{1}, x}$ is determined by the liftings $\hat{\omega}_{A^{\vee}, 1}^{\circ}=\omega_{A^{\vee} / \hat{k}, 1}^{\circ}$ of $\omega_{A^{\vee} / k, 1}^{\circ}$. So it is of dimension $n-1$. This concludes the proof of the corollary.
4.16. Geometric Jacquet-Langlands morphism. Let $\ell \neq p$ be a prime number. For $1 \leq j \leq n$, the diagram (4.2.1) gives rise to a natural morphism

$$
\begin{equation*}
\mathcal{J L}_{j}: H_{\mathrm{et}}^{0}\left(\overline{\operatorname{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right) \xrightarrow{\operatorname{pr}_{j}^{\prime *}} H_{\mathrm{et}}^{0}\left(\bar{Y}_{j}, \overline{\mathbb{Q}}_{\ell}\right) \xrightarrow{\mathrm{pr}_{j, 1}} H_{\mathrm{et}}^{2(n-1)}\left(\overline{\operatorname{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right) \tag{4.16.1}
\end{equation*}
$$

where $\mathrm{pr}_{j,!}$ is (2.10.1), whose restriction to each $H_{\mathrm{et}}^{0}\left(Y_{j, z}, \overline{\mathbb{Q}}_{\ell}\right)$ for $z \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ is the Gysin map associated to the closed immersion $Y_{j, z} \hookrightarrow \overline{\mathrm{Sh}}_{1, n-1}$. It is clear that the image of $\mathcal{J L}_{j}$ is the subspace generated by the cycle classes of $\left[Y_{j, z}\right] \in A^{n-1}\left(\overline{\operatorname{Sh}}_{1, n-1}\right)$ with $z \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$. According to [Helm 2010], $\mathcal{J L}_{j}$ should be considered as a certain geometric realization of the Jacquet-Langlands transfer from $G_{0, n}$ to $G_{1, n-1}$. Putting all the $\mathcal{J L}_{j}$ together, we get a morphism

$$
\begin{equation*}
\mathcal{J L}=\sum_{j} \mathcal{J L}_{j}: \bigoplus_{j=1}^{n} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right) \rightarrow H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right) \tag{4.16.2}
\end{equation*}
$$

Recall that we have fixed an isomorphism $G_{1, n-1}\left(\mathbb{A}^{\infty}\right) \simeq G_{0, n}\left(\mathbb{A}^{\infty}\right)$, which we write uniformly as $G\left(\mathbb{A}^{\infty}\right)$. Denote by $\mathscr{H}\left(K^{p}, \overline{\mathbb{Q}}_{\ell}\right)=\overline{\mathbb{Q}}_{\ell}\left[K^{p} \backslash G\left(\mathbb{A}^{\infty}, p\right) / K^{p}\right]$ the prime-to- $p$ Hecke algebra. Then the homomorphism (4.16.2) is a homomorphism of $\mathscr{H}\left(K^{p}, \overline{\mathbb{Q}}_{\ell}\right)$-modules.

For an irreducible admissible representation $\pi$ of $G\left(\mathbb{A}_{\infty}\right)$, we write $\pi=\pi^{p} \otimes \pi_{p}$, where $\pi^{p}$ (resp. $\pi_{p}$ ) is the prime-to- $p$ part (resp. the $p$-component) of $\pi$.

Lemma 4.17. Let $\pi_{1}$ and $\pi_{2}$ be two admissible irreducible representations of $G\left(\mathbb{A}^{\infty}\right)$, and $\left(r_{i}, s_{i}\right)$ for $i=1,2$ be two pairs of integers with $0 \leq r_{i}, s_{i} \leq n$ and $r_{1}+s_{1} \equiv r_{2}+s_{2} \bmod 2$. Assume that $\pi_{1}$ satisfies Hypothesis 2.5 with $a_{0}=\left(r_{1}, s_{1}\right)$,
and there exists an admissible irreducible representation $\pi_{2, \infty}$ of $G_{\left(r_{2}, s_{2}\right)}(\mathbb{R})$ such that $\pi_{2} \otimes \pi_{2, \infty}$ is a cuspidal automorphic representation of $G_{\left(r_{2}, s_{2}\right)}(\mathbb{A})$. If $\pi_{1}^{p}$ and $\pi_{2}^{p}$ are isomorphic as representations of $G\left(\mathbb{A}^{p, \infty}\right)$, then $\pi_{1, p} \simeq \pi_{2, p}$, and $\pi_{2} \otimes \pi_{2, \infty}$ admits a base change to a cuspidal automorphic representation of $\mathrm{GL}_{n}\left(\mathbb{A}_{E}\right) \times \mathbb{A}_{E_{0}}^{\times}$; in particular, $\pi_{2}$ satisfies Hypothesis 2.5 for $a_{0}=\left(r_{2}, s_{2}\right)$.

Proof. By assumption on $\pi_{1}$, there exists an irreducible admissible representation $\pi_{1, \infty}$ of $G_{\left(r_{1}, s_{1}\right)}(\mathbb{R})$ such that $\pi_{1} \otimes \pi_{1, \infty}$ is a cuspidal automorphic representation of $G_{r_{1}, s_{1}}(\mathbb{A})$, which base changes to a cuspidal automorphic representation $\left(\Pi_{1}, \chi_{1}\right)$ of $\mathrm{GL}_{n}\left(\mathbb{A}_{E}\right) \times \mathbb{A}_{E_{0}}^{\times}$. On the other hand, by [Shin 2014, Theorem 1.1], there exists always a base change of $\pi_{2} \otimes \pi_{2, \infty}$ to an automorphic representation $\left(\Pi_{2}, \chi_{2}\right)$ of $\operatorname{GL}_{n}\left(\mathbb{A}_{E}\right) \times \mathbb{A}_{E_{0}}^{\times}$. The base changes $\left(\Pi_{i}, \chi_{i}\right)$ with $i=1,2$ satisfy the following properties:

- $\Pi_{i}$ is conjugate self-dual,
- for every unramified rational prime $x$, the $x$-component of $\left(\Pi_{i}, \psi_{i}\right)$ depends only on the $x$-component of $\pi_{i}$ and
- if $\pi_{i, p}=\pi_{i, 0} \otimes \pi_{i, \mathfrak{p}}$ as representation of $G\left(\mathbb{Q}_{p}\right) \simeq \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$, then $\Pi_{i, p}=$ $\left(\pi_{i, \mathfrak{p}} \otimes \check{\pi}_{i, \mathfrak{p}}^{c}\right)$ as a representation of $\mathrm{GL}_{n}\left(E \otimes \mathbb{Q}_{p}\right) \cong \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right) \times \mathrm{GL}_{n}\left(E_{\overline{\mathfrak{p}}}\right)$, and $\psi_{i, p}=\pi_{i, 0} \otimes \pi_{i, 0}^{-1}$ as a representation of $\left(E_{0} \otimes \mathbb{Q}_{p}\right)^{\times}=\mathbb{Q}_{p}^{\times} \times \mathbb{Q}_{p}^{\times}$. Here, $\check{\pi}_{i, \mathfrak{p}}^{c}$ denotes the complex conjugate of the contragredient of $\pi_{i, \mathfrak{p}}$.
As $\pi_{1}^{p} \simeq \pi_{2}^{p},\left(\Pi_{1}, \psi_{1}\right)$ and $\left(\Pi_{2}, \psi_{2}\right)$ are isomorphic at almost all finite places. By the strong multiplicity one theorem for $\mathrm{GL}_{n}$ [Jacquet and Shalika 1981], we have $\left(\Pi_{1}, \psi_{1}\right) \simeq\left(\Pi_{2}, \psi_{2}\right)$; in particular, $\left(\Pi_{2}, \psi_{2}\right)$ is cuspidal. By the description of ( $\Pi_{i, p}, \psi_{i, p}$ ), it follows immediately that $\pi_{1, p} \simeq \pi_{2, p}$.

Let $\mathscr{A}_{K}$ be the set of isomorphism classes of irreducible admissible representations $\pi$ of $G\left(\mathbb{A}^{\infty}\right)$ satisfying Hypothesis 2.5 with $a_{\bullet}=(0, n)$. In particular, each $\pi \in \mathscr{A}_{K}$ is the finite part of an automorphic cuspidal representation of $G_{0, n}(\mathbb{A})$.

We fix such a $\pi \in \mathscr{A}_{K}$. Let

$$
\mathcal{J}_{\pi}: \bigoplus_{i=1}^{n} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)_{\pi^{p}} \rightarrow H_{\mathrm{et}}^{2(n-1)}\left(\overline{\operatorname{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi^{p}}
$$

denote the homomorphism on the $\left(\pi^{p}\right)^{K^{p}}$-isotypic components induced by $\mathcal{J} \mathcal{L}$, where for an $\mathscr{H}\left(K^{p}, \overline{\mathbb{Q}}_{\ell}\right)$-module $M$ we put

$$
M_{\pi^{p}}:=\operatorname{Hom}_{\left.\mathscr{H}_{\left(K^{p},\right.}, \overline{\mathbb{Q}}_{\ell}\right)}\left(\left(\pi^{p}\right)^{K^{p}}, M\right) \otimes\left(\pi^{p}\right)^{K^{p}} .
$$

Then Lemma 4.17 implies that $\pi$ is completely determined by its prime-to- $p$ part. Hence, taking the $\pi^{p}$-isotypic components is the same as taking the $\pi$-isotypic components. We can thus write $M_{\pi}$ instead of $M_{\pi^{p}}$ for a $\mathscr{H}\left(K, \bar{Q}_{\ell}\right)$-module $M$.

Recall that the image of $\mathcal{J}_{\pi}$ is included in $H_{\mathrm{et}}^{2(n-1)}\left(\overline{\operatorname{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{fin}}$, which is the maximal subspace of $H_{\mathrm{et}}^{2(n-1)}\left(\overline{\operatorname{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}$ where the action of $\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{2}}\right)$ factors through a finite quotient. Note that, at this moment, it is not clear if the target of $\mathcal{J L}_{\pi}$ is nonzero. But this will follow from the proof of our main Theorem 4.18 below.

Our main result claims that this inclusion is actually an equality under certain genericity conditions on $\pi_{p}$. To make this precise, write $\pi_{p}=\pi_{p, 0} \otimes \pi_{\mathfrak{p}}$ as a representation of $G\left(\mathbb{Q}_{p}\right) \simeq \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$. Let

$$
\rho_{\pi_{\mathfrak{p}}}: W_{\mathbb{Q}_{p^{2}}} \rightarrow \mathrm{GL}_{n}\left(\overline{\mathbb{Q}}_{\ell}\right)
$$

be the unramified representation of the Weil group of $\mathbb{Q}_{p^{2}}$ defined in (2.6.1). It induces a continuous $\ell$-adic representation of $\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{2}}\right)$, which we denote by the same notation. Then $\rho_{\pi_{\mathfrak{p}}}\left(\operatorname{Frob}_{p^{2}}\right)$ is semisimple with characteristic polynomial (2.6.2). Using this, we get an explicit description of $H_{\mathrm{et}}^{2(n-1)}\left(\overline{\operatorname{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}$ and $H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)_{\pi}$ in terms of $\rho_{\pi_{\mathfrak{p}}}$ by (2.4.1) and (2.6.3).

We can now state our main theorem.
Theorem 4.18. Fix $a \pi$ in $\mathscr{A}_{K}$. Let $\alpha_{\pi_{\mathfrak{p}}, 1}, \ldots, \alpha_{\pi_{\mathfrak{p}, n}}$ be the eigenvalues of $\rho_{\pi_{\mathfrak{p}}}\left(\operatorname{Frob}_{p^{2}}\right)$.
(1) If $\alpha_{\pi_{\mathfrak{p}}, 1}, \ldots, \alpha_{\pi_{\mathfrak{p}}, n}$ are distinct, then the map $\mathcal{J}_{\pi}$ is injective;
(2) Let $m_{1, n-1}(\pi)$ (resp. $m_{0, n}(\pi)$ ) denote the multiplicity for $\pi$ appearing in Theorem 2.6 for $a_{\bullet}=(1, n-1)\left(\right.$ resp. for $\left.a_{0}=(0, n)\right)$. Assume that $m_{1, n-1}(\pi)=$ $m_{0, n}(\pi)$ and that $\alpha_{\pi_{\mathfrak{p}}, i} / \alpha_{\pi_{\mathfrak{p}}, j}$ is not a root of unity for all $1 \leq i, j \leq n$. Then the map

$$
\mathcal{J} \mathcal{L}_{\pi}: \bigoplus_{j=1}^{n} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)_{\pi} \rightarrow H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{fin}}
$$

is an isomorphism. In other words, $H_{\mathrm{et}}^{2(n-1)}\left(\overline{\operatorname{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{fin}}$ is generated by the cycle classes of the irreducible components of $Y_{j}$ for $1 \leq j \leq n$.
The proof of this theorem will be given at the end of Section 6.
Remark 4.19. The equality $m_{1, n-1}(\pi)=m_{0, n}(\pi)$ is conjectured to be true according to Arthur's formula on the automorphic multiplicities of unitary groups, and is known to hold when $\pi$ is the finite part of an automorphic representation of $G_{1, n-1}(\mathbb{A})$ whose base change to $\mathrm{GL}_{n}\left(\mathbb{A}_{E}\right) \times \mathbb{A}_{E_{0}}^{\times}$is cuspidal, and $G_{1, n-1}$ is quasisplit at all finite places. See for instance [White 2012, Theorem E].

On the other hand, Theorem 4.18(1) gives partial results towards the equality $m_{1, n-1}(\pi)=m_{0, n}(\pi)$. Indeed, when combining with Kottwitz's description 2.6 of the $\pi$-isotypic components of the cohomology groups, Theorem 4.18(1) implies (under the assumption that the Satake parameters of $\pi_{\mathfrak{p}}$ are distinct) that $m_{1, n-1}(\pi) \geq$ $m_{0, n}(\pi)$ without using Arthur's trace formula. If we use only the fact that $\mathcal{J}_{\pi}$
is nonzero (which is an easy consequence of our computation of the intersection matrix in Theorem 6.7), we get the implication $m_{0, n}(\pi) \neq 0 \Rightarrow m_{1, n-1}(\pi) \neq 0$.

## 5. Fundamental intersection numbers

In this section, we will compute some intersection numbers on certain DeligneLusztig varieties. These numbers will play a key role in the computation in the next section of the intersection matrix of the cycles $Y_{j}$ on $\mathrm{Sh}_{1, n-1}$.
Notation 5.1. Let $X$ be an algebraic variety of pure dimension $N$ over $\overline{\mathbb{F}}_{p}$. For an integer $r \geq 0$, let $A^{r}(X)$ (resp. $A_{r}(X)$ ) denote the group of algebraic cycles on $X$ of codimension $r$ (resp. of dimension $r$ ) modulo rational equivalences. If $Y \subseteq X$ is a subscheme equidimensional of codimension $r$, we denote by $[Y] \in A^{r}(X)$ the class of $Y$. We put $A^{\star}(X)=\bigoplus_{r=0}^{N} A^{r}(X)$. For a zero-dimensional cycle $\eta \in A^{N}(X)$, we denote by

$$
\operatorname{deg}(\eta)=\int_{X} \eta
$$

the degree of $\eta$. Let $\mathcal{V}$ be a vector bundle over $X$. We denote by $c_{r}(\mathcal{V}) \in A^{r}(X)$ the $r$-th Chern class of $\mathcal{V}$ for $0 \leq r \leq N$, and put $c(\mathcal{V})=\sum_{r=0}^{N} c_{r}(\mathcal{V}) t^{r}$ in the free variable $t$.
5.2. A special Deligne-Lusztig variety. We fix an integer $n \geq 1$. For an integer $0 \leq k \leq n$, we denote by $\boldsymbol{\operatorname { G r }}(n, k)$ the Grassmannian variety over $\mathbb{F}_{p}$ classifying $k$-dimensional subspaces of $\mathbb{F}_{p}^{\oplus n}$. Given an integer $k$ with $1 \leq k \leq n$, let $Z_{k}^{\langle n\rangle}$ be the subscheme of $\boldsymbol{G r}(n, k) \times \boldsymbol{G} \boldsymbol{r}(n, k-1)$ whose $S$-valued points are isomorphism classes of pairs $\left(L_{1}, L_{2}\right)$, where $L_{1}$ and $L_{2}$ are respectively subbundles of $\mathcal{O}_{S}^{\oplus n}$ of rank $k$ and $k-1$ satisfying $L_{2} \subseteq L_{1}^{(p)}$ and $L_{2}^{(p)} \subseteq L_{1}$ (with locally free quotients). The same arguments as in Proposition 4.4 show that $Z_{k}^{\langle n\rangle}$ is a smooth variety over $\mathbb{F}_{p}$ of dimension $n-1$. We denote the natural closed immersion by

$$
i_{k}: Z_{k}^{\langle n\rangle} \hookrightarrow \boldsymbol{G r}(n, k) \times \boldsymbol{G r}(n, k-1)
$$

Let $\mathcal{L}_{1}$ and $\mathcal{L}_{2}$ denote the universal subbundles on $\boldsymbol{\operatorname { r r }}(n, k) \times \boldsymbol{G r}(n, k-1)$ coming from the two factors, and $\mathcal{Q}_{1}$ and $\mathcal{Q}_{2}$ the universal quotients, respectively. When there is no confusion, we still use $\mathcal{L}_{i}$ and $\mathcal{Q}_{i}$ for $i=1,2$ to denote their restrictions to $Z_{k}^{\langle n\rangle}$. We put

$$
\begin{equation*}
\mathcal{E}_{k}=\left(\mathcal{L}_{1} / \mathcal{L}_{2}^{(p)}\right)^{*} \otimes \mathcal{L}_{2}^{(p)} \oplus\left(\mathcal{L}_{1}^{(p)} / \mathcal{L}_{2}\right) \otimes \mathcal{Q}_{1}^{*,(p)} \tag{5.2.1}
\end{equation*}
$$

which is a vector bundle of rank $n-1$ on $Z_{k}^{\langle n\rangle}$. (This vector bundle is modeled on the description of the normal bundle $N_{Y_{j}}\left(\mathrm{Sh}_{1, n-1}\right)$ in Proposition 4.14(1), which is how our computation will be used in the next section; see Proposition 6.4.) We have
the top Chern class $c_{n-1}\left(\mathcal{E}_{k}\right) \in A^{n-1}\left(Z_{k}^{\langle n\rangle}\right)$. We define the fundamental intersection number on $Z_{k}^{\langle n\rangle}$ as

$$
\begin{equation*}
N(n, k):=\int_{Z_{k}^{(n)}} c_{n-1}\left(\mathcal{E}_{k}\right) \tag{5.2.2}
\end{equation*}
$$

The main theorem we prove in this section is the following:
Theorem 5.3. For integers $n, r$ with $0 \leq r \leq n$, let

$$
\binom{n}{r}_{q}=\frac{\left(q^{n}-1\right)\left(q^{n-1}-1\right) \cdots\left(q^{n-r+1}-1\right)}{(q-1)\left(q^{2}-1\right) \cdots\left(q^{r}-1\right)}
$$

be the Gaussian binomial coefficients, and let d $(n, k)=(2 k-1) n-2 k(k-1)-1$ denote the dimension of $\boldsymbol{G r}(n, k) \times \boldsymbol{G r}(n, k-1)$. Then, for $1 \leq k \leq n$, we have

$$
\begin{equation*}
N(n, k)=(-1)^{n-1} \sum_{\delta=0}^{\min \{k-1, n-k\}}(n-2 \delta) p^{d(n-2 \delta, k-\delta)}\binom{n}{\delta}_{p^{2}} . \tag{5.3.1}
\end{equation*}
$$

Remark 5.4. We point out that this theorem seems to be more than a technical result. It is at the heart of the understanding of these cycles we constructed.

Proof. We first claim that $N(n, k)=N(n, n+1-k)$ for $1 \leq k \leq n$. Let $\left(L_{1}, L_{2}\right)$ be an $S$-valued point of $\boldsymbol{G r}(n, k) \times \boldsymbol{G r}(n, k-1)$, and $Q_{i}=\mathcal{O}_{S}^{\oplus n} / L_{i}$ for $i=1,2$ be the corresponding quotient bundles. Then $\left(L_{1}, L_{2}\right) \mapsto\left(Q_{2}^{*}, Q_{1}^{*}\right)$ defines a duality isomorphism

$$
\theta: \boldsymbol{G r}(n, k) \times \boldsymbol{G r}(n, k-1) \xrightarrow{\sim} \boldsymbol{G r}(n, n+1-k) \times \boldsymbol{G r}(n, n-k)
$$

Since $L_{2}^{(p)} \subseteq L_{1}$ (resp. $L_{2} \subseteq L_{1}^{(p)}$ ) is equivalent to $Q_{1}^{*} \subseteq Q_{2}^{*,(p)}$ (resp. to $Q_{1}^{*,(p)} \subseteq Q_{2}^{*}$ ), $\theta$ induces an isomorphism between $Z_{k}^{(n)}$ and $Z_{n+1-k}^{(n)}$. It is also direct to check that $\mathcal{E}_{k}=\theta^{*}\left(\mathcal{E}_{n+1-k}\right)$. This verifies the claim. Now since the right hand side of (5.3.1) is also invariant under replacing $k$ by $n+1-k$, it suffices to prove the theorem when $k \leq \frac{1}{2}(n+1)$.

We reduce the proof of the theorem to an analogous situation where the twists are given on one of the $L_{i}$. Let $\tilde{Z}_{k}^{\langle n\rangle}$ be the subscheme of $\boldsymbol{G r}(n, k) \times \boldsymbol{G r}(n, k-1)$ whose $S$-valued points are the isomorphism classes of pairs ( $\tilde{L}_{1}, \tilde{L}_{2}$ ), where $\tilde{L}_{1}$ and $\tilde{L}_{2}$ are respectively subbundles of $\mathcal{O}_{S}^{\oplus n}$ of rank $k$ and $k-1$ satisfying $\tilde{L}_{2} \subseteq \tilde{L}_{1}$ and $\tilde{L}_{2}^{\left(p^{2}\right)} \subseteq \tilde{L}_{1}$. The relative Frobenius morphisms on the two Grassmannian factors induce two morphisms

$$
\begin{gathered}
Z_{k}^{\langle n\rangle} \xrightarrow{\varphi}\left(\tilde{Z}_{k}^{\langle n\rangle} \xrightarrow{\hat{\varphi}}\left(Z_{k}^{\langle n\rangle}\right)^{(p)}\right. \\
\left(L_{1}, L_{2}\right) \longmapsto\left(L_{1}^{(p)}, L_{2}\right) \\
\left(\tilde{L}_{1}, \tilde{L}_{2}\right) \longmapsto\left(\tilde{L}_{1}, \tilde{L}_{2}^{(p)}\right),
\end{gathered}
$$

such that the composition is the relative Frobenius on $\tilde{Z}_{k}^{\langle n\rangle}$. Using a simple deformation computation, we see that $\varphi$ has degree $p^{n-k}$ and $\hat{\varphi}$ has degree $p^{k-1}$. Let $\tilde{\mathcal{L}}_{1}$ and $\tilde{\mathcal{L}}_{2}$ denote the universal subbundles on $\boldsymbol{G r}(n, k) \times \boldsymbol{G r}(n, k-1)$ when restricted to $\tilde{Z}_{k}^{\langle n\rangle}$; let $\tilde{\mathcal{Q}}_{1}$ and $\tilde{\mathcal{Q}}_{2}$ denote the universal quotients, respectively. We put

$$
\begin{equation*}
\tilde{\mathcal{E}}_{k}=\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)^{*} \otimes \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} \oplus\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*} \tag{5.4.1}
\end{equation*}
$$

which is a vector bundle of rank $n-1$ on $\tilde{Z}_{k}^{\langle n\rangle}$.
Note that

$$
\varphi^{*}\left(\tilde{\mathcal{E}}_{k}\right)=\left(\mathcal{L}_{1}^{(p)} / \mathcal{L}_{2}^{\left(p^{2}\right)}\right)^{*} \otimes \mathcal{L}_{2}^{\left(p^{2}\right)} \oplus\left(\mathcal{L}_{1}^{(p)} / \mathcal{L}_{2}\right) \otimes \mathcal{Q}_{1}^{*,(p)}
$$

Comparing with $\mathcal{E}_{k}$, we see that $c_{n-1}\left(\varphi^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right)=p^{k-1} c_{n-1}\left(\mathcal{E}_{k}\right)$, where the factor $p^{k-1}$ comes from the Frobenius twist on the first factor. Thus, we have

$$
\begin{align*}
\int_{\tilde{Z}_{k}^{(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k}\right) & =(\operatorname{deg} \varphi)^{-1} \int_{Z_{k}^{(n)}} c_{n-1}\left(\varphi^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right) \\
& =p^{k-n} \int_{Z_{k}^{(n)}} p^{k-1} c_{n-1}\left(\mathcal{E}_{k}\right)=p^{2 k-n-1} N(n, k) \tag{5.4.2}
\end{align*}
$$

Since $d(n-2 \delta, k-\delta)+2 k-n-1=2(k-\delta-1)(n-k-\delta+1)$, the theorem is in fact equivalent to the following (for each fixed $k$ ).

Proposition 5.5. For $1 \leq k \leq(n+1) / 2$, we have

$$
\begin{equation*}
\int_{\tilde{Z}_{k}^{(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k}\right)=(-1)^{n-1} \sum_{\delta=0}^{k-1}(n-2 \delta) p^{2(k-\delta-1)(n-k-\delta+1)}\binom{n}{\delta}_{p^{2}} \tag{5.5.1}
\end{equation*}
$$

Remark 5.6. Before giving the proof of this proposition, we point out a variant of the construction of $\tilde{Z}_{k}^{\langle n\rangle}$. Let $\tilde{Z}_{k}^{\prime(n)}$ be the subscheme of $\boldsymbol{\operatorname { G r }}(n, k) \times \boldsymbol{\operatorname { G r }}(n, k-1)$ whose $S$-valued points are the isomorphism classes of pairs $\left(\tilde{L}_{1}^{\prime}, \tilde{L}_{2}^{\prime}\right)$, where $\tilde{L}_{1}^{\prime}$ and $\tilde{L}_{2}^{\prime}$ are respectively subbundles of $\mathcal{O}_{S}^{\oplus n}$ of rank $k$ and $k-1$ satisfying $\tilde{L}_{2}^{\prime} \subseteq \tilde{L}_{1}^{\prime}$ and $\tilde{L}_{2}^{\prime} \subseteq \tilde{L}_{1}^{\prime\left(p^{2}\right)}$ (Note that the twist is on $L_{1}^{\prime}$ as opposed to be on $L_{2}^{\prime}$ ). This is again a certain partial-Frobenius twist of $Z_{k}^{\langle n\rangle}$; it is smooth of dimension $n-1$. Define the universal subbundles and quotient bundles $\tilde{\mathcal{L}}_{1}^{\prime}, \tilde{\mathcal{L}}_{2}^{\prime}, \tilde{\mathcal{Q}}_{1}^{\prime}$, and $\tilde{\mathcal{Q}}_{2}^{\prime}$ similarly. We put

$$
\tilde{\mathcal{E}}_{k}^{\prime}=\left(\tilde{\mathcal{L}}_{1}^{\prime} / \tilde{\mathcal{L}}_{2}^{\prime}\right)^{*} \otimes \tilde{\mathcal{L}}_{2}^{\prime} \oplus\left(\tilde{\mathcal{L}}_{1}^{\prime\left(p^{2}\right)} / \tilde{\mathcal{L}}_{2}^{\prime}\right) \otimes\left(\tilde{\mathcal{Q}}_{1}^{\prime *}\right)^{\left(p^{2}\right)}
$$

Using the same argument as above, we see that, for every fixed $k$,

$$
\int_{\tilde{Z}_{k}^{(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k}^{\prime}\right)=p^{n+1-2 k} N(n, k)
$$

Note that the exponent is different from (5.4.2). So Proposition 5.5 for each fixed $k$ is equivalent to

$$
\int_{\tilde{Z}_{k}^{\prime(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k}^{\prime}\right)=(-1)^{n-1} \sum_{\delta=0}^{k-1}(n-2 \delta) p^{2(k-\delta)(n-k-\delta)}\binom{n}{\delta}_{p^{2}},
$$

as $2(k-\delta)(n-k-\delta)=d(n-2 \delta, k-\delta)+n-2 k+1$.
Proof of Proposition 5.5. We first prove it in the case of $k=1,2$ and then we explain an inductive process to deal with the general case.

When $k=1, \tilde{Z}_{1}^{\langle n\rangle}$ classifies a line subbundle $\tilde{L}_{1}$ of $\mathcal{O}_{S}^{\oplus n}$ with no additional condition (as $\tilde{L}_{2}$ is zero); so $\tilde{Z}_{1}^{\langle n\rangle} \cong \mathbb{P}^{n-1}$ and $\tilde{\mathcal{L}}_{1}=\mathcal{O}_{\mathbb{P}^{n-1}}(-1)$. The vector bundle $\tilde{\mathcal{E}}_{1}$ is equal to $\tilde{\mathcal{L}}_{1} \otimes \tilde{\mathcal{Q}}_{1}^{*}$. It is straightforward to check that

$$
c\left(\tilde{\mathcal{E}}_{1}\right)=\left(1+c_{1}\left(\mathcal{O}_{\mathbb{P}^{n-1}}(-1)\right)\right)^{n} \quad \text { and hence } \quad \int_{\tilde{Z}_{1}^{(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{1}\right)=(-1)^{n-1} n
$$

the proposition is proved in this case.
When $k=2$, we consider a forgetful morphism

$$
\psi: \quad \tilde{Z}_{2}^{\langle n\rangle} \rightarrow \tilde{Z}_{1}^{\langle n\rangle}, \quad\left(\tilde{L}_{1}, \tilde{L}_{2}\right) \mapsto \tilde{L}_{2}
$$

This morphism is an isomorphism over the closed points $x \in \tilde{Z}_{1}^{\langle n\rangle}\left(\overline{\mathbb{F}}_{p}\right)$ for which $\tilde{L}_{2, x} \neq \tilde{L}_{2, x}^{\left(p^{2}\right)}$, because in this case $\tilde{L}_{1, x}$ is forced to be $\tilde{L}_{2, x}+\tilde{L}_{2, x}^{\left(p^{2}\right)}$. On the other hand, for a closed point $x \in \tilde{Z}_{1}^{\langle n\rangle}\left(\overline{\mathbb{F}}_{p}\right)$ where $\tilde{L}_{2, x}=\tilde{L}_{2, x}^{\left(p^{2}\right)}$, i.e., for $x \in$ $\tilde{Z}_{1}^{\langle n\rangle}\left(\mathbb{F}_{p^{2}}\right) \cong \mathbb{P}^{n-1}\left(\mathbb{F}_{p^{2}}\right), \psi^{-1}(x)$ is the space classifying a line $\tilde{L}_{1}$ in $\mathbb{F}_{p}^{\oplus n} / \tilde{L}_{2, x}$; so $\psi^{-1}(x) \simeq \mathbb{P}^{n-2}$. A simple tangent space computation shows that $\psi$ is the blowup morphism of $\tilde{Z}_{1}^{\langle n\rangle} \cong \mathbb{P}^{n-1}$ at all of its $\mathbb{F}_{p^{2}}$-points. We use $E$ to denote the exceptional divisors, which is a disjoint union of $\binom{n}{1}_{p^{2}}$ copies of $\mathbb{P}^{n-2}$.

Note that the vanishing of the morphism $\tilde{\mathcal{L}}_{2} \rightarrow \tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}$ defines the divisor $E$ (as we can see using deformation); so

$$
\mathcal{O}_{\tilde{Z}_{2}^{(n)}}(E) \cong \tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} \otimes \tilde{\mathcal{L}}_{2}^{-1}
$$

Put $\eta=c_{1}\left(\tilde{\mathcal{L}}_{2}\right)=\psi^{*} c_{1}\left(\mathcal{O}_{\mathbb{P}^{n-1}}(-1)\right)$ and $\xi=c_{1}(E)$. Then

$$
\begin{align*}
c\left(\tilde{\mathcal{E}}_{2}\right) & =c\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)^{*} \otimes \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right) \cdot c\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*}\right) \\
& =\left(1-\xi+\left(p^{2}-1\right) \eta\right) \cdot\left(1+\xi+p^{2} \eta\right)^{n} /\left(1+\xi+\left(p^{2}-1\right) \eta\right) \tag{5.6.1}
\end{align*}
$$

where the computation of the second term comes from the following two exact sequences

$$
\begin{aligned}
& 0 \rightarrow\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*} \rightarrow\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right)^{\oplus n} \rightarrow\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{L}}_{1}^{*} \rightarrow 0 \\
& 0 \rightarrow \mathcal{O}_{\tilde{\mathcal{Z}}_{2}^{(n)}} \rightarrow\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{L}}_{1}^{*} \rightarrow\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{L}}_{2}^{*} \rightarrow 0
\end{aligned}
$$

Note that $\int_{\tilde{Z}_{2}^{(n)}} \xi^{i} \eta^{j}=0$ unless $(i, j)=(n-1,0)$ or $(0, n-1)$, in which case

$$
\int_{\tilde{Z}_{2}^{(n)}} \eta^{n-1}=(-1)^{n-1} \quad \text { and } \quad \int_{\tilde{Z}_{2}^{(n)}} \xi^{n-1}=(-1)^{n}\binom{n}{1}_{p^{2}} .
$$

Here, to prove the last formula, we used the fact that the restriction of $\mathcal{O}_{\tilde{Z}_{2}^{(n)}}(E)$ to each irreducible component $\mathbb{P}^{n-2}$ of $E$ is isomorphic to $\mathcal{O}_{\mathbb{P}^{n-2}}(-1)$. So it suffices to compute

- the $\xi^{n-1}$-coefficient of (5.6.1), which is the same as the $\xi^{n-1}$-coefficient of $(1-\xi)(1+\xi)^{n-1}$ and is equal to $2-n$; and
- the $\eta^{n-1}$-coefficient of (5.6.1), which is the same as the $\eta^{n-1}$-coefficient of $\left(1+\left(p^{2}-1\right) \eta\right)\left(1+p^{2} \eta\right)^{n} /\left(1+\left(p^{2}-1\right) \eta\right)=\left(1+p^{2} \eta\right)^{n}$ and is equal to $n p^{2(n-1)}$.
To sum up, we see that

$$
\int_{\tilde{Z}_{2}^{(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{2}\right)=(-1)^{n-1} n p^{2(n-1)}+(-1)^{n}(2-n)\binom{n}{1}_{p^{2}},
$$

which is exactly (5.5.1) for $k=2$.
In general, we make an induction on $k$. Assume that the proposition is proved for $k-1 \geq 1$ and we now prove the proposition for $k$ (assuming that $k \leq \frac{1}{2}(n+1)$ ). By Remark 5.6, we get the similar intersection formula for $\tilde{\mathcal{E}}_{k-1}^{\prime}$ on $\tilde{Z}_{k-1}^{\overline{\langle n})^{2}}$ :

$$
\begin{equation*}
\int_{\tilde{Z}_{k-1}^{\prime(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k-1}^{\prime}\right)=(-1)^{n-1} \sum_{\delta=0}^{k-2}(n-2 \delta) p^{2(k-\delta-1)(n-k-\delta+1)}\binom{n}{\delta}_{p^{2}} \tag{5.6.2}
\end{equation*}
$$

We consider the moduli space $W$ over $\mathbb{F}_{p^{2}}$ whose $S$-points are tuples $\left(\tilde{L}_{1}, \tilde{L}_{2}=\right.$ $\tilde{L}_{2}^{\prime}, \tilde{L}_{3}^{\prime}$ ), where $\tilde{L}_{1}, \tilde{L}_{2}$ and $\tilde{L}_{3}^{\prime}$ are respectively subbundles of $\mathcal{O}_{S}^{\oplus n}$ of rank $k, k-1$ and $k-2$ satisfying $\tilde{L}_{3}^{\prime} \subset \tilde{L}_{2} \subset \tilde{L}_{1}$ and $\tilde{L}_{3}^{\prime} \subset \tilde{L}_{2}^{\left(p^{2}\right)} \subset \tilde{L}_{1}$. It is easy to use deformation theory to check that $W$ is a smooth variety of dimension $n-1$. There are two natural morphisms


Let $E$ denote the subspace of $W$ whose closed points $x \in W\left(\overline{\mathbb{F}}_{p}\right)$ are those such that $\tilde{L}_{2, x}=\tilde{L}_{2, x}^{\left(p^{2}\right)}$, i.e., $\tilde{L}_{2, x}$ is an $\mathbb{F}_{p^{2}}$-rational subspace of $\mathbb{F}_{p^{2}}^{\oplus n}$ of dimension $k-1$. It is clear that $E$ is a disjoint union of $\binom{n}{k-1}$ pres (corresponding to the choices of $\tilde{L}_{2}$ ) of $\mathbb{P}^{n-k} \times \mathbb{P}^{k-2}$ (corresponding to the choice of $\tilde{L}_{1}$ and $\tilde{L}_{3}^{\prime}$ respectively). It gives rise to a smooth divisor on $W$.

For a point $x \in(W \backslash E)\left(\bar{F}_{p}\right)$, we have $\tilde{L}_{2, x} \neq \tilde{L}_{2, x}^{\left(p^{2}\right)}$ and hence it uniquely determines both $\tilde{L}_{1, x}$ and $\tilde{L}_{3, x}^{\prime}$; so $\psi_{12}$ and $\psi_{23}$ are isomorphisms restricted to $W \backslash E$. On
the other hand, when restricted to $E, \psi_{12}$ contracts each copy of $\mathbb{P}^{n-k} \times \mathbb{P}^{k-2}$ of $E$ into the first factor $\mathbb{P}^{n-k}$; whereas $\psi_{23}$ contracts each copy of $\mathbb{P}^{n-k} \times \mathbb{P}^{k-2}$ of $E$ into the second factor $\mathbb{P}^{k-2}$. It is clear from this (with a little bit of help from a deformation argument) that $\psi_{12}$ is the blowup of $\tilde{Z}_{k}^{\langle n\rangle}$ along $\psi_{12}(E)$ and $\psi_{23}$ is the blowup of $\tilde{Z}_{k-1}^{\prime\langle n\rangle}$ along $\psi_{23}(E)$; the divisor $E$ is the exceptional divisor for both blowups.

A simple deformation theory argument shows that the normal bundle of $E$ in $W$ when restricted to each component $\mathbb{P}^{n-k} \times \mathbb{P}^{k-2}$ is $\mathcal{O}_{\mathbb{P}^{n-k}}(-1) \otimes \mathcal{O}_{\mathbb{P}^{k-2}}(-1)$. Moreover, we can characterize $E$ as the zero locus of either one of the following natural homomorphisms

$$
\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime} \rightarrow \tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}, \quad \tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime} \rightarrow \tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}
$$

So as a line bundle over $W$, we have

$$
\mathcal{O}_{W}(E) \cong\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \cong\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)
$$

We want to compare

$$
\begin{align*}
\int_{\tilde{Z}_{k}^{(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k}\right) & =\int_{W} c_{n-1}\left(\psi_{12}^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right) \quad \text { and }  \tag{5.6.3}\\
\int_{\tilde{Z}_{k-1}^{\prime(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k-1}^{\prime}\right) & =\int_{W} c_{n-1}\left(\psi_{23}^{*}\left(\tilde{\mathcal{E}}_{k-1}^{\prime}\right)\right)
\end{align*}
$$

We will show that they differ by $(2 k-n-2)(-1)^{n}\binom{n}{k-1}_{p^{2}}$ and this will conclude the proof of the proposition by inductive hypothesis (5.6.2). Indeed, we have

$$
\begin{align*}
& c\left(\psi_{12}^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right)=c\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)^{*} \otimes \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right) \cdot c\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*}\right)  \tag{5.6.4}\\
& c\left(\psi_{23}^{*}\left(\tilde{\mathcal{E}}_{k-1}^{\prime}\right)\right)=c\left(\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{*} \otimes \tilde{\mathcal{L}}_{3}^{\prime}\right) \cdot c\left(\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime}\right) \otimes \tilde{\mathcal{Q}}_{2}^{*,\left(p^{2}\right)}\right) \tag{5.6.5}
\end{align*}
$$

where $\tilde{\mathcal{Q}}_{1}$ and $\tilde{\mathcal{Q}}_{2}$ are the universal quotient vector bundles. Consider the following two exact sequences where the two last terms are identified:

$$
\begin{aligned}
& \mathcal{O}_{W}(E) \otimes\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)^{-1} \otimes \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} \\
& \uparrow \cong \\
& 0 \longrightarrow\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes \tilde{\mathcal{L}}_{3}^{\prime} \longrightarrow\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} \longrightarrow\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime}\right) \longrightarrow 0 \\
& \uparrow \cong \\
& 0 \longrightarrow\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*} \longrightarrow\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{2}^{*,\left(p^{2}\right)} \longrightarrow\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes\left(\tilde{\mathcal{Q}}_{2}^{*,\left(p^{2}\right)} / \tilde{\mathcal{Q}}_{1}^{*}\right) \longrightarrow 0 . \\
& \uparrow \cong \\
& \mathcal{O}_{W}(E) \otimes\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime}\right) \otimes \tilde{\mathcal{Q}}_{2}^{*,\left(p^{2}\right)}
\end{aligned}
$$

Here the right vertical isomorphism is given by

$$
\begin{aligned}
\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes\left(\tilde{\mathcal{Q}}_{2}^{*,\left(p^{2}\right)} / \tilde{\mathcal{Q}}_{1}^{*}\right) & \cong\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)^{-1} \\
& \cong\left(\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime}\right) \otimes \mathcal{O}_{W}(E)\right) \otimes\left(\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right) \otimes \mathcal{O}_{W}(E)\right)^{-1} \\
& \cong\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime}\right)
\end{aligned}
$$

From these two exact sequences we see that

$$
\begin{aligned}
& c\left(\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes \tilde{\mathcal{L}}_{3}^{\prime}\right) \cdot c\left(\mathcal{O}_{W}(E) \otimes\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime}\right) \otimes \tilde{\mathcal{Q}}_{2}^{*,\left(p^{2}\right)}\right) \\
& \quad=c\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*}\right) \cdot c\left(\mathcal{O}_{W}(E) \otimes\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)^{-1} \otimes \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)
\end{aligned}
$$

Comparing this with (5.6.5) and (5.6.4), we get

$$
\begin{aligned}
& c_{n-1}\left(\psi_{12}^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right)-c_{n-1}\left(\psi_{23}^{*}\left(\tilde{\mathcal{E}}_{k-1}^{\prime}\right)\right) \\
& =\left(c_{k-1}\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)^{-1} \otimes \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)-c_{k-1}\left(\mathcal{O}_{W}(E) \otimes\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)^{-1} \otimes \tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right)\right) \cdot c_{n-k}\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*}\right) \\
& \left.-c_{k-2}\left(\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes \tilde{\mathcal{L}}_{3}^{\prime}\right) \cdot\left(c_{n-k+1}\left(\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)} / \tilde{\mathcal{L}}_{3}^{\prime}\right) \otimes \tilde{\mathcal{Q}}_{2}^{*\left(p^{2}\right)}\right)-c_{n-k+1}\left(\mathcal{O}_{W}(E) \otimes\left(\tilde{\mathcal{L}}_{2}^{\left(p^{2}\right)}\right) \tilde{\mathcal{L}}_{3}^{\prime}\right) \otimes \tilde{\mathcal{Q}}_{2}^{*,\left(p^{2}\right)}\right)\right) .
\end{aligned}
$$

Recall that $E$ is the exceptional divisor for the blow-up $\psi_{12}$ centered at a disjoint union of $\mathbb{P}^{n-k}$; so $c_{1}(E)$ kills $\psi_{12}^{*}\left(A^{i}\left(\tilde{Z}_{k}^{\langle n\rangle}\right)\right)$ for $i \geq n-k+1$. Similarly, $c_{1}(E)$ kills $\psi_{23}^{*}\left(A^{i}\left(\tilde{Z}_{k-1}^{\prime\langle n\rangle}\right)\right)$ for $i \geq k-1$. As a result, we can rewrite the above complicated formula as

$$
\begin{aligned}
& c_{n-1}\left(\psi_{12}^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right)-c_{n-1}\left(\psi_{23}^{*}\left(\tilde{\mathcal{E}}_{k-1}^{\prime}\right)\right) \\
& \quad=-\left.c_{1}(E)^{k-2}\right|_{E} \cdot c_{n-k}\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*}\right)+\left.c_{k-2}\left(\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes \tilde{\mathcal{L}}_{3}^{\prime}\right) \cdot c_{1}(E)^{n-k}\right|_{E} \\
& \quad=\left.(-1)^{k-1} c_{n-k}\left(\left(\tilde{\mathcal{L}}_{1} / \tilde{\mathcal{L}}_{2}\right) \otimes \tilde{\mathcal{Q}}_{1}^{*}\right)\right|_{\psi_{12}(E)}+\left.(-1)^{n-k} c_{k-2}\left(\left(\tilde{\mathcal{L}}_{2} / \tilde{\mathcal{L}}_{3}^{\prime}\right)^{-1} \otimes \tilde{\mathcal{L}}_{3}^{\prime}\right)\right|_{\psi_{23}(E)}
\end{aligned}
$$

For the first term, over each $\mathbb{P}^{n-k}$ of $\psi_{12}(E)$, it is to take the top Chern class of the canonical subbundle of rank $n-k$ twisted by $\mathcal{O}_{\mathbb{P}^{n-k}}(-1)$; so the degree of the first term is $(-1)^{n-k}(n-k+1)$ on each $\mathbb{P}^{n-k}$. Similarly, for the second term, over each $\mathbb{P}^{k-2}$, it is the top Chern class of the canonical subbundle of rank $k-2$ twisted by $\mathcal{O}_{\mathbb{P}^{k-2}}(-1)$; so the degree of the second term is $(-1)^{k-2}(k-1)$ on each $\mathbb{P}^{k-2}$. To sum up, we have

$$
\begin{align*}
\int_{W} c_{n-1}( & \left(\psi_{12}^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right)-\int_{W} c_{n-1}\left(\psi_{23}^{*}\left(\tilde{\mathcal{E}}_{k-1}^{\prime}\right)\right) \\
& =(-1)^{k-1}(-1)^{n-k}(n-k+1)\binom{n}{k-1}_{p^{2}}+(-1)^{n-k}(-1)^{k-2}(k-1)\binom{n}{k-1}_{p^{2}} \\
& =(-1)^{n-1}(n-2 k+2)\binom{n}{k-1}_{p^{2}} . \tag{5.6.6}
\end{align*}
$$

So by the inductive hypothesis,

$$
\begin{aligned}
\int_{\tilde{z}_{k}^{(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k}\right) & \stackrel{(5.6 .3)}{=} \int_{W} c_{n-1}\left(\psi_{12}^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right) \\
& \stackrel{(5.6 .6)}{=} \int_{W} c_{n-1}\left(\psi_{23}^{*}\left(\tilde{\mathcal{E}}_{k}\right)\right)+(-1)^{n-1}(n-2 k+2)\binom{n}{k-1}_{p^{2}} \\
& \stackrel{(5.6 .3)}{=} \int_{\tilde{z}_{k-1}^{\prime(n)}} c_{n-1}\left(\tilde{\mathcal{E}}_{k-1}^{\prime}\right)+(-1)^{n-1}(n-2 k+2)\binom{n}{k-1}_{p^{2}} \\
& =(-1)^{n-1} \sum_{\delta=0}^{k-2}(n-2 \delta) p^{2(k-\delta-1)(n-k-\delta+1)}\binom{n}{\delta}_{p^{2}} \\
& +(-1)^{n-1}(n-2 k+2)\binom{n}{k-1}_{p^{2}} \\
& =(-1)^{n-1} \sum_{\delta=0}^{k-1}(n-2 \delta) p^{2(k-\delta-1)(n-k-\delta+1)}\binom{n}{\delta}_{p^{2}}
\end{aligned}
$$

This shows the statement of the proposition for $k$ and hence concludes the proof.

## 6. Intersection matrix of supersingular cycles on $\mathrm{Sh}_{1, n-1}$

Throughout this section, we fix an integer $n \geq 2$ and keep the notation as in Section 4. We will study the intersection theory of cycles $Y_{j}$ for $1 \leq j \leq n$ on $\mathrm{Sh}_{1, n-1}$ considered in Section 4. For this, we may assume the following:

Hypothesis 6.1. We assume that the tame level structure $K^{p}$ is taken sufficiently small so that Lemma 4.13 holds with $N=2$.
6.2. Hecke correspondences on $\mathbf{S h}_{\mathbf{0}, \boldsymbol{n}}$. Recall that we have an isomorphism

$$
G\left(\mathbb{Q}_{p}\right) \simeq \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right) \cong \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(\mathbb{Q}_{p^{2}}\right)
$$

Put $K_{\mathfrak{p}}=\mathrm{GL}_{n}\left(\mathcal{O}_{E_{\mathfrak{p}}}\right)$ and $K_{p}=\mathbb{Z}_{p}^{\times} \times K_{\mathfrak{p}}$. The Hecke algebra $\mathbb{Z}\left[K_{\mathfrak{p}} \backslash \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right) / K_{\mathfrak{p}}\right]$ can be viewed as a subalgebra of $\mathbb{Z}\left[K_{p} \backslash G\left(\mathbb{Q}_{p}\right) / K_{p}\right]$ (with trivial factor at the $\mathbb{Q}_{p}^{\times}$-component).

For $\gamma \in \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$, the double coset $T_{\mathfrak{p}}(\gamma):=K_{\mathfrak{p}} \gamma K_{\mathfrak{p}}$ defines a Hecke correspondence on $\mathrm{Sh}_{0, n}$. It induces a set theoretic Hecke correspondence

$$
T_{\mathfrak{p}}(\gamma): \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right) \rightarrow \mathcal{S}\left(\operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)\right)
$$

where $\mathcal{S}\left(\mathrm{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)\right)$ denotes the set of subsets of $\mathrm{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$. By Remark 4.12, $\mathrm{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ is a union of $\# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{0, n}\right)$-isogeny classes of abelian varieties. Fix a base point $z_{0} \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$. Let

$$
\Theta_{z_{0}}: \operatorname{Isog}\left(z_{0}\right) \xrightarrow{\sim} G_{0, n}(\mathbb{Q}) \backslash\left(G\left(\mathbb{A}^{\infty, p}\right) \times G\left(\mathbb{Q}_{p}\right)\right) / K^{p} \times K_{p} .
$$

be the bijection constructed as in Corollary 4.11. Write $K_{\mathfrak{p}} \gamma K_{\mathfrak{p}}=\coprod_{i \in I} \gamma_{i} K_{\mathfrak{p}}$. If $z \in \operatorname{Isog}\left(z_{0}\right)$ corresponds to the class of $\left(g^{p}, g_{p}\right) \in G\left(\mathbb{A}^{\infty, p}\right) \times G\left(\mathbb{Q}_{p}\right)$ with $g_{p}=\left(g_{p, 0}, g_{\mathfrak{p}}\right)$, then $T_{\mathfrak{p}}(\gamma)(z)$ consists of points in $\operatorname{Isog}\left(z_{0}\right)$ corresponding to the class of $\left(g^{p},\left(g_{p, 0}, g_{\mathfrak{p}} \gamma_{i}\right)\right)$ for all $i \in I$.

Alternatively, $T_{\mathfrak{p}}(\gamma)$ has the following description. Write $z=(A, \lambda, \eta)$, and let $\mathbb{L}_{z}$ denote the $\mathbb{Z}_{p^{2}}$-free module $\tilde{\mathcal{D}}(A)_{1}^{\circ, F^{2}=p}$. Then a point $z^{\prime}=\left(B, \lambda^{\prime}, \eta^{\prime}\right) \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ belongs to $T_{\mathfrak{p}}(\gamma)(z)$ if and only if there exists an $\mathcal{O}_{D}$-equivariant $p$-quasi-isogeny $\phi: B^{\prime} \rightarrow B$ (i.e., $p^{m} \phi$ is an isogeny of $p$-power order for some integer $m$ ) such that
(1) $\phi^{\vee} \circ \lambda \circ \phi=\lambda^{\prime}$,
(2) $\phi \circ \eta^{\prime}=\eta$,
(3) $\phi_{*}\left(\mathbb{Z}_{z^{\prime}}\right)$ is a lattice of $\mathbb{L}_{z}[1 / p]=\mathbb{Z}_{z} \otimes_{\mathbb{Z}_{p^{2}}} \mathbb{Q}_{p^{2}}$ with the property: there exists a $\mathbb{Z}_{p^{2}}$-basis $\left(e_{1}, \ldots, e_{n}\right)$ for $\mathbb{L}_{z}$ such that $\left(e_{1}, \ldots, e_{n}\right) \gamma$ is a $\mathbb{Z}_{p^{2}}$-basis for $\phi_{*}\left(\mathbb{L}_{z^{\prime}}\right)$.
When $\gamma=\operatorname{Diag}\left(p^{a_{1}}, \ldots, p^{a_{n}}\right)$ with $a_{i} \in\{-1,0,1\}$, For given $z$ and $z^{\prime}$, such a $\phi$ is necessarily unique if it exists, by Lemma 4.13 (with $N=2$ ). Therefore, $T_{\mathfrak{p}}(\gamma)(z)$ is in natural bijection with the set of $\mathbb{Z}_{p^{2}}$-lattices $\mathbb{L}^{\prime} \subseteq \mathbb{L}_{z}[1 / p]$ satisfying property (3) above.

For each integer $i$ with $0 \leq i \leq n$, we put

$$
T_{\mathfrak{p}}^{(i)}=T_{\mathfrak{p}}(\operatorname{Diag}(\underbrace{p, \ldots, p}_{i}, \underbrace{1, \ldots, 1}_{n-i}))
$$

By the discussion above, one has a natural bijection

$$
T_{\mathfrak{p}}^{(i)}(z) \xrightarrow{\sim}\left\{\mathbb{L}_{z^{\prime}} \subseteq \mathbb{L}_{z}[1 / p] \mid p \mathbb{L}_{z} \subseteq \mathbb{L}_{z^{\prime}} \subseteq \mathbb{L}_{z}, \operatorname{dim}_{\mathbb{F}_{p^{2}}}\left(\mathbb{L}_{z} / \mathbb{L}_{z^{\prime}}\right)=i\right\}
$$

for $z \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$. Note that $T_{\mathfrak{p}}^{(0)}=\mathrm{id}$, and we put $S_{\mathfrak{p}}:=T_{\mathfrak{p}}^{(n)}$. Then the Satake isomorphism implies $\mathbb{Z}\left[K_{\mathfrak{p}} \backslash \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right) / K_{\mathfrak{p}}\right] \cong \mathbb{Z}\left[T_{\mathfrak{p}}^{(1)}, \ldots, T_{\mathfrak{p}}^{(n-1)}, S_{\mathfrak{p}}, S_{\mathfrak{p}}^{-1}\right]$. More generally, for $0 \leq a \leq b \leq n$, we put

$$
R_{\mathfrak{p}}^{(a, b)}=T_{\mathfrak{p}}(\operatorname{Diag}(\underbrace{p^{2}, \ldots, p^{2}}_{a}, \underbrace{p, \ldots, p}_{b-a}, \underbrace{1, \ldots, 1}_{n-b}) .
$$

Note that $R_{\mathfrak{p}}^{(0, i)}=T_{\mathfrak{p}}^{(i)}$, and $R_{\mathfrak{p}}^{(a, b)} S_{\mathfrak{p}}^{-1}$ is the Hecke operator

$$
T_{\mathfrak{p}}(\operatorname{Diag}(\underbrace{p, \ldots, p}_{a}, \underbrace{1, \ldots, 1}_{b-a}, \underbrace{p^{-1}, \ldots, p^{-1}}_{n-b}))
$$

For the explicit relations between $R_{\mathfrak{p}}^{(a, b)}$ and $T_{\mathfrak{p}}^{(i)}$, see Proposition A.1.
6.3. Refined Gysin homomorphism. For an algebraic variety $X$ over $\overline{\mathbb{F}}_{p}$ of pure dimension $N$ and any integer $r \geq 0$, we write $A_{r}(X)=A^{N-r}(X)$ to denote the group of dimension $r$ (codimension $N-r$ ) cycles in $X$ modulo rational equivalence. Recall that the restriction of $\mathrm{pr}_{j}: Y_{j} \rightarrow \operatorname{Sh}_{1, n-1}$ to each $Y_{j, z}$ for $z \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ and
$1 \leq j \leq n$ is a regular closed immersion (into $\overline{\operatorname{Sh}}_{1, n-1}$ ). There is a well-defined Gysin homomorphism

$$
\begin{equation*}
\operatorname{pr}_{j}^{\prime}: A_{n-1}\left(\overline{\mathrm{Sh}}_{1, n-1}\right) \rightarrow A_{0}\left(\bar{Y}_{j}\right)=\bigoplus_{z \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)} A_{0}\left(Y_{j, z}\right) \tag{6.3.1}
\end{equation*}
$$

whose composition with the natural projection $A_{0}\left(\bar{Y}_{j}\right) \rightarrow A_{0}\left(Y_{j, z}\right)$ is the refined Gysin map $\left(\left.\mathrm{pr}_{j}\right|_{Y_{j, z}}\right)$ ! defined in [Fulton 1998, 6.2] for regular immersions. Let $X \subseteq \overline{\mathrm{Sh}}_{1, n-1}$ be a closed subvariety of dimension $n-1$. Consider the Cartesian diagram


Assume that the restriction of $g_{X}$ to each $Y_{j, z} \times \times_{\overline{S h}_{1, n-1}} X$ with $z \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ is a regular closed immersion as well. Then $\operatorname{pr}_{j}^{\prime}([X]) \in A_{0}\left(\bar{Y}_{j}\right)$ can be described as follows. Put $N_{Y_{j, 2}}\left(\overline{\operatorname{Sh}}_{1, n-1}\right):=\operatorname{pr}_{j}^{*}\left(T_{\overline{\operatorname{Sh}}_{1, n-1}}\right) / T_{Y_{j, z}}$, and we define $N_{Y_{j, z} \times_{\overline{S h}_{1, n-1}} X}(X)$ in a similar way. We define the excess vector bundle as

$$
\mathcal{E}\left(Y_{j, z}, X\right):=g_{j}^{*} N_{Y_{j, z}}\left(\overline{\operatorname{Sh}}_{1, n-1}\right) / N_{Y_{j, z} \times_{\overline{\mathrm{Sh}}_{1, n-1}} X}(X)
$$

This is a vector bundle on $Y_{j, z} \times{\overline{\text { Shen }_{1, n-1}}} X$. Let $r$ be its rank function, which is equal to the dimension of $\bar{Y}_{j} \times \overline{S h}_{1, n-1} X$ on each of its connected component. Then the excess intersection formula [Fulton 1998, 6.3] shows that

$$
\begin{equation*}
\operatorname{pr}_{j}^{!}([X])=\sum_{z \in \operatorname{Sh}_{0, n}\left(\bar{F}_{p}\right)} \int_{Y_{j, z} \times \times_{\operatorname{Sh}_{1, n-1}} X} c_{r}\left(\mathcal{E}\left(Y_{j, z}, X\right)\right) \tag{6.3.2}
\end{equation*}
$$

where $c_{r}\left(\mathcal{E}\left(Y_{j, z}, X\right)\right)$ is the top Chern class of $\mathcal{E}\left(Y_{j, z}, X\right)$ over $Y_{j, z} \times \overline{S h}_{1, n-1} X$. The integration should be understood as the sum over all connected components of $Y_{j, z} \times{ }_{\overline{S h}_{1, n-1}} X$ of the degrees of $c_{r}\left(\mathcal{E}\left(Y_{j, z}, X\right)\right)$.
Proposition 6.4. Let $i, j$ be integers with $1 \leq i \leq j \leq n$ and $z, z^{\prime} \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$.
(1) The subvarieties $Y_{i, z}$ and $Y_{j, z^{\prime}}$ of $\overline{\mathrm{Sh}}_{1, n-1}$ have nonempty intersection if and only if there exists an integer $\delta$ with $0 \leq \delta \leq \min \{n-j, i-1\}$ such that $z^{\prime} \in R_{\mathfrak{p}}^{(j-i+\delta, n-\delta)} S_{\mathfrak{p}}^{-1}(z)$, or equivalently $z \in R_{\mathfrak{p}}^{(\bar{\delta}, n+i-j-\delta)} S_{\mathfrak{p}}^{-1}\left(z^{\prime}\right)$, where $R_{\mathfrak{p}}^{(a, b)}$ and $S_{\mathfrak{p}}$ are the Hecke operators defined in Section 6.2.
(2) If the condition in (1) is satisfied for some $\delta$, then $Y_{i, z} \times{\overline{\operatorname{Sh}_{1, n-1}}} Y_{j, z^{\prime}}$ is isomorphic to the variety $\bar{Z}_{i-\delta}^{\langle n+i-j-2 \delta\rangle}$ defined in Section 5.2. Moreover, the excess vector bundles $\mathcal{E}\left(Y_{i, z}, Y_{j, z^{\prime}}\right)$ and $\mathcal{E}\left(Y_{j, z^{\prime}}, Y_{i, z}\right)$ are both isomorphic to the vector bundle (5.2.1) on $\bar{Z}_{i-\delta}^{\langle n+i-j-2 \delta\rangle}$.

Proof. Let $\left(\mathcal{B}_{z}, \lambda_{z}, \eta_{z}\right)$ and $\left(\mathcal{B}_{z^{\prime}}, \lambda_{z^{\prime}}, \eta_{z^{\prime}}\right)$ be the universal polarized abelian varieties on $\overline{\mathrm{Sh}}_{0, n}$ at $z$ and $z^{\prime}$, respectively. Then $Y_{i, z} \times{\overline{\operatorname{Sh}_{1, n-1}}} Y_{j, z^{\prime}}$ is the moduli space of tuples $\left(A, \lambda, \eta, \phi, \phi^{\prime}\right)$ where $\phi: \mathcal{B}_{z} \rightarrow A$ and $\phi^{\prime}: \mathcal{B}_{z^{\prime}} \rightarrow A$ are isogenies such that $\left(A, \lambda, \eta, \mathcal{B}_{z}, \lambda_{z}, \eta_{z}, \phi\right)$ and $\left(A, \lambda, \eta, \mathcal{B}_{z^{\prime}}, \eta_{z^{\prime}}, \phi^{\prime}\right)$ are points of $Y_{i, z}$ and $Y_{j, z^{\prime}}$ respectively.

Assume first that $Y_{i, z} \times{\overline{\operatorname{Sh}_{1, n-1}}} Y_{j, z^{\prime}}$ is nonempty, and let $\left(A, \lambda, \eta, \phi, \phi^{\prime}\right)$ be an $\overline{\mathbb{F}}_{p}$-valued point of it. Denote by $\tilde{\omega}_{A^{\vee}{ }^{\circ}{ }_{k} \subseteq}^{\circ} \subseteq \tilde{\mathcal{D}}(A)_{k}^{\circ}$ for $k=1,2$ the inverse image of $\omega_{A^{\vee} / \overline{\mathbb{F}}_{p}, k}^{\circ} \subseteq H_{1}^{\mathrm{dR}}\left(A / \overline{\mathbb{F}}_{p}\right)_{k}^{\circ} \cong \tilde{\mathcal{D}}(A)_{k}^{\circ} / p \tilde{\mathcal{D}}(A)_{k}^{\circ}$. We identify $\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{k}^{\circ}$ and $\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{k}^{\circ}$ with their images in $\tilde{\mathcal{D}}(A)_{k}^{\circ}$ via $\phi_{z, *, k}$ and $\phi_{z^{\prime}, *, k}$. Then we have a diagram of inclusions of $W\left(\overline{\mathbb{F}}_{p}\right)$-modules:


Here the numbers on the arrows indicate the $\overline{\mathbb{F}}_{p}$-dimensions of the cokernel of the corresponding inclusions, which we shall compute below. By the definition of $Y_{i}$ and $Y_{j}$, we have

$$
\operatorname{dim}_{\overline{\mathbb{F}}_{p}}\left(\tilde{\mathcal{D}}(A)_{1}^{\circ} / \tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ}\right)=\operatorname{dim}_{\overline{\mathbb{F}}_{p}} \operatorname{Coker}\left(\phi_{*, 1}\right)=i-1
$$

and similarly, $\operatorname{dim}_{\overline{\mathbb{F}}_{p}}\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ} / \tilde{\omega}_{A^{\vee}, 1}^{\circ}\right)=n-j$. Therefore, if we put

$$
\delta=\operatorname{dim}_{\overline{\mathbb{F}}_{p}}\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ}+\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ}\right) / \tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ}=\operatorname{dim}_{\overline{\mathbb{F}}_{p}} \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ} /\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ} \cap \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ}\right)
$$

we have $0 \leq \delta \leq \min \{i-1, n-j\}$. Moreover, the quasi-isogeny $\phi_{z, z^{\prime}}=\phi^{-1} \circ \phi^{\prime}$ : $\mathcal{B}_{z^{\prime}} \rightarrow \mathcal{B}_{z}$ makes $\mathcal{B}_{z^{\prime}}$ an element of $\operatorname{Isog}(z)$. We identify $\mathbb{L}_{z^{\prime}}$ defined in (4.11.1) with a $\mathbb{Z}_{p^{2}}$-lattice of $\mathbb{Z}_{z}[1 / p]$ via $\phi_{z^{\prime}, z, *, 1}$. Then

$$
\operatorname{dim}_{\mathbb{F}_{p^{2}}}\left(\mathbb{Q}_{z} \cap \mathbb{L}_{z^{\prime}}\right) / p \mathbb{L}_{z}=\operatorname{dim}_{\overline{\mathbb{F}}_{p}}\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ} \cap \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ}\right) / p \tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ}=n+i-j-\delta .
$$

Take a $\mathbb{Z}_{p^{2}}$-basis $\left(e_{1}, \ldots, e_{n}\right)$ of $\mathbb{L}_{z}$ such that the image of $\left(e_{j-i+\delta+1}, \ldots, e_{n}\right)$ in $\mathbb{L}_{z} / p \mathbb{L}_{z}$ form a basis of $\left(\mathbb{L}_{z} \cap \mathbb{L}_{z^{\prime}}\right) / p \mathbb{L}_{z}$ and such that $p^{-1} e_{n-\delta+1}, \ldots, p^{-1} e_{n}$ form a basis of $\left(\mathbb{L}_{z}+\mathbb{L}_{z^{\prime}}\right) / \mathbb{L}_{z}$. Then

$$
\begin{equation*}
\left(p e_{1}, \ldots, p e_{j-i+\delta}, e_{j-i+\delta+1}, \ldots, e_{n-\delta}, p^{-1} e_{n-\delta+1}, \ldots, p^{-1} e_{n}\right) \tag{6.4.2}
\end{equation*}
$$

is a basis of $\mathbb{L}_{z^{\prime}}$, that is $z^{\prime} \in R_{\mathfrak{p}}^{(j-i+\delta, n-\delta)} S_{\mathfrak{p}}^{-1}(z)$ according to the convention of Section 6.2.

Conversely, assume that there exists $\delta$ with $1 \leq \delta \leq \min \{i-1, n-j\}$ such that the point $z^{\prime} \in R_{\mathfrak{p}}^{(j-i+\delta, n-\delta)} S_{\mathfrak{p}}^{-1}(z)$. We have to prove statement (2), then the nonemptiness of $Y_{i, z} \times{\overline{S h_{1, n-1}}} Y_{j, z^{\prime}}$ will follow automatically. Let $\phi_{z^{\prime}, z}: \mathcal{B}_{z^{\prime}} \rightarrow \mathcal{B}_{z}$ be the unique quasi-isogeny which identifies $\mathbb{Z}_{z^{\prime}}$ with a $\mathbb{Z}_{p^{2}}$ lattice of $\mathbb{L}_{z}[1 / p]$. By the definition of $R_{\mathfrak{p}}^{(j-i+\delta, n-\delta)} S_{\mathfrak{p}}^{-1}$, there exists a basis $e_{1}, \ldots, e_{n}$ of $\mathbb{L}_{z}$ such that (6.4.2) is a basis of $\mathbb{L}_{z^{\prime}}$. One checks easily that $p\left(\mathbb{L}_{z}+\mathbb{L}_{z^{\prime}}\right) \subseteq \mathbb{L}_{z} \cap \mathbb{L}_{z^{\prime}}$. We put

$$
M_{k}=\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{k}^{\circ} \cap \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{k}^{\circ}\right) / p\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{k}^{\circ}+\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{k}^{\circ}\right)
$$

for $k=1,2$. Then one has

$$
\operatorname{dim}_{\overline{\mathscr{F}}_{p}}\left(M_{k}\right)=\operatorname{dim}_{\mathbb{F}_{p^{2}}}\left(\mathbb{L}_{z} \cap \mathbb{Z}_{z^{\prime}}\right) / p\left(\mathbb{L}_{z}+\mathbb{L}_{z^{\prime}}\right)=n+i-j-2 \delta .
$$

The Frobenius and Verschiebung on $\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)$ induce two bijective Frobenius semilinear maps $F: M_{1} \rightarrow M_{2}$ and $V^{-1}: M_{2} \rightarrow M_{1}$. We denote their linearizations by the same notation if no confusions arise. Let $Z_{\delta}\left(M_{\bullet}\right)$ be the moduli space which attaches to each locally noetherian $\overline{\mathbb{F}}_{p}$-scheme $S$ the set of isomorphism classes of pairs $\left(L_{1}, L_{2}\right)$, where $L_{1} \subseteq M_{1} \otimes_{\overline{\mathbb{F}}_{p}} \mathcal{O}_{S}$ and $L_{2} \subseteq M_{2} \otimes_{\overline{\mathbb{F}}_{p}} \mathcal{O}_{S}$ are subbundles of rank $i-\delta$ and $i-1-\delta$ respectively such that

$$
L_{2} \subseteq F\left(L_{1}^{(p)}\right), \quad V^{-1}\left(L_{2}^{(p)}\right) \subseteq L_{1}
$$

Note that there exists a basis $\left(\varepsilon_{k, 1}, \ldots, \varepsilon_{k, n+i-j-2 \delta}\right)$ of $M_{k}$ for $k=1,2$ under which the matrices of $F$ and $V^{-1}$ are both identity. Indeed, by solving a system of equations of Artin-Schreier type, one can take a basis $\left(\varepsilon_{1, \ell}\right)_{1 \leq \ell \leq n+i-j-2 \delta}$ for $M_{1}$ such that

$$
V^{-1}\left(F\left(\varepsilon_{1, \ell}\right)\right)=\varepsilon_{1, \ell} \quad \text { for all } 1 \leq \ell \leq n+i-j-2 \delta
$$

We put $\varepsilon_{2, \ell}=F\left(\varepsilon_{1, \ell}\right)$. Using these bases to identify both $M_{1}$ and $M_{2}$ with $\overline{\mathbb{F}}_{p}^{n+i-j-2 \delta}$, it is clear that $Z_{\delta}\left(M_{\bullet}\right)$ is isomorphic to the variety $\bar{Z}_{i-\delta}^{\langle n+i-j-2 \delta\rangle}$ considered in Section 5.2.

We have to establish an isomorphism between $Z_{\delta}\left(M_{\mathbf{\bullet}}\right)$ and $Y_{i, z} \times{ }_{\overline{S h}_{1, n-1}} Y_{j, z^{\prime}}$. Let ( $L_{1}, L_{2}$ ) be an $S$-point of $Z_{\delta}\left(M_{\bullet}\right)$. Note that there is a natural surjection

$$
\left(\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{k}^{\circ} \cap \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{k}^{\circ}\right) / p \tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{k}^{\circ}\right) \otimes_{\overline{\mathbb{F}}_{p}} \mathcal{O}_{S} \rightarrow M_{k} \otimes_{\overline{\mathbb{F}}_{p}} \mathcal{O}_{S}
$$

We define $H_{z, k}$ for $k=1,2$ to be the inverse image of $L_{k}$ under this surjection. Then $H_{z, k}$ can be naturally viewed as a subbundle of $\mathcal{D}\left(\mathcal{B}_{z}\right)_{k}^{\circ} \otimes_{\overline{\mathbb{F}}_{p}} \mathcal{O}_{S}$ of rank $i+1-k$, and we have $H_{z, 2} \subseteq F\left(H_{z, 1}^{(p)}\right)$ and $V^{-1}\left(H_{z, 2}^{(p)}\right) \subseteq H_{z, 1}$ since the pair $\left(L_{1}, L_{2}\right)$ verifies similar properties. Therefore, $\left(L_{1}, L_{2}\right) \mapsto\left(\mathcal{B}_{z, S}, \lambda_{z, S}, \eta_{z, S}, H_{z, 1}, H_{z, 2}\right)$ gives rise to a well-defined map $\varphi_{i, z}^{\prime}: Z_{\delta}\left(M_{\bullet}\right) \rightarrow Y_{i, z}^{\prime}$, where $\left(\mathcal{B}_{z, S}, \lambda_{z, S}, \eta_{z, S}\right)$ is the base change of $\left(\mathcal{B}_{z}, \lambda_{z}, \eta_{z}\right)$ to $S$. Similarly, we have a morphism $\varphi_{j, z^{\prime}}^{\prime}: Z_{\delta}\left(M_{\mathbf{0}}\right) \rightarrow Y_{j, z^{\prime}}^{\prime}$ defined by $\left(L_{1}, L_{2}\right) \mapsto\left(\mathcal{B}_{z^{\prime}, S}, \lambda_{z^{\prime}, S}, \eta_{z^{\prime}, S}, H_{z^{\prime}, 1}, H_{z^{\prime}, 2}\right)$, where $H_{z^{\prime}, k}$ is the inverse image of
$L_{k}$ under the natural surjection:

$$
\left(\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{k}^{\circ} \cap \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{k}^{\circ}\right) / p \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{k}^{\circ}\right) \otimes_{\overline{\mathbb{F}}_{p}} \mathcal{O}_{S} \rightarrow M_{k} \otimes_{\overline{\mathbb{F}}_{p}} \mathcal{O}_{S}
$$

By Proposition 4.8, we get two morphisms

$$
\varphi_{i, z}: Z_{\delta}\left(M_{\bullet}\right) \rightarrow Y_{i, z}, \quad \varphi_{j, z^{\prime}}: Z_{\delta}\left(M_{\bullet}\right) \rightarrow Y_{j, z^{\prime}}
$$

We claim that $\mathrm{pr}_{i} \circ \varphi_{i, z}=\operatorname{pr}_{j} \circ \varphi_{j, z}$, so that $\left(\varphi_{i, z}, \varphi_{j, z^{\prime}}\right)$ defines a map

$$
\varphi: Z_{\delta}\left(M_{\bullet}\right) \rightarrow Y_{i, z} \times_{\overline{S h}_{1, n-1}} Y_{j, z^{\prime}} .
$$

Since $Y_{i, z} \times{ }_{\overline{S h}_{1, n-1}} Y_{j, z^{\prime}}$ is separated, the locus where $\mathrm{pr}_{i} \circ \varphi_{i, z}$ coincides with $\mathrm{pr}_{j} \circ \varphi_{j, z}$ is a closed subscheme of $Z_{\delta}\left(M_{\bullet}\right)$. As $Z_{\delta}\left(M_{\bullet}\right)$ is reduced, it is enough to show $\operatorname{pr}_{i}\left(\varphi_{i, z}(x)\right)=\operatorname{pr}_{j}\left(\varphi_{j, z}(x)\right)$ for each closed geometric point $x=\left(L_{1}, L_{2}\right) \in$ $Z_{\delta}\left(M_{.}\right)\left(\overline{\mathbb{F}}_{p}\right)$. Let $\left(A, \lambda, \eta, \mathcal{B}_{z}, \lambda_{z}, \eta_{z}, \phi\right)$ and $\left(A^{\prime}, \lambda^{\prime}, \eta^{\prime}, \mathcal{B}_{z^{\prime}}, \lambda_{z^{\prime}}, \eta_{z^{\prime}}^{\prime}, \phi^{\prime}\right)$ be respectively the image of $\left(L_{1}, L_{2}\right)$ under $\varphi_{i, z}$ and $\varphi_{j, z^{\prime}}$. To prove the claim, we have to show that there is an isomorphism $(A, \lambda, \eta) \cong\left(A^{\prime}, \lambda^{\prime}, \eta^{\prime}\right)$ as objects of $\overline{\operatorname{Sh}}_{1, n-1}$. We identify $\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right), \tilde{\mathcal{D}}(A), \tilde{\mathcal{D}}\left(A^{\prime}\right)$ with $W\left(\overline{\mathbb{F}}_{p}\right)$-lattices of $\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)[1 / p]$ via the quasiisogenies $\phi_{z^{\prime}, z}: \mathcal{B}_{z^{\prime}} \rightarrow \mathcal{B}_{z}, \phi^{-1}: A \rightarrow \mathcal{B}_{z}$ and $\phi_{z^{\prime}, z}^{-1} \circ \phi^{\prime}: A^{\prime} \rightarrow \mathcal{B}_{z}$. Then by the construction of $A$ (cf., the proof of Proposition 4.8), $\tilde{\mathcal{D}}(A)_{1}^{\circ}$ and $\tilde{\omega}_{A^{\vee}, 1}^{\circ}$ fit into the diagram (6.4.1) such that there is a canonical isomorphism

$$
\begin{align*}
L_{1} & \cong \tilde{\omega}_{A^{\vee}, 1}^{\circ} / p\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ}+\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ}\right) \\
& \subseteq\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ} \cap \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ}\right) / p\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ}+\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ}\right)=M_{1} \tag{6.4.3}
\end{align*}
$$

Similarly, we have

$$
\begin{align*}
L_{2} & \cong p \tilde{\omega}_{A^{\vee}, 2}^{\circ} / p\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{2}^{\circ}+\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{2}^{\circ}\right) \\
& \subseteq\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{2}^{\circ} \cap \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{2}^{\circ}\right) / p\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{2}^{\circ}+\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{2}^{\circ}\right)=M_{2} \tag{6.4.4}
\end{align*}
$$

It is easy to see that such relations determine $\tilde{\mathcal{D}}(A)$ uniquely from $\left(L_{1}, L_{2}\right)$. But the same argument shows that the same relations are satisfied with $A$ replaced by $A^{\prime}$. Hence, we see that the quasi-isogeny $f$ induces an isomorphism between the Dieudonné modules of $A$ and $A^{\prime}$. As $f$ is a $p$-quasi-isogeny, this implies immediately that $f$ is an isomorphism of abelian varieties, proving the claim.

It remains to prove that $\varphi: Z_{\delta}\left(M_{\bullet}\right) \xrightarrow{\sim} Y_{i, z} \times \overline{S \overline{S h}}_{1, n-1} Y_{j, z^{\prime}}$ is an isomorphism. It suffices to show that $\varphi$ induces bijections on closed points and tangents spaces. The argument is similar to the proof of Proposition 4.8. Indeed, given a closed point $x=\left(A, \lambda, \eta, \phi, \phi^{\prime}\right)$ of $Y_{i, z} \times{ }_{\text {Sh }_{1, n-1}} Y_{j, z^{\prime}}$, one can construct a unique point $y=\left(L_{1}, L_{2}\right)$ of $Z_{\delta}\left(M_{.}\right)$with $\varphi(y)=x$ by the relations (6.4.3) and (6.4.4). It follows immediately that $\varphi$ induces a bijection on closed points. Let $x$ and $y$ be as above. By the same argument as in Proposition 4.4, the tangent space of $Z_{\delta}\left(M_{\bullet}\right)$ at $y$ is given by

$$
T_{Z_{\delta}\left(M_{\bullet}\right), y} \cong\left(L_{1} / V^{-1}\left(L_{2}^{(p)}\right)\right)^{*} \otimes\left(M_{1} / L_{1}\right) \oplus L_{2}^{*} \otimes F\left(L_{1}^{(p)}\right) / L_{2}
$$

On the other hand, using Grothendieck-Messing deformation theory, one sees easily that the tangent space of $Y_{i, z} \times{\overline{\operatorname{Sh}_{1, n-1}}} Y_{j, z^{\prime}}$ at $x$ is given by

$$
\begin{aligned}
T_{Y_{i, z} \times \overline{S \overline{S h}} 1, n-1} Y_{j, z^{\prime}, x} \cong \operatorname{Hom}_{\overline{\mathbb{T}}_{p}} & \left(\omega_{A^{\vee}, 1}^{\circ},\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{1}^{\circ} \cap \tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ}\right) / \tilde{\omega}_{A^{\vee}, 1}^{\circ}\right) \\
& \oplus \operatorname{Hom}_{\overline{\mathbb{F}}_{p}}\left(\tilde{\omega}_{A^{\vee}, 2}^{\circ} /\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{2}^{\circ}+\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{2}^{\circ}\right), \tilde{\mathcal{D}}(A)_{2}^{\circ} / \tilde{\omega}_{A^{\vee}, 2}^{\circ}\right) .
\end{aligned}
$$

From (6.4.3) and (6.4.4), we see easily that

$$
\begin{aligned}
\omega_{A^{\vee}, 1}^{\circ} \cong L_{1} / V^{-1}\left(L_{2}^{(p)}\right), & & \left.\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}\right)_{1}^{\circ}\right) / \tilde{\omega}_{A^{\vee}, 1}^{\circ} \cong M_{1} / L_{1} \\
\tilde{\omega}_{A^{\vee}, 2}^{\circ} /\left(\tilde{\mathcal{D}}\left(\mathcal{B}_{z}\right)_{2}^{\circ}+\tilde{\mathcal{D}}\left(\mathcal{B}_{z^{\prime}}^{\circ}\right)_{2}^{\circ}\right) \cong L_{2}, & & \tilde{\mathcal{D}}(A)_{2}^{\circ} / \tilde{\omega}_{A^{\vee}, 2}^{\circ} \cong F\left(L_{1}^{(p)}\right) / L_{2}
\end{aligned}
$$

It follows that $\varphi$ induces a bijection between $T_{Z_{\delta}\left(M_{\bullet}\right), y}$ and $T_{Y_{i, z} \times \overline{S h}_{1, n-1}} Y_{j, z^{\prime}, x}$. This finishes the proof of Proposition 6.4.
6.5. Applications to cohomology. Recall that we have a morphism $\mathcal{J} \mathcal{L}_{j}(4.16 .1)$ for each $j=1, \ldots, n$. We consider another map in the opposite direction:

$$
v_{j}: H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right) \xrightarrow{\mathrm{pr}_{j}^{*}} H_{\mathrm{et}}^{2(n-1)}\left(\bar{Y}_{j}, \overline{\mathbb{Q}}_{\ell}\right) \xrightarrow{\sim} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right),
$$

where the second isomorphism is induced by the trace map

$$
\operatorname{Tr}_{\mathrm{pr}_{j}^{\prime}}: R^{2(n-1)} \operatorname{pr}_{j, *}^{\prime} \overline{\mathbb{Q}}_{\ell}(n-1) \xrightarrow{\sim} \overline{\mathbb{Q}}_{\ell}
$$

For $1 \leq i, j \leq n$, we define
$m_{i, j}=v_{j} \circ \mathcal{J} \mathcal{L}_{i}: H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right) \xrightarrow{\mathcal{J} \mathcal{L}_{i}} H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right) \xrightarrow{v_{j}} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)$.
Putting all the morphisms $\mathcal{J}_{i}$ and $v_{j}$ together, we get a sequence of morphisms:

$$
\begin{align*}
& \bigoplus_{i=1}^{n} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right) \xrightarrow{\mathcal{J L}} H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{S}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right) \\
& \xrightarrow{\nu=\left(v_{1}, \ldots, v_{n}\right)} \bigoplus_{j=1}^{n} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{~S}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right) . \tag{6.5.1}
\end{align*}
$$

We see that the composed morphism above is given by the matrix $M=\left(m_{i, j}\right)_{1 \leq i, j \leq n}$, and we call it the intersection matrix of cycles $Y_{j}$ on $\mathrm{Sh}_{1, n-1}$. All these morphisms are equivariant under the natural action of the Hecke algebra $\mathscr{H}\left(K^{p}, \overline{\mathbb{Q}}_{\ell}\right)$. We describe the intersection matrix in terms of the Hecke action of $\overline{\mathbb{Q}}_{\ell}\left[K_{\mathfrak{p}} \backslash \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right) / K_{\mathfrak{p}}\right]$ on $H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)$.

The group $H_{\mathrm{et}}^{0}\left(\overline{\operatorname{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)$ is the space of functions on $\operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$ with values in $\overline{\mathbb{Q}}_{\ell}$. For $z \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$, let $e_{z}$ denote the characteristic function of $z$. Then the image of $z$ under $K_{\mathfrak{p}} \gamma K_{\mathfrak{p}}$ for $\gamma \in \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$ is

$$
\left[K_{\mathfrak{p}} \gamma K_{\mathfrak{p}}\right]_{*}\left(e_{z}\right)=\sum_{z^{\prime} \in T_{\mathfrak{p}}(\gamma)(z)} e_{z^{\prime}}
$$

where $T_{\mathfrak{p}}(\gamma)(z)$ means the set theoretic Hecke correspondence defined in Section 6.2. In the sequel, we will use the same notation $T_{\mathfrak{p}}(\gamma)$ to denote the action of [ $K_{\mathfrak{p}} \gamma K_{\mathfrak{p}}$ ] on $H_{\mathrm{et}}^{0}\left(\overline{\mathrm{~S}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)$. In particular, we have Hecke operators $T_{\mathfrak{p}}^{(i)}, S_{\mathfrak{p}}, R_{\mathfrak{p}}^{(a, b)}, \ldots$.

Proposition 6.6. For $1 \leq i \leq j \leq n$, we have

$$
\begin{aligned}
& m_{i, j}=\sum_{\delta=0}^{\min \{i-1, n-j\}} N(n+i-j-2 \delta, i-\delta) R_{\mathfrak{p}}^{(j-i+\delta, n-\delta)} S_{\mathfrak{p}}^{-1}, \\
& m_{j, i}=\sum_{\delta=0}^{\min \{i-1, n-j\}} N(n+i-j-2 \delta, i-\delta) R_{\mathfrak{p}}^{(\delta, n+i-j-\delta)} S_{\mathfrak{p}}^{-1},
\end{aligned}
$$

where $N(n+i-j-2 \delta, i-\delta)$ are the fundamental intersection numbers defined by (5.2.2).

Proof. We have a commutative diagram:


Here, the vertical arrows are cycle class maps, and $\mathrm{pr}_{j}^{\prime}$ is the refined Gysin map defined in (6.3.1). For $z \in \operatorname{Sh}_{0, n}\left(\overline{\mathbb{F}}_{p}\right)$, the image of $e_{z}$ under $m_{i, j}$ is given by

$$
\begin{aligned}
m_{i, j}\left(e_{z}\right) & =\operatorname{Tr}_{\operatorname{pr}_{j}^{\prime}} \operatorname{pr}_{j}^{*} \operatorname{Gys}_{\operatorname{pr}_{i}} \operatorname{cl}\left(\left[Y_{i, z}\right]\right)=\operatorname{Tr}_{\mathrm{pr}_{j}^{\prime}}\left(\operatorname{cl}^{\left.\left(\operatorname{pr}_{j}^{\prime} \operatorname{pr}_{i, *}\left[Y_{i, z}\right]\right)\right)}\right. \\
& =\operatorname{Tr}_{\mathrm{pr}_{j}^{\prime}}\left(\sum_{z^{\prime} \in \operatorname{Sh}_{0, n}\left(\bar{F}_{p}\right)} \operatorname{cl}\left(c_{r\left(z, z^{\prime}\right)}\left(\mathcal{E}\left(Y_{j, z^{\prime}}, Y_{i, z}\right)\right)\right) \cdot \operatorname{cl}\left(Y_{j, z^{\prime}} \times_{\overline{\operatorname{Sh}}_{1, n-1}} Y_{i, z}\right)\right) \\
& =\sum_{z^{\prime} \in \operatorname{Sh}_{0, n}\left(\mathbb{F}_{p}\right)}\left(\int_{Y_{j, z^{\prime}} \times \operatorname{Sh}_{1, n-1} Y_{i, z}} c_{r\left(z, z^{\prime}\right)}\left(\mathcal{E}\left(Y_{j, z^{\prime}}, Y_{i, z}\right)\right)\right) e_{z^{\prime}},
\end{aligned}
$$

where $r\left(z, z^{\prime}\right)$ is the rank of $\mathcal{E}\left(Y_{j, z^{\prime}}, Y_{i, z}\right)$, and we used (6.3.2) in the second step. Indeed, Proposition 6.4(1) says that the schematic intersection $Y_{i, z} \times_{\overline{S h}_{1, n-1}} Y_{j, z^{\prime}}$ is smooth, so the closed immersion $Y_{i, z} \times_{\overline{S h}_{1, n-1}} Y_{j, z^{\prime}} \hookrightarrow Y_{j, z^{\prime}}$ is a regular immersion and the assumptions for (6.3.2) are thus satisfied here.

By Proposition 6.4(1), $e_{z^{\prime}}$ has a nonzero contribution to the summation above if and only if there exists an integer $\delta$ with $0 \leq \delta \leq \min \{i-1, n-j\}$ such that $z^{\prime} \in R_{\mathfrak{p}}^{(j-i+\delta, n-\delta)} S_{\mathfrak{p}}^{-1}(z)$. In that case, Proposition 6.4(2) implies that the coefficient of $e_{z^{\prime}}$ is nothing but the fundamental intersection number $N(n+i-j-2 \delta, i-\delta)$ defined in (5.2.2). The formula for $m_{i, j}$ now follows immediately. The formula for $m_{j, i}$ is proved in the same manner.

If we express $m_{i, j}$ in terms of the elementary Hecke operators $T_{\mathfrak{p}}^{(k)}$, we get the following.

Theorem 6.7. Put $d(n, k)=(2 k-1) n-2 k(k-1)-1$ for integers $1 \leq k \leq n$. Then, for $1 \leq i \leq j \leq n$, we have

$$
\begin{aligned}
m_{i, j} & =\sum_{\delta=0}^{\min \{i-1, n-j\}}(-1)^{n+1+i-j}(n+i-j-2 \delta) p^{d(n+i-j-2 \delta, i-\delta)} T_{\mathfrak{p}}^{(j-i+\delta)} T_{\mathfrak{p}}^{(n-\delta)} S_{\mathfrak{p}}^{-1}, \\
m_{j, i} & =\sum_{\delta=0}^{\min \{i-1, n-j\}}(-1)^{n+1+i-j}(n+i-j-2 \delta) p^{d(n+i-j-2 \delta, i-\delta)} T_{\mathfrak{p}}^{(\delta)} T_{\mathfrak{p}}^{(n+i-j-\delta)} S_{\mathfrak{p}}^{-1} .
\end{aligned}
$$

Proof. We prove only the statement for $m_{i, j}$, and that for $m_{j, i}$ is similar. By Proposition A. 1 in Appendix A, the right hand side of the first formula above is

$$
\begin{aligned}
& \sum_{\delta=0}^{\min \{i-1, n-j\}}(-1)^{n+1-i-j}(n+i-j-2 \delta) p^{d(n+i-j-2 \delta, i-\delta)} \\
& \cdot\left(\sum_{k=0}^{\delta}\binom{n+i-j-2 \delta+2 k}{k}_{p^{2}} R_{\mathfrak{p}}^{(j-i+\delta-k, n-\delta+k)} S_{\mathfrak{p}}^{-1}\right) \\
& \\
& =\sum_{r=0}^{\min \{i-1, n-j\}}(\star) R_{\mathfrak{p}}^{(j-i+r, n-r)} S_{\mathfrak{p}}^{-1} .
\end{aligned}
$$

Here, we have put $r=\delta-k$, and the expression $\star$ in the parentheses is

$$
\begin{aligned}
& \star=\sum_{k=0}^{\min \{i-1-r, n-j-r\}}(-1)^{n+1+i-j}(n+i-j-2 r-2 k) \\
&=N(n+i-j-2 r, i-r) .
\end{aligned} \quad \cdot p^{d(n+i-j-2 r-2 k, i-r-k)}\binom{n+i-j-2 r}{k}_{p^{2}} .
$$

Here, the last equality is Theorem 5.3. The statement for $m_{i, j}$ now follows from Proposition 6.6.

Example 6.8. We write down explicitly the intersection matrices when $n$ is small.
(1) Consider first the case $n=2$. This case is essentially the same as the Hilbert quadratic case studied in [Tian and Xiao 2014], and the intersection matrix can be written:

$$
M=\left(\begin{array}{cc}
-2 p & T_{\mathfrak{p}}^{(1)} \\
T_{\mathfrak{p}}^{(1)} S_{\mathfrak{p}}^{-1} & -2 p
\end{array}\right)
$$

(2) When $n=3$, Theorem 6.7 gives

$$
M=\left(\begin{array}{ccc}
3 p^{2} & -2 p T_{\mathfrak{p}}^{(1)} & T_{\mathfrak{p}}^{(2)} \\
-2 p T_{\mathfrak{p}}^{(2)} S_{\mathfrak{p}}^{-1} & 3 p^{4}+T_{\mathfrak{p}}^{(1)} T_{\mathfrak{p}}^{(2)} S_{\mathfrak{p}}^{-1} & -2 p T_{\mathfrak{p}}^{(1)} \\
T_{\mathfrak{p}}^{(1)} S_{\mathfrak{p}}^{-1} & -2 p T_{\mathfrak{p}}^{(2)} S_{\mathfrak{p}}^{-1} & 3 p^{2}
\end{array}\right)
$$

(3) The intersection matrix for $n=4$ can be written:

$$
M=\left(\begin{array}{cccc}
-4 p^{3} & 3 p^{2} T_{\mathfrak{p}}^{(1)} & -2 p T_{\mathfrak{p}}^{(2)} & T_{\mathfrak{p}}^{(3)} \\
3 p^{2} T_{\mathfrak{p}}^{(3)} S_{\mathfrak{p}}^{-1} & -4 p^{7}-2 p T_{\mathfrak{p}}^{(1)} T_{\mathfrak{p}}^{(3)} S_{\mathfrak{p}}^{-1} & 3 p^{4} T_{\mathfrak{p}}^{(1)}+T_{\mathfrak{p}}^{(2)} T_{\mathfrak{p}}^{(3)} S_{\mathfrak{p}}^{-1} & -2 p T_{\mathfrak{p}}^{(2)} \\
-2 p T_{\mathfrak{p}}^{(2)} S_{\mathfrak{p}}^{-1} 3 p^{4} T_{\mathfrak{p}}^{(3)} S_{\mathfrak{p}}^{-1}+T_{\mathfrak{p}}^{(1)} T_{\mathfrak{p}}^{(2)} S_{\mathfrak{p}}^{-1} & -4 p^{7}-2 p T_{\mathfrak{p}}^{(1)} T_{\mathfrak{p}}^{(3)} S_{\mathfrak{p}}^{-1} 3 p^{2} T_{\mathfrak{p}}^{(1)} \\
T_{\mathfrak{p}}^{(1)} S_{\mathfrak{p}}^{-1} & -2 p T_{\mathfrak{p}}^{(2)} S_{\mathfrak{p}}^{-1} & 3 p^{2} T_{\mathfrak{p}}^{(3)} S_{\mathfrak{p}}^{-1} & -4 p^{3}
\end{array}\right)
$$

6.9. Proof of Theorem 4.18(1). Let $\pi \in \mathscr{A}_{K}$ as in the statement of Theorem 4.18(1). Consider the $\left(\pi^{p}\right)^{K^{p}}$-isotypic direct factor of the $\mathscr{H}\left(K^{p}, \overline{\mathbb{Q}}_{\ell}\right)$-equivariant sequence (6.5.1):

$$
\begin{align*}
& \bigoplus_{i=1}^{n} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)_{\pi^{p}} \xrightarrow{\mathcal{J \mathcal { L }}_{\pi}} H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi^{p}} \\
& \xrightarrow{\nu_{\pi}} \bigoplus_{j=1}^{n} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)_{\pi^{p}} . \tag{6.9.1}
\end{align*}
$$

In particular, when $i=j=1, \nu_{1} \circ \mathcal{J L}_{1}$ is given by multiplication by $-n p^{n-1}$. So the $\pi^{p}$-isotypic component of (6.9.1) is nonzero. This implies that $\pi^{p}$ appears in $H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)$, i.e., there exist admissible irreducible representations $\pi_{p}^{\prime}$ of $G_{1, n-1}\left(\mathbb{Q}_{p}\right)$ and $\pi_{\infty}^{\prime}$ of $G_{1, n-1}(\mathbb{R})$, which is cohomological in degree $n-1$, such that $\pi^{p} \otimes \pi_{p}^{\prime} \otimes \pi_{\infty}^{\prime}$ is a cuspidal automorphic representation $\pi^{\prime} \otimes \pi_{\infty}^{\prime}$ of $G_{1, n-1}\left(\mathbb{A}_{\mathbb{Q}}\right)$. By Lemma 4.17, $\pi^{\prime} \simeq \pi$ satisfies Hypothesis $2.5(2)$ for $a_{0}=(1, n-1)$. Thus, taking the $\pi^{p}$-isotypic component of (6.9.1) is the same as taking its $\pi$ isotypic component. From now on, we use subscript $\pi$ in places of subscript $\pi^{p}$.

If $a_{\mathfrak{p}}^{(i)}$ denotes the eigenvalues of $T_{\mathfrak{p}}^{(i)}$ on $\pi_{\mathfrak{p}}^{K_{\mathfrak{p}}}$ for each $1 \leq i \leq n$, then $T_{\mathfrak{p}}^{(i)}$ acts as the scalar $a_{\mathfrak{p}}^{(i)}$ on all the terms in (6.9.1). Therefore, $v_{\pi} \circ \mathcal{J} \mathcal{L}_{\pi}$ is given by the matrix $M_{\pi}$, which is obtained by replacing $T_{\mathfrak{p}}^{(i)}$ by $a_{\mathfrak{p}}^{(i)}$ in each entry of $M$. By definition, the $\alpha_{\pi_{\mathrm{p}}, i}$ are the roots of the Hecke polynomial (2.6.2):

$$
X^{n}+\sum_{i=1}^{n}(-1)^{i} p^{i(i-1)} a_{\mathfrak{p}}^{(i)} X^{n-i}
$$

Then Theorem 4.18(1) follows easily from the following.
Lemma 6.10. We have

$$
\operatorname{det}\left(M_{\pi}\right)= \pm p^{\frac{n\left(n^{2}-1\right)}{3}} \frac{\prod_{i<j}\left(\alpha_{\pi_{\mathfrak{p}}, i}-\alpha_{\pi_{\mathfrak{p}}, j}\right)^{2}}{\left(\prod_{i=1}^{n} \alpha_{\pi_{\mathfrak{p}}, i}\right)^{n-1}}
$$

Here, $\pm$ means that the formula holds up to sign. In particular, $v_{\pi} \circ \mathcal{J L}_{\pi}$ is an isomorphism if the $\alpha_{\pi_{\mathfrak{p}}, i}$ are distinct.

Proof. Put $\beta_{i}=\alpha_{\pi_{\mathfrak{p}}, i} / p^{n-1}$ for $1 \leq i \leq n$. For $i=1, \ldots, n$, let $s_{i}$ be the $i$-th elementary symmetric polynomial in $\beta_{1}, \ldots, \beta_{n}$. Then we have $a_{\mathfrak{p}}^{(i)}=p^{i(n-i)} s_{i}$. It follows from Theorem 6.7 that the $(i, j)$-entry of $M_{\pi}$ with $1 \leq i \leq j \leq n$ is given by

$$
\begin{aligned}
& m_{i, j}(\pi)=s_{n}^{-1} \sum_{\delta=0}^{\min \{i-1, n-j\}}(-1)^{n+1+i-j}(n+i-j-2 \delta) \\
& \cdot p^{d(n+i-j-2 \delta, i-\delta)+(j-i+\delta)(n+i-j-\delta)+\delta(n-\delta)} s_{j-i+\delta} s_{n-\delta}
\end{aligned}
$$

A direct computation shows that the exponent index on $p$ in each term above is independent of $\delta$, and is equal to $e(i, j):=(n+1)(i+j-1)-\left(i^{2}+j^{2}\right)$. The same holds when $i>j$. In summary, we get $m_{i, j}(\pi)=s_{n}^{-1} p^{e(i, j)} m_{i, j}^{\prime}(\pi)$ with

$$
m_{i, j}^{\prime}(\pi)= \begin{cases}\sum_{\delta=0}^{\min \{i-1, n-j\}}(-1)^{n+1+i-j}(n+i-j-2 \delta) s_{j-i+\delta} s_{n-\delta}, & \text { if } i \leq j, \\ \sum_{\delta=0}^{\min \{j-1, n-i\}}(-1)^{n+1+j-i}(n+j-i-2 \delta) s_{\delta} s_{n+j-i-\delta}, & \text { if } i>j\end{cases}
$$

For any $n$-permutation $\sigma$, we have

$$
\sum_{i=1}^{n} e(i, \sigma(i))=\frac{n\left(n^{2}-1\right)}{3}
$$

Thus we get $\operatorname{det}\left(M_{\pi}\right)=p^{n\left(n^{2}-1\right) / 3} s_{n}^{-n} \operatorname{det}\left(m_{i, j}^{\prime}(\pi)\right)$. The rest of the computation is purely combinatorial, which is the case $q=-1$ of Theorem B. 1 in Appendix B.

Remark 6.11. We point out that the determinant of the intersection matrix computed by Theorem B. 1 holds with an auxiliary variable $q$. A similar phenomenon also appeared in the case of Hilbert modular varieties [Tian and Xiao 2014], where the computation was related to the combinatorial model of periodic semimeanders. These motivate us to ask, out of curiosity, whether there might be some quantum version of the construction of cycles, or even Conjecture 2.12, possibly for the geometric Langlands setup.
6.12. Proof of Theorem 4.18(2). Given Theorem 4.18(1), it suffices to prove that

$$
\begin{equation*}
n \operatorname{dim} H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)_{\pi} \geq \operatorname{dim} H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}(n-1)\right)_{\pi}^{\mathrm{fin}} \tag{6.12.1}
\end{equation*}
$$

Actually, by (2.4.1) and (2.6.3), we have
$H_{\mathrm{et}}^{0}\left(\overline{\mathrm{Sh}}_{0, n}, \overline{\mathbb{Q}}_{\ell}\right)_{\pi}=\pi^{K} \otimes R_{(0, n), \ell}(\pi), \quad H_{\mathrm{et}}^{2(n-1)}\left(\overline{\mathrm{Sh}}_{1, n-1}, \overline{\mathbb{Q}}_{\ell}\right)_{\pi}=\pi^{K} \otimes R_{(1, n-1), \ell}(\pi)$.
Write $\pi_{p}=\pi_{p, 0} \otimes \pi_{\mathfrak{p}}$ as a representation of $G\left(\mathbb{Q}_{p}\right) \simeq \mathbb{Q}_{p}^{\times} \times \mathrm{GL}_{n}\left(E_{\mathfrak{p}}\right)$. Let $\chi_{\pi_{p, 0}}$ : $\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{2}}\right) \rightarrow \overline{\mathbb{Q}}_{\ell}^{\times}$denote the character sending $\operatorname{Frob}_{p^{2}}$ to $\pi_{p, 0}\left(p^{2}\right)$, and let $\rho_{\pi_{\mathfrak{p}}}$
be as in (2.6.1). According to (2.6.3), up to semisimplification, we have

$$
\begin{align*}
{\left[R_{(0, n), \ell}(\pi)\right]=\# } & \operatorname{ker}^{1}\left(\mathbb{Q}, G_{0, n}\right) m_{0, n}(\pi)  \tag{6.12.2}\\
& {\left[\wedge^{n} \rho_{\pi_{\mathfrak{p}}} \otimes \chi_{\pi_{p, 0}}^{-1} \otimes \overline{\mathbb{Q}}_{\ell}\left(\frac{1}{2} n(n-1)\right)\right] } \\
{\left[R_{(1, n-1), \ell}(\pi)\right]=\# } & \operatorname{ker}^{1}\left(\mathbb{Q}, G_{1, n-1}\right) m_{0, n}(\pi)  \tag{6.12.3}\\
& {\left[\rho_{\pi_{\mathfrak{p}}} \otimes \wedge^{n-1} \rho_{\pi_{\mathfrak{p}}} \otimes \chi_{\pi_{p, 0}}^{-1} \otimes \overline{\mathbb{Q}}_{\ell}\left(\frac{1}{2}(n-1)(n-2)\right)\right] . }
\end{align*}
$$

Note that

$$
\begin{aligned}
& \operatorname{dim}\left(\rho_{\pi_{\mathfrak{p}}} \otimes \wedge^{n-1} \rho_{\pi_{\mathfrak{p}}} \otimes \chi_{\pi_{p, 0}}^{-1} \otimes \overline{\mathbb{Q}}_{\ell}\left(\frac{(n-1)(n-2)}{2}\right)\right)^{\mathrm{fin}} \\
&=\sum_{\zeta} \operatorname{dim}\left(\rho_{\pi_{\mathfrak{p}}} \otimes \wedge^{n-1} \rho_{\pi_{\mathfrak{p}}}\right)^{\mathrm{Frob}_{p^{2}}=p^{n(n-1)} \zeta},
\end{aligned}
$$

where the superscript "fin" means taking the subspace on which $\operatorname{Gal}\left(\overline{\mathbb{F}}_{p} / \mathbb{F}_{p^{2}}\right)$ acts through a finite quotient, and $\zeta$ runs through all roots of unity. If $\alpha_{\pi_{\mathfrak{p}}, i} / \alpha_{\pi_{\mathfrak{p}}, j}$ is not a root of unity for any pair $i \neq j$, the right hand side above is equal to the sum of the multiplicities of $\prod_{i=1}^{n} \alpha_{\pi_{\mathfrak{p}}, i}=p^{n(n-1)} \zeta$ as eigenvalues of $\left(\rho_{\pi_{\mathfrak{p}}} \otimes \wedge^{n-1} \rho_{\pi_{\mathfrak{p}}}\right)\left(\operatorname{Frob}_{p^{2}}\right)$, which is $n$. Therefore, under these conditions on the $\alpha_{\pi_{\mathfrak{p}}, i}$, we have by (6.12.3)

$$
\operatorname{dim} R_{(1, n-1), \ell}(\pi)^{\mathrm{fin}} \leq n \cdot \# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{1, n-1}\right) \cdot m_{1, n-1}(\pi)
$$

and the equality holds if $\mathrm{Frob}_{p^{2}}$ is semisimple on $R_{(1, n-1), \ell}(\pi)$. On the other hand, we have from (6.12.2)

$$
\operatorname{dim} R_{(0, n), \ell}(\pi)=\# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{0, n}\right) \cdot m_{0, n}(\pi)
$$

By a result of White [2012, Theorem E], the multiplicity $m_{a_{\mathbf{\bullet}}}(\pi)$ above is equal to 1 for $a_{\bullet}=(1, n-1)$ and $a_{\bullet}=(0, n)$. Now the inequality $(6.12 .1)$ follows immediately from this and the fact that $\# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{1, n-1}\right)=\# \operatorname{ker}^{1}\left(\mathbb{Q}, G_{0, n}\right)$. This finishes the proof of Theorem 4.18(2).

## 7. Construction of cycles in the case of $G(U(r, s) \times U(s, r))$

We keep the notation of Section 3.6. In this section, we will give the construction of certain cycles on Shimura varieties for $G(U(r, s) \times U(s, r))$. We always assume that $s \geq r$.
7.1. Description of the cycles in terms of Dieudonné modules. Let $\delta$ be a nonnegative integer with $\delta \leq r$. We consider the case of Conjecture 2.12 when $n=r+s$, $a_{1}=r, a_{2}=s, b_{1}=r-\delta$, and $b_{2}=s+\delta$. The representation $r_{a_{0}}$ of $\mathrm{GL}_{n}$ involved is

$$
r_{a_{0}}=\wedge^{r} \operatorname{Std} \otimes \wedge^{s} \operatorname{Std}
$$

The weight $\lambda$ of Conjecture 2.12 is

$$
\lambda=(\underbrace{2, \ldots, 2}_{r-\delta}, \underbrace{1, \ldots, 1}_{s-r+2 \delta}, \underbrace{0, \ldots, 0}_{r-\delta}) .
$$

By elementary calculation of representations of $\mathrm{GL}_{n}$, the multiplicity of $\lambda$ in $r_{a_{0}}$ is $m_{\lambda}\left(a_{0}\right)=\binom{s-r+2 \delta}{\delta}$. Then Conjecture 2.12 thus predicts the existence of $\binom{s-r+2 \delta}{\delta}$ cycles $Y_{j}$ on $\mathrm{Sh}_{r, s}$, each of dimension

$$
\frac{1}{2}\left(\operatorname{dim} \mathrm{Sh}_{r, s}+\operatorname{dim} \mathrm{Sh}_{r-\delta, s+\delta}\right)=\frac{1}{2}(2 r s+2(r-\delta)(s+\delta))=2 r s-(s-r) \delta-\delta^{2}
$$

and each admits a rational map to $\mathrm{Sh}_{r-\delta, s+\delta}$. The principal goal of this section is to construct these cycles, at least conjecturally. We start with the description in terms of the Dieudonné modules at closed points.

Consider the interval $[r-\delta, s+\delta]$; it contains $s-r+2 \delta$ unit segments with integer endpoints. We will parametrize the cycles on the Shimura variety by the subsets of these $s-r+2 \delta$ unit segments of cardinality $\delta$. There are exactly $\binom{s-r+2 \delta}{\delta}$ such subsets. Let $\boldsymbol{j}$ be one of them. Then we can write the union of all the segments in $\boldsymbol{j}$ as

$$
\begin{equation*}
\left[j_{1,1}, j_{1,2}\right] \cup\left[j_{2,1}, j_{2,2}\right] \cup \cdots \cup\left[j_{\epsilon, 1}, j_{\epsilon, 2}\right] \tag{7.1.1}
\end{equation*}
$$

such that all $j_{\alpha, i}$ are integers,

$$
r-\delta \leq j_{1,1}<j_{1,2}<j_{2,1}<j_{2,2}<\cdots<j_{\epsilon, 1}<j_{\epsilon, 2} \leq s+\delta
$$

and we have $\sum_{\alpha=1}^{\epsilon}\left(j_{\alpha, 2}-j_{\alpha, 1}\right)=\delta$. For notational convenience, we put $j_{0,1}=$ $j_{0,2}=0$.

We define $Z_{j}$ to be the subset of $\overline{\mathbb{F}}_{p}$-points $z$ of $\mathrm{Sh}_{r, s}$ such that the reduced Dieudonné modules $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}$ and $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ}$ contain submodules $\tilde{\mathcal{E}}_{1}$ and $\tilde{\mathcal{E}}_{2}$ satisfying (3.2.1) for $m=\epsilon$, i.e.,

$$
p^{\epsilon} \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{i}^{\circ} \subseteq \tilde{\mathcal{E}}_{i}, \quad F\left(\tilde{\mathcal{E}}_{i}\right) \subseteq \tilde{\mathcal{E}}_{3-i}, \quad \text { and } \quad V\left(\tilde{\mathcal{E}}_{i}\right) \subseteq \tilde{\mathcal{E}}_{3-i}, \quad \text { for } i=1,2
$$

and the following condition for $i=1,2$ :

$$
\begin{align*}
& \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{i}^{\circ} / \tilde{\mathcal{E}}_{i} \simeq\left(W\left(\overline{\mathbb{F}}_{p}\right) / p^{\epsilon}\right)^{\oplus j_{1, i}} \oplus\left(W\left(\overline{\mathbb{F}}_{p}\right) / p^{\epsilon-1}\right)^{\oplus\left(j_{2, i}-j_{1, i}\right)} \oplus \cdots \\
& \cdots \oplus\left(W\left(\overline{\mathbb{F}}_{p}\right) / p\right)^{\oplus\left(j_{\epsilon, i}-j_{\epsilon-1, i}\right)} \tag{7.1.2}
\end{align*}
$$

We refer to the toy model discussed in Example 7.3 for the motivation of this condition. For technical reasons, we will not prove the set $Z_{j}$ is the set of $\overline{\mathbb{F}}_{p}$-points of a closed subscheme of $\mathrm{Sh}_{r, s}$; instead we prove that a closely related subset of $Z_{j}$ is. See Remark 7.5.

Applying Proposition 3.2 with $m=\delta$, the submodules $\tilde{\mathcal{E}}_{1}$ and $\tilde{\mathcal{E}}_{2}$ give rise to a polarized abelian variety $\left(\mathcal{A}_{z}^{\prime}, \lambda_{z}^{\prime}\right)$ over $z$ with an $\mathcal{O}_{D}$-action and an $\mathcal{O}_{D}$-equivariant
isogeny $\mathcal{A}_{z}^{\prime} \rightarrow \mathcal{A}_{z}$. Moreover, by (3.2.2), we have

$$
\begin{aligned}
\operatorname{dim} \omega_{\mathcal{A}_{z}^{\vee \vee} / \overline{\mathbb{F}}_{p}, 1}^{\circ} & =\operatorname{dim} \omega_{\mathcal{A}_{z}^{\vee} / \overline{\mathbb{F}}_{p}, 1}^{\circ}+\sum_{\alpha=0}^{\epsilon-1}\left((\epsilon-\alpha)\left(j_{\alpha+1,1}-j_{\alpha, 1}\right)-(\epsilon-\alpha)\left(j_{\alpha+1,2}-j_{\alpha, 2}\right)\right) \\
& =r-\delta
\end{aligned}
$$

and similarly $\operatorname{dim} \omega_{\mathcal{A}_{z}^{\sim} / \bar{F}_{p}, 2}^{\circ}=s+\delta$. So $\mathcal{A}_{z}^{\prime}$ satisfies the moduli problem for $\mathrm{Sh}_{r-\delta, s+\delta}$; this suggests a geometric relationship between $Z_{j}$ and $\mathrm{Sh}_{r-\delta, s+\delta}$ that we make precise in Definition 7.4.

We make an immediate remark that when $\delta=r$, the abelian variety $\mathcal{A}_{z}$ coming from a point $z$ of $Z_{j}$ is isogenous to an abelian variety $\mathcal{A}_{z}^{\prime}$ that is a moduli object for the Shimura variety $\mathrm{Sh}_{0, n}$. Thus both $\mathcal{A}_{z}^{\prime}$ and $\mathcal{A}_{z}$ are supersingular. So every $Z_{j}$ is contained in the supersingular locus of $\mathrm{Sh}_{r, s}$. In fact, we shall show in Theorem 7.8 that the supersingular locus of $\mathrm{Sh}_{r, s}$ is exactly the union of these $Z_{j}$.
7.2. Towards a moduli interpretation. We need to reinterpret in a more geometric manner the Dieudonné-theoretic condition defining $Z_{j}$. For $\alpha=0, \ldots, \epsilon$, we define submodules

$$
\tilde{\mathcal{E}}_{\alpha, 1}:=\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha}} \tilde{\mathcal{E}}_{1} \quad \text { and } \quad \tilde{\mathcal{E}}_{\alpha, 2}:=\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha}} \tilde{\mathcal{E}}_{2}
$$

of $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}$ and $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ}$. They are easily seen to satisfy condition (3.2.1) with $m=\alpha$. Thus, Proposition 3.2 generates a polarized abelian variety $\left(A_{\alpha}, \lambda_{\alpha}\right)$ with $\mathcal{O}_{D}$-action and an $\mathcal{O}_{D}$-equivariant isogeny $A_{\alpha} \rightarrow \mathcal{A}_{z}$, where

$$
\begin{align*}
& r_{\alpha}:=\operatorname{dim} \omega_{A_{\alpha}^{\vee} / \mathbb{F}_{p}, 1}^{\circ}=r-\sum_{\alpha^{\prime}=1}^{\alpha}\left(j_{\alpha^{\prime}, 2}-j_{\alpha^{\prime}, 1}\right) \quad \text { and }  \tag{7.2.1}\\
& s_{\alpha}:=\operatorname{dim} \omega_{A_{\alpha}^{\vee} / \mathbb{F}_{p}, 2}^{\circ}=n-\operatorname{dim} \omega_{A_{\alpha}^{\vee} / \overline{\mathbb{F}}_{p}, 1}^{\circ}
\end{align*}
$$

by the formula (3.2.2). In particular $r_{0}=r, s_{0}=s, r_{\epsilon}=r-\delta$ and $s_{\epsilon}=s+\delta$.
In fact, applying Proposition 3.2 (with $m=1$ ) to the sequence of inclusions

$$
\tilde{\mathcal{E}}_{i}=\tilde{\mathcal{E}}_{\epsilon, i} \subset \tilde{\mathcal{E}}_{\epsilon-1, i} \subset \cdots \subset \tilde{\mathcal{E}}_{0, i}=\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{i}^{\circ}
$$

we obtain a sequence of isogenies (each with $p$-torsion kernels):

$$
\begin{equation*}
\mathcal{A}_{z}^{\prime}=A_{\epsilon} \xrightarrow{\phi_{\epsilon}} A_{\epsilon-1} \xrightarrow{\phi_{\epsilon-1}} \ldots \xrightarrow{\phi_{1}} A_{0}=\mathcal{A}_{z} . \tag{7.2.2}
\end{equation*}
$$

We have $\operatorname{ker} \phi_{\alpha} \subseteq A_{\alpha}[p]$, so that there exists a unique isogeny $\psi_{\alpha}: A_{\alpha-1} \rightarrow A_{\alpha}$ such that $\psi_{\alpha} \phi_{\alpha}=p \cdot \mathrm{id}_{A_{\alpha}}$ and $\phi_{\alpha} \psi_{\alpha}=p \cdot \mathrm{id}_{A_{\alpha-1}}$.

For each $\alpha$, the cokernel of the induced map on cohomology

$$
\begin{aligned}
& \phi_{\alpha, *, i}: H_{1}^{\mathrm{dR}}\left(A_{\alpha} / \overline{\mathbb{F}}_{p}\right)_{i}^{\circ} \rightarrow H_{1}^{\mathrm{dR}}\left(A_{\alpha-1} / \overline{\mathbb{F}}_{p}\right)_{i}^{\circ} \\
&\left(\operatorname{resp.} \psi_{\alpha, *, i}: H_{1}^{\mathrm{dR}}\left(A_{\alpha-1} / \overline{\mathbb{F}}_{p}\right)_{i}^{\circ} \rightarrow H_{1}^{\mathrm{dR}}\left(A_{\alpha} / \overline{\mathbb{F}}_{p}\right)_{i}^{\circ}\right)
\end{aligned}
$$

is canonically isomorphic to $\tilde{\mathcal{E}}_{\alpha-1, i} / \tilde{\mathcal{E}}_{\alpha, i}$ (resp. $\tilde{\mathcal{E}}_{\alpha, i} / p \tilde{\mathcal{E}}_{\alpha-1, i}$ ), which has dimension $j_{\alpha, i}$ (resp. $n-j_{\alpha, i}$ ) over $\overline{\mathbb{F}}_{p}$ by a straightforward computation using (7.1.2).

The upshot is that all these numeric information of the chain of isogenies (7.2.2) can be used to reconstruct $\tilde{\mathcal{E}}_{i}$ inside $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{i}^{\circ}$. This idea will be made precise after this important example.

Example 7.3. We give a good toy model for the isogenies of Dieudonné modules. This is the inspiration of the construction of this section. We start with the Dieudonné module $\tilde{\mathcal{D}}\left(A_{\epsilon}\right)_{1}^{\circ}=\bigoplus_{i=1}^{n} W\left(\overline{\mathbb{F}}_{p}\right) \boldsymbol{e}_{j}$ and $\tilde{\mathcal{D}}\left(A_{\epsilon}\right)_{2}^{\circ}=\bigoplus_{j=1}^{n} W\left(\overline{\mathbb{F}}_{p}\right) \boldsymbol{f}_{j}$. The maps $V_{1}$ : $\tilde{\mathcal{D}}\left(A_{\epsilon}\right)_{1}^{\circ} \rightarrow \tilde{\mathcal{D}}\left(A_{\epsilon}\right)_{2}^{\circ}$ and $V_{2}: \tilde{\mathcal{D}}\left(A_{\epsilon}\right)_{2}^{\circ} \rightarrow \tilde{\mathcal{D}}\left(A_{\epsilon}\right)_{1}^{\circ}$, with respect to the given bases, are given by the diagonal matrices

$$
\operatorname{Diag}(\underbrace{1, \ldots, 1}_{s+\delta}, \underbrace{p, \ldots, p}_{r-\delta}) \text { and } \operatorname{Diag}(\underbrace{1, \ldots, 1}_{r-\delta}, \underbrace{p, \ldots, p}_{s+\delta}) \text {, }
$$

respectively. Using the isogenies $\phi_{\alpha}$ we may naturally identify $\tilde{\mathcal{D}}\left(A_{\alpha}\right)_{i}^{\circ}$ as lattices in $\tilde{\mathcal{D}}\left(A_{\epsilon}\right)_{i}^{\circ}[1 / p]$ with induced Frobenius and Verschiebung morphisms. For our toy model, we choose

$$
\begin{array}{r}
\tilde{\mathcal{D}}\left(A_{\alpha}\right)_{1}^{\circ}=\operatorname{Span}_{W\left(\mathbb{F}_{p}\right)}\left\{\frac{1}{p^{\epsilon-\alpha}} \boldsymbol{e}_{1}, \ldots, \frac{1}{p^{\epsilon-\alpha}} \boldsymbol{e}_{j_{\alpha+1,1}}, \frac{1}{p^{\epsilon-\alpha-1}} \boldsymbol{e}_{j_{\alpha+1,1}+1}, \ldots, \frac{1}{p^{\epsilon-\alpha-1}} \boldsymbol{e}_{j_{\alpha+2,1}},\right. \\
\left.\frac{1}{p^{\epsilon-\alpha-2}} \boldsymbol{e}_{j_{\alpha+2,1}+1}, \ldots, \frac{1}{p} \boldsymbol{e}_{j_{\epsilon, 1}}, \boldsymbol{e}_{j_{\epsilon, 1}+1}, \ldots, \boldsymbol{e}_{n}\right\} ; \\
\tilde{\mathcal{D}}\left(A_{\alpha}\right)_{2}^{\circ}=\operatorname{Span}_{W\left(\bar{F}_{p}\right)}\left\{\frac{1}{p^{\epsilon-\alpha}} \boldsymbol{f}_{1}, \ldots, \frac{1}{p^{\epsilon-\alpha}} \boldsymbol{f}_{j_{\alpha+1,2}}, \frac{1}{p^{\epsilon-\alpha-1}} \boldsymbol{f}_{j_{\alpha+1,2}+1}, \ldots, \frac{1}{p^{\epsilon-\alpha-1}} \boldsymbol{f}_{j_{\alpha+2,2}},\right. \\
\left.\frac{1}{p^{\epsilon-\alpha-2}} \boldsymbol{f}_{j_{\alpha+2,2}+1}, \ldots, \frac{1}{p} \boldsymbol{f}_{j_{\epsilon, 2}}, \boldsymbol{f}_{j_{\epsilon, 2}+1}, \ldots, \boldsymbol{f}_{n}\right\} .
\end{array}
$$

In particular, the Verschiebung $V_{1}: \tilde{\mathcal{D}}\left(A_{\alpha}\right)_{1}^{\circ} \rightarrow \tilde{\mathcal{D}}\left(A_{\alpha}\right)_{2}^{\circ}$ with respect to the bases above is given by

$$
\operatorname{Diag}(\underbrace{1, \ldots, 1}_{j_{\alpha+1,1}}, * * * \cdots * * *, \underbrace{p, \ldots, p}_{r-\delta})
$$

where the $* * *$ part is $p$ if the place is in $\left[j_{\alpha^{\prime}, 1}+1, j_{\alpha^{\prime}, 2}\right]$ for some $\alpha^{\prime} \geq \alpha$, and is 1 otherwise. Similarly, the Verschiebung $V_{2}: \tilde{\mathcal{D}}\left(A_{0}\right)_{2}^{\circ} \rightarrow \tilde{\mathcal{D}}\left(A_{0}\right)_{1}^{\circ}$ with respect to the bases above is given by

$$
\operatorname{Diag}(\underbrace{1, \ldots, 1}_{r-\delta}, \underbrace{p, \ldots, p}_{j_{\alpha+1,1}-r+\delta}, * * * \cdots * * *, \underbrace{p, \ldots, p}_{n-j_{\alpha_{\epsilon, 2}}})
$$

where the $* * *$ part is 1 if the place is in $\left[j_{\alpha^{\prime}, 1}+1, j_{\alpha^{\prime}, 2}\right]$ for some $\alpha^{\prime} \geq \alpha$, and is $p$ otherwise.

So the sheaf of differentials is given by

$$
\begin{aligned}
& \omega_{A_{\alpha}^{\vee} / \mathbb{F}_{p}, 1}^{\circ}= \operatorname{span}_{\overline{\mathbb{F}}_{p}}\left\{\frac{1}{p^{\epsilon-\alpha}} \boldsymbol{e}_{1}, \ldots, \frac{1}{p^{\epsilon-\alpha}} \boldsymbol{e}_{r-\delta}, \frac{1}{p^{\epsilon-\alpha}} \boldsymbol{e}_{j_{\alpha, 1}+1}, \ldots, \frac{1}{p^{\epsilon-\alpha}} \boldsymbol{e}_{j_{\alpha, 2}},\right. \\
&\left.\frac{1}{p^{\epsilon-\alpha-1}} \boldsymbol{e}_{j_{\alpha+1,1}+1}, \ldots, \frac{1}{p} \boldsymbol{e}_{j_{\epsilon-1,2}}, \boldsymbol{e}_{j_{\epsilon, 1}+1}, \ldots, \boldsymbol{e}_{j_{\epsilon, 2}}\right\} \\
& \omega_{A_{\alpha}^{\vee} / \mathbb{F}_{p}, 2}^{\circ}=\operatorname{Span}_{\overline{\mathbb{F}}_{p}}\left\{\frac{1}{p^{\epsilon-\alpha}} \boldsymbol{f}_{1}, \ldots, \frac{1}{p^{\epsilon-\alpha}} \boldsymbol{f}_{j_{\alpha+1,1}}, \frac{1}{p^{\epsilon-\alpha-1}} \boldsymbol{f}_{j_{\alpha+1,2}+1}, \ldots, \frac{1}{p^{\epsilon-\alpha-1}} \boldsymbol{f}_{j_{\alpha+2,1}}\right. \\
&\left.\frac{1}{p^{\epsilon-\alpha-2}} \boldsymbol{f}_{j_{\alpha+2,1}}, \ldots, \frac{1}{p} \boldsymbol{f}_{j_{\epsilon, 1}}, \boldsymbol{f}_{j_{\epsilon, 2}+1}, \ldots, \boldsymbol{f}_{s+\delta-1}\right\}
\end{aligned}
$$

Definition 7.4. Let $\boldsymbol{j}$ be as above. Define the numbers $j_{\alpha, i}$ as in (7.1.1) and the numbers $r_{\alpha}, s_{\alpha}$ as in (7.2.1). Let $\underline{Y}_{\boldsymbol{j}}$ be the functor taking a locally noetherian $\mathbb{F}_{p^{2}}$-scheme $S$ to the set of isomorphism classes of tuples

$$
\begin{equation*}
\left(A_{0}, \ldots, A_{\epsilon}, \lambda_{0}, \ldots, \lambda_{\epsilon}, \eta_{0}, \ldots, \eta_{\epsilon}, \phi_{1}, \ldots, \phi_{\epsilon}, \psi_{1}, \ldots, \psi_{\epsilon}\right) \tag{7.4.1}
\end{equation*}
$$

such that:
(1) for each $\alpha,\left(A_{\alpha}, \lambda_{\alpha}, \eta_{\alpha}\right)$ is an $S$-point of $\mathrm{Sh}_{r_{\alpha}, s_{\alpha}}$;
(2) for each $\alpha, \phi_{\alpha}$ is an $\mathcal{O}_{D}$-isogeny $A_{\alpha} \rightarrow A_{\alpha-1}$, with kernel contained in $A_{\alpha}[p]$, which is compatible with the polarizations in the sense that $p \lambda_{\alpha}=\phi_{\alpha}^{\vee} \circ \lambda_{\alpha-1} \circ \phi_{\alpha}$ and with the tame level structures in the sense that $\phi_{\alpha} \circ \eta_{\alpha}=\eta_{\alpha-1}$;
(3) $\psi_{\alpha}$ is the isogeny $A_{\alpha-1} \rightarrow A_{\alpha}$ such that $\phi_{\alpha} \psi_{\alpha}=p \cdot \mathrm{id}_{A_{\alpha}}$ and $\psi_{\alpha} \phi_{\alpha}=p \cdot \mathrm{id}_{A_{\alpha-1}}$;
(4) the cokernel of the induced map $\phi_{\alpha, *, i}^{\mathrm{dR}}: H_{1}^{\mathrm{dR}}\left(A_{\alpha} / S\right)_{i}^{\circ} \rightarrow H_{1}^{\mathrm{dR}}\left(A_{\alpha-1} / S\right)_{i}^{\circ}$ is a locally free $\mathcal{O}_{S}$-module of rank $j_{\alpha, i}$ for each $\alpha$ and $i=1,2$;
(5) the cokernel of the induced map $\psi_{\alpha, *, i}^{\mathrm{dR}}: H_{1}^{\mathrm{dR}}\left(A_{\alpha-1} / S\right)_{i}^{\circ} \rightarrow H_{1}^{\mathrm{dR}}\left(A_{\alpha} / S\right)_{i}^{\circ}$ is a locally free $\mathcal{O}_{S}$-module of rank $n-j_{\alpha, i}$ for each $\alpha$ and $i=1,2 ;{ }^{18}$
(6) for each $\alpha, \operatorname{Ker}\left(\phi_{\alpha, *, 2}^{\mathrm{dR}}\right)$ is contained in $\omega_{A_{\alpha}^{\vee} / S, 2}^{\circ}$;
(7) for each $\alpha$, the ( $r_{\alpha-1}-r_{\alpha}+r_{\epsilon}+1$ )-st Fitting ideal of the cokernel of $\phi_{\alpha, *, 1}^{\mathrm{dR}}$ : $\omega_{A_{\alpha}^{\vee} / S, 1}^{\circ} \rightarrow \omega_{A_{\alpha-1}^{\vee} / S, 1}^{\circ}$ is zero, or equivalently, Zariski locally on $S$, if we represent the map $\phi_{\alpha, *, 1}^{\mathrm{dR}}: \omega_{A_{\alpha}^{\vee} / S, 1}^{\circ} \rightarrow \omega_{A_{\alpha-1}^{\vee} / S, 1}^{\circ}$ by an $r_{\alpha-1} \times r_{\alpha}$-matrix (after choosing local bases), then all $\left(r_{\alpha}-r_{\epsilon} \stackrel{\alpha-1}{+}\right) \times\left(r_{\alpha}-r_{\epsilon}+1\right)$-minors vanish.
(8) the $\left(r_{\alpha}-r_{\epsilon}+1\right)$-st Fitting ideal of the cokernel of $\psi_{\alpha,,, 1}^{\mathrm{dR}}: \omega_{A_{\alpha-1}^{\vee} / S, 1}^{\circ} \rightarrow \omega_{A_{\alpha}^{\vee} / S, 1}^{\circ}$ is zero for each $\alpha$.

Note that conditions (6)-(8) are all closed conditions. So the moduli problem $\underline{Y}_{j}$ is represented by a proper scheme $Y_{j}$ of finite type over $\mathbb{F}_{p^{2}}$. The moduli space $Y_{j}$ admits natural maps to $\mathrm{Sh}_{r, s}$ and $\mathrm{Sh}_{r-\delta, s+\delta}$ by sending the tuple (7.4.1) to

[^13]$\left(A_{0}, \lambda_{0}, \eta_{0}\right)$ and $\left(A_{\epsilon}, \lambda_{\epsilon}, \eta_{\epsilon}\right)$, respectively.


We also point out that conditions (2) and (3) together imply that, for each $\alpha$ and $i=1$, 2, we have $\operatorname{Im}\left(\psi_{\alpha, *, i}^{\mathrm{dR}}\right)=\operatorname{Ker}\left(\phi_{\alpha, *, i}^{\mathrm{dR}}\right)$ and $\operatorname{Im}\left(\phi_{\alpha, *, i}^{\mathrm{dR}}\right)=\operatorname{Ker}\left(\psi_{\alpha, *, i}^{\mathrm{dR}}\right)$. We shall freely use this property later.
Remark 7.5. Conditions (6)-(8) in Definition 7.4 are satisfied by the toy model in Example 7.3. They did not appear in moduli problem in Section 4.2 because they trivially hold in that case. The purpose of keeping these conditions in the moduli problem and carefully formulating them is so that the moduli space may hope to have the correct irreducible components. We think the picture is the following: $Z_{j}$ is probably or at least heuristically the set of $\overline{\mathbb{F}}_{p}$-points of a closed subscheme of $\mathrm{Sh}_{r, s}$. But this scheme has many irreducible components, which may have overlaps with other $Z_{j^{\prime}}$. Conditions (6)-(8) will help select one irreducible component that is "special" for $\boldsymbol{j}$. When taking the union of all images of the $Y_{\boldsymbol{j}}$, we should still get the union of the $Z_{j}$. This is verified in the case of supersingular locus (i.e., $r=\delta$ ) in Theorem 7.8.

Notation 7.6. Let $Y_{j}$ as above. It will be convenient to introduce some dummy notation:

- $\phi_{0}$ is the identity map on $A_{0}$;
- $\psi_{\epsilon}$ is the identity map on $A_{\epsilon}$.

We use $Y_{j}^{\circ}$ to denote the open subscheme of $Y_{j}$ representing the functor that takes a locally noetherian $\mathbb{F}_{p^{2}}$-scheme $S$ to the subset of isomorphism classes of tuples

$$
\left(A_{0}, \ldots, A_{\epsilon}, \lambda_{0}, \ldots, \lambda_{\epsilon}, \eta_{0}, \ldots, \eta_{\epsilon}, \phi_{1}, \ldots, \phi_{\epsilon}, \psi_{1}, \ldots, \psi_{\epsilon}\right)
$$

of $Y_{j}(S)$ such that
(i) for each $\alpha=1, \ldots, \epsilon$, the $\operatorname{sum} \phi_{\alpha, *, 2}\left(\omega_{A_{\alpha}^{\vee} / S, 2}^{\circ}\right)+\operatorname{Ker}\left(\phi_{\alpha-1, *, 2}^{\mathrm{dR}}\right)$ is an $\mathcal{O}_{S^{-}}$ subbundle of $H_{1}^{\mathrm{dR}}\left(A_{\alpha-1} / S\right)_{2}^{\circ}$ of rank $\operatorname{rank} \omega_{A_{\alpha}^{\vee} / S, 2}-\operatorname{rank} \operatorname{Ker}\left(\phi_{\alpha, *, 2}^{\mathrm{dR}}\right)+\operatorname{rank} \operatorname{Ker}\left(\phi_{\alpha-1, *, 2}^{\mathrm{dR}}\right)=s_{\alpha}-j_{\alpha, 2}+j_{\alpha-1,2}$,
(ii) for each $\alpha=1, \ldots, \epsilon, \operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\mathrm{dR}}\right)+\operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\mathrm{dR}}\right)$ is an $\mathcal{O}_{S^{\text {-subbundle of rank }} \text {. }}$

$$
\operatorname{rank} \operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\mathrm{dR}}\right)+\operatorname{rank} \operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\mathrm{dR}}\right)=j_{\alpha, 1}+\left(n-j_{\alpha+1,1}\right),
$$

(iii) for each $\alpha$, the cokernel of $\phi_{\alpha, *, 1}^{\mathrm{dR}}: \omega_{A_{\alpha}^{\vee} / S, 1}^{\circ} \rightarrow \omega_{A_{\alpha-1}^{\vee} / S, 1}^{\circ}$ is a locally free $\mathcal{O}_{S}$-module of rank $r_{\alpha-1}-\left(r_{\alpha}-r_{\epsilon}\right)$,
(iv) for each $\alpha$, the cokernel of $\psi_{\alpha,,, 1}^{\mathrm{dR}}: \omega_{A_{\alpha-1}^{\vee} / S, 1}^{\circ} \rightarrow \omega_{A_{\alpha}^{\vee} / S, 1}^{\circ}$ is a locally free $\mathcal{O}_{S}$-module of rank $r_{\alpha}-r_{\epsilon}$.

We note that the ranks in conditions (i) and (ii) are maximal possible and the ranks in conditions (iii) and (iv) are minimal possible, under the conditions in Definition 7.4. So $Y_{j}^{\circ}$ is an open subscheme of $Y_{j}$.

We point out an additional benefit of having conditions (ii)-(iv). By (iii), $\omega_{A_{\alpha}^{\vee} / S, 1}^{\circ} \cap \operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\mathrm{dR}}\right)$ is an $\mathcal{O}_{S}$-subbundle of $\omega_{A_{\alpha}^{\vee} / S, 1}^{\circ}$ of rank $r_{\epsilon}$, for $\alpha=1, \ldots, \epsilon$; by (iv), $\omega_{A_{\alpha}^{\vee} / S, 1}^{\circ} \cap \operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\mathrm{dR}}\right)$ is an $\mathcal{O}_{S}$-subbundle of $\omega_{A_{\alpha}^{\vee} / S, 1}^{\circ}$ of rank $r_{\alpha}-r_{\epsilon}$, for $\alpha=0, \ldots, \epsilon-1$. Combining these two rank estimates and condition (ii) which implies that $\operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\mathrm{dR}}\right)$ and $\operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\mathrm{dR}}\right)$ are disjoint subbundles, we arrive at a direct sum decomposition

$$
\begin{equation*}
\omega_{A_{\alpha}^{\vee} / S, 1}^{\circ}=\left(\omega_{A_{\alpha}^{\vee} / S, 1}^{\circ} \cap \operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\mathrm{dR}}\right)\right) \oplus\left(\omega_{A_{\alpha}^{\vee} / S, 1}^{\circ} \cap \operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\mathrm{dR}}\right)\right), \tag{7.6.1}
\end{equation*}
$$

for $\alpha=1, \ldots, \epsilon-1$; and we know that $\omega_{A_{0}^{\vee} / S, 1}^{\circ} \cap \operatorname{Ker}\left(\psi_{1, *, 1}^{\mathrm{dR}}\right)$ has rank $r_{0}-r_{\epsilon}=\delta$ and $\omega_{A_{\epsilon} / S, 1}^{\circ} \subseteq \operatorname{Ker}\left(\phi_{\epsilon, * 1}^{\mathrm{dR}}\right)$.

We shall show below in Theorem 7.7 that $Y_{j}^{\circ}$ is smooth. Unfortunately, we do not know how to prove the nonemptiness of $Y_{j}^{\circ}$, nor do we know if some $Y_{j}$ is completely contained in some other $Y_{j}$; but the fact that the Dieudonné modules in Example 7.3 satisfy conditions (i)-(iv) above is good evidence for this nonemptiness. Of course, if one can compute the intersection matrix in the sense of Theorem 6.7 and calculate the determinant, one can then probably show that these $Y_{j}$ are essentially different. But the difficulties of this computation lie in understanding the singularities at $Y_{\boldsymbol{j}} \backslash Y_{\boldsymbol{j}}^{\circ}$, which seems to be very combinatorially involved.
Theorem 7.7. Each $Y_{j}^{\circ}$ is smooth of dimension $r s+(r-\delta)(s+\delta)$ (if not empty). Proof. Let $\hat{R}$ be a noetherian $\mathbb{F}_{p^{2}}$-algebra and $\hat{I} \subset \hat{R}$ an ideal such that $\hat{I}^{2}=0$. Put $R=\hat{R} / \hat{I}$. Say we want to lift an $R$-point

$$
\left(A_{0}, \ldots, A_{\epsilon}, \lambda_{0}, \ldots, \lambda_{\epsilon}, \eta_{0}, \ldots, \eta_{\epsilon}, \phi_{1}, \ldots, \phi_{\epsilon}, \psi_{1}, \ldots, \psi_{\epsilon}\right)
$$

of $Y_{j}^{\circ}$ an $\hat{R}$-point and we try to compute the corresponding tangent space. By SerreTate and Grothendieck-Messing deformation theory we recalled in Theorem 3.4, it is enough to lift, for $i=1,2$ and each $\alpha=0, \ldots, \epsilon$, the differentials $\omega_{A_{\alpha}^{\vee} / R, i}^{\circ} \subseteq$ $H_{1}^{\mathrm{dR}}\left(A_{\alpha} / R\right)_{i}^{\circ}$ to a subbundle $\hat{\omega}_{\alpha, i} \subseteq H_{1}^{\text {cris }}\left(A_{\alpha} / \hat{R}\right)_{i}^{\circ}$ such that
(a) $\phi_{\alpha, *, i}^{\mathrm{cris}}\left(\hat{\omega}_{\alpha, i}\right) \subseteq \hat{\omega}_{\alpha-1, i}$ and $\psi_{\alpha, *, i}^{\mathrm{cris}}\left(\hat{\omega}_{\alpha-1, i}\right) \subseteq \hat{\omega}_{\alpha, i}$ (so that both $\phi_{\alpha}$ and $\psi_{\alpha}$ are lifted, which would automatically imply $\left.\operatorname{Ker}\left(\phi_{\alpha}\right) \in A_{\alpha}[p]\right)$,
(b) $\hat{\omega}_{\alpha, 2} \supseteq \operatorname{Ker}\left(\phi_{\alpha, *, 2}^{\text {cris }}\right)$, and
(c) the $\hat{R}$-modules $\hat{\omega}_{\alpha-1,1} / \phi_{\alpha, *, 1}^{\text {cris }}\left(\hat{\omega}_{\alpha, 1}\right)$ and $\hat{\omega}_{\alpha, 1} / \psi_{\alpha, *, 1}^{\text {cris }}\left(\hat{\omega}_{\alpha-1,1}\right)$ are flat and of rank $r_{\alpha-1}-\left(r_{\alpha}-r_{\epsilon}\right)$ and $r_{\alpha}-r_{\epsilon}$, respectively.
We shall see that condition (i) of Notation 7.6 is automatic. Also, condition (ii) already holds: since $H_{1}^{\text {cris }}\left(A_{\alpha} / \hat{R}\right)_{1}^{\circ} /\left(\operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\text {cris }}\right)+\operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\text {cris }}\right)\right)$ is locally generated by $j_{\alpha+1,1}-j_{\alpha, 1}$ elements after modulo $\hat{I}$, it is so prior to modulo $\hat{I}$
by Nakayama's lemma. Note that rank of $\operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\text {cris }}\right)$ and $\operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\text {cris }}\right)$ and the number of the generators of the quotient above add up to exactly $n$; it follows that $\operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\text {cris }}\right)+\operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\text {cris }}\right)$ is a direct sum and the sum is a subbundle of $H_{1}^{\text {cris }}\left(A_{\alpha} / \hat{R}\right)_{1}^{\circ}$.

We separate the discussion of lifts at $q_{1}$ and $q_{2}$, and show that the tangent space $T_{Y_{j}^{\circ}}$ is isomorphic to $T_{1} \oplus T_{2}$ for the contributions $T_{1}$ and $T_{2}$ from the two places. We first look at $q_{2}$, as it is easier. Note that condition (b) $\hat{\omega}_{\alpha, 2} \supseteq \operatorname{Ker}\left(\phi_{\alpha, *, 2}^{\text {cris }}\right)=\operatorname{Im}\left(\psi_{\alpha, *, 2}^{\text {cris }}\right)$ automatically implies that $\psi_{\alpha, *, 2}^{\text {cris }}\left(\hat{\omega}_{\alpha-1,2}\right) \subseteq \hat{\omega}_{\alpha, 2}$; so we can proceed as follows:
Step 0: First lift $\omega_{A_{\epsilon} / R, 2}^{\circ}$ to a subbundle $\hat{\omega}_{\epsilon, 2}$ of $H_{1}^{\text {cris }}\left(A_{\epsilon} / \hat{R}\right)_{2}^{\circ}$ so that it contains $\operatorname{Ker}\left(\phi_{\epsilon, * 2}^{\text {cris }}\right)$,
Step 1: then lift $\omega_{A_{\epsilon-1}^{\vee} / R, 2}^{\circ}$ to a subbundle $\hat{\omega}_{\epsilon-1,2}$ of $H_{1}^{\text {cris }}\left(A_{\epsilon-1} / \hat{R}\right)_{2}^{\circ}$ so that it contains $\phi_{\epsilon, * 2}^{\text {cris }}\left(\hat{\omega}_{\epsilon, 2}\right)+\operatorname{Ker}\left(\phi_{\epsilon-1, *, 2}^{\text {cris }}\right)$,
Step(s) $\alpha$ : then lift $\omega_{A_{\epsilon-\alpha}^{\vee} / R, 2}^{\circ}$ to a subbundle $\hat{\omega}_{\epsilon-\alpha, 2}$ of $H_{1}^{\text {cris }}\left(A_{\epsilon-\alpha} / \hat{R}\right)_{2}^{\circ}$ so that it contains $\phi_{\epsilon-\alpha+1, *, 2}^{\text {cris }}\left(\hat{\omega}_{\epsilon-\alpha+1,2}^{\circ}\right)+\operatorname{Ker}\left(\phi_{\epsilon-\alpha, *, 2}^{\text {cris }}\right)$,
Step $\boldsymbol{\epsilon}$ : finally lift $\omega_{A_{0}^{\vee} / R, 2}^{\circ}$ to a subbundle $\hat{\omega}_{0,2}$ of $H_{1}^{\text {cris }}\left(A_{0} / \hat{R}\right)_{2}^{\circ}$ so that it contains $\phi_{1, *, 2}^{\text {cris }}\left(\hat{\omega}_{1,2}\right)$.

At Step 0, the choices form a torsor for the group

$$
\operatorname{Hom}_{R}\left(\omega_{A_{\epsilon}^{\vee} / R, 2}^{\circ} / \operatorname{Ker}\left(\phi_{\epsilon, *, 2}^{\mathrm{dR}}\right), \operatorname{Lie}_{A_{\epsilon} / R, 2}^{\circ}\right) \otimes_{R} \hat{I}
$$

the Hom space is a locally free $R$-module of rank $\left(s_{\epsilon}-j_{\epsilon, 2}\right) r_{\epsilon}$.
At Step $\alpha=1, \ldots, \epsilon$, we observe that condition (i) of the moduli problem $\underline{Y}_{j}^{\circ} \mathrm{im}$ plies $\phi_{\epsilon-\alpha+1, *, 2}\left(\omega_{A_{\epsilon-\alpha+1}^{\vee} / R, 2}^{\circ}\right)+\operatorname{Ker}\left(\phi_{\epsilon-\alpha, *, 2}^{\mathrm{dR}}\right)$ is an $R$-subbundle of $H_{1}^{\mathrm{dR}}\left(A_{\epsilon-\alpha} / R\right)_{2}^{\circ}$ of rank

$$
\begin{equation*}
s_{\epsilon-\alpha+1}-j_{\epsilon-\alpha+1,2}+j_{\epsilon-\alpha, 2}=s_{\epsilon-\alpha}+\left(j_{\epsilon-\alpha+1,1}-j_{\epsilon-\alpha, 2}\right) \quad \text { if } \alpha=1, \ldots, \epsilon-1 \tag{7.7.1}
\end{equation*}
$$

and of rank $s_{1}-j_{1,2}$ if $\alpha=\epsilon$. So $\phi_{\epsilon-\alpha+1, *, 2}^{\text {cris }}\left(\hat{\omega}_{\epsilon-\alpha+1,2}\right)+\operatorname{Ker}\left(\phi_{\epsilon-\alpha, *, 2}^{\text {cris }}\right)$ is an $\hat{R}-$ subbundle of $H_{1}^{\text {cris }}\left(A_{\epsilon-\alpha} / \hat{R}\right)_{2}^{\circ}$ of the same rank. The choices of the lifts $\hat{\omega}_{\epsilon-\alpha, 2}$ form a torsor for the group

$$
\operatorname{Hom}_{R}\left(\omega_{A_{\epsilon-\alpha}^{\vee} / R, 2}^{\circ} /\left(\phi_{\epsilon-\alpha+1, *, 2}\left(\omega_{A_{\epsilon-\alpha+1} / R, 2}^{\circ}\right)+\operatorname{Ker}\left(\phi_{\epsilon-\alpha, *, 2}^{\mathrm{dR}}\right)\right), \operatorname{Lie}_{A_{\epsilon-\alpha} / R, 2}^{\circ}\right) \otimes_{R} \hat{I}
$$

By (7.7.1), this Hom space is a locally free $R$-module of $\operatorname{rank}\left(j_{\epsilon-\alpha+1,1}-j_{\epsilon-\alpha, 2}\right) r_{\epsilon-\alpha}$ if $\alpha=1, \ldots, \epsilon-1$ and of $\operatorname{rank}\left(s_{0}-\left(s_{1}-j_{1,2}\right)\right) r_{0}$ if $\alpha=\epsilon$. This implies that the contribution $T_{2}$ to the tangent space $T_{Y_{j}^{\circ}}$ at $q_{2}$ admits a filtration such that the subquotients are

$$
\mathcal{H o m}\left(\omega_{\mathcal{A}_{\epsilon-\alpha}^{\vee}, 2}^{\circ} /\left(\phi_{\epsilon-\alpha+1, *, 2}\left(\omega_{\mathcal{A}_{\epsilon-\alpha+1}, 2}^{\circ}\right)+\operatorname{Ker}\left(\phi_{\epsilon-\alpha, *, 2}^{\mathrm{dR}}\right)\right), \operatorname{Lie}_{\mathcal{A}_{\epsilon-\alpha}, 2}^{\circ}\right)
$$

where the $\mathcal{A}_{\epsilon-\alpha}$ are the universal abelian varieties and $\phi_{\epsilon+1, *, 2}\left(\omega_{\mathcal{A}_{\epsilon+1}^{\vee}, 2}^{\circ}\right)$ is interpreted as zero. In particular, $T_{2}$ is a locally free sheaf on $Y_{j}^{\circ}$ of rank

$$
\begin{align*}
\left(s_{\epsilon}-j_{\epsilon, 2}\right) r_{\epsilon}+\left(s_{0}-\left(s_{1}-\right.\right. & \left.\left.j_{1,2}\right)\right) r_{0}+\sum_{\alpha=1}^{\epsilon-1}\left(j_{\epsilon-\alpha+1,1}-j_{\epsilon-\alpha, 2}\right) r_{\epsilon-\alpha} \\
& =\left(s_{\epsilon}-j_{\epsilon, 2}\right) r_{\epsilon}+j_{1,1} r_{0}+\sum_{\alpha=1}^{\epsilon-1}\left(j_{\alpha+1,1}-j_{\alpha, 2}\right) r_{\alpha} \tag{7.7.2}
\end{align*}
$$

We now look at the place $q_{1}$. By condition (ii), $\phi_{\alpha, *, 1}^{\text {cris }}$ when restricted to $\operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\text {cris }}\right)$ is a saturated injection of $\hat{R}$-bundles; and $\psi_{\alpha, *, 1}^{\text {cris }}$ when restricted to $\operatorname{Ker}\left(\phi_{\alpha-1, *, 1}^{\text {cris }}\right)$ is also a saturated injection of $\hat{R}$-bundles. We first recall from the discussion in Notation 7.6 especially (7.6.1) that, when $\alpha=1, \ldots, \epsilon-1, \omega_{A_{\alpha}^{\vee} / R, 1}^{\circ}$ is the direct sum of

$$
\omega_{A_{\alpha}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \phi}:=\omega_{A_{\alpha}^{\vee} / R, 1}^{\circ} \cap \operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\mathrm{dR}}\right) \quad \text { and } \quad \omega_{A_{\alpha}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \psi}:=\omega_{A_{\alpha}^{\vee} / R, 1}^{\circ} \cap \operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\mathrm{dR}}\right),
$$

which are locally free $R$-modules of rank $r_{\epsilon}$ and $r_{\alpha}-r_{\epsilon}$, respectively. Similarly, put

$$
\omega_{A_{\epsilon} / R, 1}^{\circ, \operatorname{Ker} \phi}:=\omega_{A_{\epsilon}^{\vee} / R, 1}^{\circ}, \quad \omega_{A_{\epsilon}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \psi}:=0, \quad \text { and } \quad \omega_{A_{0}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \psi}=\omega_{A_{0}^{\vee} / R, 1}^{\circ} \cap \operatorname{Ker}\left(\psi_{1, *, 1}^{\mathrm{dR}}\right) ;
$$

they have ranks $r_{\epsilon}, 0$, and $r_{0}-r_{\epsilon}$, respectively. We shall avoid talking about $\omega_{A_{0}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \phi}$ (as it does not make sense) but only psychologically understand it as the process that enlarges $\omega_{A_{0}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \phi}$ to $\omega_{A_{0}^{\vee} / R, 1^{\circ}}^{\circ}$.

For $\alpha=1, \ldots, \epsilon$, the lift $\hat{\omega}_{\alpha, 1}$ takes the form of $\hat{\omega}_{\alpha, 1}^{\operatorname{Ker} \phi} \oplus \hat{\omega}_{\alpha, 1}^{\operatorname{Ker} \psi}$, where the two direct summands are $\hat{R}$-subbundles of $\operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\text {cris }}\right)$ and $\operatorname{of} \operatorname{Ker}\left(\psi_{\alpha+1, *, 1}^{\text {cris }}\right)$, lifting $\omega_{A_{\alpha}^{\alpha} / R, 1}^{\circ} \operatorname{Ker} \phi \quad$ and $\omega_{A_{\alpha}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \psi}$, respectively. Whereas, the lift $\hat{\omega}_{0,1}$ contains the lift $\hat{\omega}_{0,1}^{\operatorname{Ker} \psi}$ of $\omega_{A_{0}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \psi}$ as an $\hat{R}$-subbundle of $\operatorname{Ker}\left(\psi_{1, *, 1}^{\text {cris }}\right)$. Now the compatibility conditions $\phi_{\alpha, *, 1}^{\text {cris }}\left(\hat{\omega}_{\alpha, 1}\right) \subseteq \hat{\omega}_{\alpha-1,1}$ and $\psi_{\alpha, *, 1}^{\text {cris }}\left(\hat{\omega}_{\alpha-1,1}\right) \subseteq \hat{\omega}_{\alpha, 1}$ together with the condition (c) are equivalent to

$$
\phi_{\alpha, *, 1}^{\mathrm{cris}}\left(\hat{\omega}_{\alpha, 1}^{\operatorname{Ker} \psi}\right) \subseteq \hat{\omega}_{\alpha-1,1}^{\operatorname{Ker} \psi} \quad \text { and } \quad \psi_{\alpha, *, 1}^{\mathrm{cris}}\left(\hat{\omega}_{\alpha-1,1}^{\operatorname{Ker} \phi}\right) \subseteq \hat{\omega}_{\alpha, 1}^{\operatorname{Ker} \phi}
$$

(The condition (c) on ranks of the quotients are also automatic.) In particular, the tangent space $T_{1}$ has three contributions, coming from the lifts $\hat{\omega}_{\alpha, 1}^{\operatorname{Ker} \phi}$ (for $\alpha=1, \ldots, \epsilon$ ), from the lifts $\hat{\omega}_{\alpha, 1}^{\mathrm{Ker} \psi}$ (for $\alpha=0, \ldots, \epsilon$ ), and from lifting $\omega_{A_{0}^{\vee} / R, 1}^{\circ}$ to an $\hat{R}$-subbundle $\hat{\omega}_{0,1}$ of $H_{1}^{\text {cris }}\left(A_{0} / \hat{R}\right)_{1}^{\circ}$ containing $\hat{\omega}_{0,1}^{\operatorname{Ker} \psi}$. We shall use $T_{1}^{\text {Ker } \phi,}$ $T_{1}^{\operatorname{Ker} \psi}$, and $T_{1}^{\operatorname{Ker} \phi, 0}$ to denote these three parts of the tangent space; and they will sit in an exact sequence

$$
\begin{equation*}
0 \rightarrow T_{1}^{\mathrm{Ker} \phi, 0} \rightarrow T_{1} \rightarrow T_{1}^{\mathrm{Ker} \phi} \oplus T_{1}^{\mathrm{Ker} \psi} \rightarrow 0 \tag{7.7.3}
\end{equation*}
$$

We first determine the lifts $\hat{\omega}_{\alpha, 1}^{\text {Ker } \phi}$ for $\alpha=1, \ldots, \epsilon$. For $\hat{\omega}_{1,1}^{\text {Ker } \phi}$, it lifts $\omega_{A_{1}^{\vee} / R, 1}^{0, \text { Ker } \phi}$ as an $\hat{R}$-subbundle of $H_{1}^{\text {cris }}\left(A_{1} / \hat{R}\right)_{1}^{\circ}$ of rank $r_{\epsilon}$ (with no further constraint). Then due
to the rank constraint (and the injectivity of $\psi_{\alpha, *, 1}^{\text {cris }}$ when restricted to $\left.\operatorname{Ker}\left(\phi_{\alpha+1, *, 1}^{\text {cris }}\right)\right)$, the lift $\hat{\omega}_{\alpha, 1}^{\mathrm{Ker} \phi}$ for each $\alpha=2, \ldots, \epsilon$ is then forced to be equal to the image

$$
\psi_{\alpha, *, 1}^{\mathrm{cris}} \circ \cdots \circ \psi_{1, *, 1}^{\mathrm{cris}}\left(\hat{\omega}_{1,1}^{\operatorname{Ker} \phi}\right)
$$

So it suffices to consider the choices of the lift $\hat{\omega}_{1,1}^{\mathrm{Ker} \phi}$, which form a torsor for the group

$$
\operatorname{Hom}_{R}\left(\omega_{A_{1}^{\wedge} / R, 1}^{\circ, \operatorname{Ker} \phi}, \operatorname{Ker}\left(\phi_{1, *, 1}^{\mathrm{dR}}\right) / \omega_{A_{1}^{\prime} / R, 1}^{\circ, \operatorname{Ker} \phi}\right) \otimes_{R} \hat{I} .
$$

This Hom space is a locally free $R$-module of rank

$$
\begin{equation*}
r_{\epsilon}\left(j_{1,1}-r_{\epsilon}\right) \tag{7.7.4}
\end{equation*}
$$

It follows that the tangent space $T_{1}^{\mathrm{Ker} \phi}$ is simply just

$$
\mathcal{H o m}\left(\omega_{\mathcal{A}_{0}^{\vee}, 1}^{\circ, \operatorname{Ker} \phi}, \operatorname{Ker}\left(\phi_{1, *, 1}^{\mathrm{dR}}\right) / \omega_{\mathcal{A}_{0}^{\vee}, 1}^{\circ, \operatorname{Ker} \phi}\right)
$$

We now determine the lifts $\hat{\omega}_{\alpha, 1}^{\mathrm{Ker} \psi}$ for $\alpha=0, \ldots, \epsilon$ following the steps below:
Step 0: We start with putting $\hat{\omega}_{\epsilon, 1}^{\text {Ker } \psi}=0$ because $\omega_{A_{\epsilon} / R, 1}^{\circ \text { Ker } \psi}$ is,
$\operatorname{Step}(\mathbf{s}) \alpha: \operatorname{lift} \omega_{A_{\epsilon-\alpha}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \psi}$ to a subbundle $\hat{\omega}_{\epsilon-\alpha, 1}^{\operatorname{Ker} \psi} \operatorname{of} \operatorname{Ker}\left(\psi_{\epsilon-\alpha+1, *, 1}^{\text {cris }}\right)$ so that it contains $\phi_{\epsilon-\alpha+1}^{\text {cris }}\left(\hat{\omega}_{\epsilon-\alpha+1,1}^{\operatorname{Ker} \psi}\right)$,
Step $\boldsymbol{\epsilon}$ : finally lift $\omega_{A_{0}^{\gamma} / R, 1}^{\circ, \operatorname{Ker} \psi}$ to a subbundle $\hat{\omega}_{0,1}^{\operatorname{Ker} \psi} \operatorname{of} \operatorname{Ker}\left(\psi_{1, *, 1}^{\text {cris }}\right)$ so that it contains $\phi_{1, *, 1}^{\mathrm{cris}}\left(\hat{\omega}_{1,1}^{\mathrm{Ker} \psi}\right)$.

At Step $\alpha=1, \ldots, \epsilon$, the choices of the lifts $\hat{\omega}_{\epsilon-\alpha, 1}^{\mathrm{Ker} \psi}$ form a torsor for the group

$$
\operatorname{Hom}_{R}\left(\omega_{A_{\epsilon-\alpha} / R, 1}^{\circ, \operatorname{Ker} \psi} / \phi_{\epsilon-\alpha+1, *, 1}\left(\omega_{A_{\epsilon-\alpha+1}^{\prime}}^{\circ, \operatorname{Ker} \psi} / R, 1\right), \operatorname{Ker}\left(\psi_{\epsilon-\alpha+1, *, 1}^{\mathrm{KR}}\right) / \omega_{A_{\epsilon-\alpha}^{\circ} / R, 1}^{\circ \mathrm{Ker} \psi}\right) \otimes_{R} \hat{I}
$$

This Hom space is a locally free $R$-module of rank

$$
\left(\left(r_{\epsilon-\alpha}-r_{\epsilon}\right)-\left(r_{\epsilon-\alpha+1}-r_{\epsilon}\right)\right)\left(\left(n-j_{\epsilon-\alpha+1,1}\right)-\left(r_{\epsilon-\alpha}-r_{\epsilon}\right)\right) .
$$

This implies that the tangent space $T_{1}^{\mathrm{Ker} \psi}$ admits a filtration such that the subquotients are

$$
\mathcal{H o m}\left(\omega_{\mathcal{A}_{\epsilon-\alpha}^{\vee}, 1}^{\circ, \operatorname{Ker} \psi} / \phi_{\epsilon+1-\alpha, *, 1}\left(\omega_{\mathcal{A}_{\epsilon+1-\alpha}^{\vee}, 1}^{\circ, \operatorname{Ker} \psi}\right), \operatorname{Ker}\left(\psi_{\epsilon+1-\alpha, *, 1}^{\mathrm{VR}}\right) / \omega_{\mathcal{A}_{\epsilon-\alpha}^{\vee}, 1}^{\circ, \operatorname{Ker} \psi}\right)
$$

In particular, $T_{1}^{\mathrm{Ker} \psi}$ is a locally free sheaf on $Y_{j}^{\circ}$ of rank

$$
\begin{align*}
& \sum_{\alpha=1}^{\epsilon}\left(\left(r_{\epsilon-\alpha}-r_{\epsilon}\right)-\left(r_{\epsilon-\alpha+1}-r_{\epsilon}\right)\right)\left(\left(n-j_{\epsilon-\alpha+1,1}\right)-\left(r_{\epsilon-\alpha}-r_{\epsilon}\right)\right) \\
&=\sum_{\alpha=0}^{\epsilon-1}\left(r_{\alpha}-r_{\alpha+1}\right)\left(s_{\alpha}-j_{\alpha+1,1}+r_{\epsilon}\right) \tag{7.7.5}
\end{align*}
$$

Finally, we discuss the $\hat{R}$-module $\hat{\omega}_{0,1}$ that lifts $\omega_{A_{0}^{\vee} / R, 1}^{\circ}$ and contains $\hat{\omega}_{0,1}^{\mathrm{Ker} \psi}$ we obtained earlier. The lift is subject to one condition: $\hat{\omega}_{0,1} \subseteq\left(\psi_{1, *, 1}^{\text {cris }}\right)^{-1}\left(\hat{\omega}_{1,1}^{\text {Ker } \phi}\right)$. So the choices of the lift form a torsor for the group

$$
\operatorname{Hom}_{R}\left(\omega_{A_{0}^{\vee} / R, 1}^{\circ} / \omega_{A_{0}^{\vee} / R, 1}^{\circ, \operatorname{Ker} \psi},\left(\psi_{1, *, 1}^{\mathrm{dR}}\right)^{-1}\left(\omega_{A_{1}^{\vee} / R, 1}^{\circ} \mathrm{Ker} \phi\right) / \omega_{A_{0}^{\vee} / R, 1}^{\circ}\right) \otimes_{R} \hat{I}
$$

This implies that

$$
T_{1}^{\operatorname{Ker} \phi, 0}=\mathcal{H} \operatorname{com}\left(\omega_{\mathcal{A}_{0}^{\vee}, 1}^{\circ} / \omega_{\mathcal{A}_{0}^{\vee}, 1}^{\circ, \operatorname{Ker} \psi},\left(\psi_{1, *, 1}^{\mathrm{dR}}\right)^{-1}\left(\omega_{\mathcal{A}_{1}^{\vee}, 1}^{\circ, \operatorname{Ker} \phi}\right) / \omega_{\mathcal{A}_{0}^{\vee}, 1}^{\circ}\right)
$$

which is locally free of rank

$$
\begin{equation*}
\left(r_{0}-\left(r_{0}-r_{\epsilon}\right)\right)\left(\left(r_{\epsilon}+n-j_{1,1}\right)-r_{0}\right)=r_{\epsilon}\left(s_{0}+r_{\epsilon}-j_{1,1}\right) \tag{7.7.6}
\end{equation*}
$$

To sum up, the tangent space $T_{Y_{j}^{\circ}}$, as the direct sum $T_{1} \oplus T_{2}$ with $T_{1}$ sitting in the exact sequence (7.7.3), is a locally free sheaf of rank given by (7.7.6) + (7.7.4) + (7.7.5) + (7.7.2), that is,

$$
\begin{aligned}
& r_{\epsilon}\left(s_{0}+r_{\epsilon}-j_{1,1}\right)+r_{\epsilon}\left(j_{1,1}-r_{\epsilon}\right)+ \sum_{\alpha=0}^{\epsilon-1}\left(r_{\alpha}-r_{\alpha+1}\right)\left(s_{\alpha}-j_{\alpha+1,1}+r_{\epsilon}\right) \\
&+\left(s_{\epsilon}-j_{\epsilon, 2}\right) r_{\epsilon}+j_{1,1} r_{0}+\sum_{\alpha=1}^{\epsilon-1}\left(j_{\alpha+1,1}-j_{\alpha, 2}\right) r_{\alpha} \\
&= r_{\epsilon} s_{0}+\sum_{\alpha=0}^{\epsilon-1} r_{\alpha}\left(s_{\alpha}-j_{\alpha+1,1}+r_{\epsilon}\right)-\sum_{\alpha=1}^{\epsilon} r_{\alpha}\left(s_{\alpha-1}-j_{\alpha, 1}+r_{\epsilon}\right) \\
&+\left(s_{\epsilon}-j_{\epsilon, 2}\right) r_{\epsilon}+j_{1,1} r_{0}+\sum_{\alpha=1}^{\epsilon-1}\left(j_{\alpha+1,1}-j_{\alpha, 2}\right) r_{\alpha} \\
&=r_{\epsilon} s_{0}+r_{0}\left(s_{0}-j_{1,1}+r_{\epsilon}\right)+r_{\epsilon}\left(s_{\epsilon-1}-j_{\epsilon, 1}+r_{\epsilon}\right)+\left(s_{\epsilon}-j_{\epsilon, 2}\right) r_{\epsilon}+j_{1,1} r_{0} \\
& \quad+\sum_{\alpha=1}^{\epsilon-1} r_{\alpha}\left(\left(s_{\alpha}-j_{\alpha+1,1}+r_{\epsilon}\right)-\left(s_{\alpha-1}-j_{\alpha, 1}+r_{\epsilon}\right)+\left(j_{\alpha+1,1}-j_{\alpha, 2}\right)\right)
\end{aligned}
$$

One easily checks that the first line adds up to $r_{\epsilon} s_{\epsilon}+r_{0} s_{0}$, and the second line cancels to zero. This concludes the proof.

In the special case of $\delta=r$, each abelian variety $A_{\alpha}$ appearing in the moduli problem of $Y_{j}$ is isogenous to $A_{\epsilon}$, which is a certain abelian variety parameterized by the discrete Shimura variety $\mathrm{Sh}_{0, n}$ and is hence supersingular (by Remark 3.7). So in particular, the image $\operatorname{pr}_{j}\left(Y_{j}\right)$ in this case is contained in the supersingular locus of $\mathrm{Sh}_{r, s}$. In fact, the converse is also true.
Theorem 7.8. Assume $\delta=r$. The supersingular locus of $\mathrm{Sh}_{r, s}$ is the union of all $\operatorname{pr}_{j}\left(Y_{j}\right)$.

Proof. We say a finite torsion $W\left(\overline{\mathbb{F}}_{p}\right)$-module has divisible sequence $\left(a_{1}, a_{2}, \ldots, a_{\epsilon}\right)$ with nonnegative integers $a_{1} \leq \cdots \leq a_{\epsilon}$ if it is isomorphic to

$$
\left(W\left(\overline{\mathbb{F}}_{p}\right) / p^{\epsilon}\right)^{\oplus a_{1}} \oplus\left(W\left(\overline{\mathbb{F}}_{p}\right) / p^{\epsilon-1}\right)^{\oplus\left(a_{2}-a_{1}\right)} \oplus \cdots \oplus\left(W\left(\overline{\mathbb{F}}_{p}\right) / p\right)^{\oplus\left(a_{\epsilon}-a_{\epsilon-1}\right)}
$$

The following is an elementary linear algebra fact, whose proof we omit.
Claim: If $M_{1} \subseteq M_{2}$ are two torsion $W\left(\overline{\mathbb{F}}_{p}\right)$-modules with divisible sequences $\left(a_{1, i}, \ldots, a_{\epsilon, i}\right)$ for $i=1,2$ respectively, then $a_{\alpha, 1} \leq a_{\alpha, 2}$ for all $\alpha=1, \ldots, \epsilon$.

The proof of the theorem is similar to the proof of Proposition 4.14(3), which is a special case of this theorem. It suffices to look at the closed points of $\mathrm{Sh}_{r, s}$. Let $z=\left(\mathcal{A}_{z}, \lambda, \eta\right) \in \operatorname{Sh}_{r, s}\left(\overline{\mathbb{F}}_{p}\right)$ be a supersingular point. Consider

$$
\mathbb{L}_{\mathbb{Q}}=\left(\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}[1 / p]\right)^{F^{2}=p}=\left\{a \in \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}[1 / p] \mid F^{2}(a)=p a\right\} .
$$

Since $x$ is supersingular, $\mathbb{L}_{\mathbb{Q}}$ is a $\mathbb{Q}_{p^{2}}$-vector space of dimension $n$, and $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}[1 / p]$ may be identified with the extension of scalars of $\mathbb{Q}_{\mathbb{Q}}$ from $\mathbb{Q}_{p^{2}}$ to $W\left(\overline{\mathbb{F}}_{p}\right)[1 / p]$. Put

$$
\tilde{\mathcal{E}}_{1}^{\circ}=\left(\mathbb{Q}_{\mathbb{Q}} \cap \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}\right) \otimes_{\mathbb{Z}_{p^{2}}} W\left(\overline{\mathbb{F}}_{p}\right) \quad \text { and } \quad \tilde{\mathcal{E}}_{2}^{\circ}=F\left(\tilde{\mathcal{E}}_{1}^{\circ}\right)=V\left(\tilde{\mathcal{E}}_{1}^{\circ}\right) \subseteq \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ}
$$

Then we have

$$
\begin{align*}
& \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{i}^{\circ} / \tilde{\mathcal{E}}_{i} \simeq\left(W\left(\overline{\mathbb{F}}_{p}\right) / p^{\epsilon}\right)^{\oplus j_{1, i}} \oplus\left(W\left(\overline{\mathbb{F}}_{p}\right) / p^{\epsilon-1}\right)^{\oplus\left(j_{2, i}-j_{1, i}\right)} \oplus \cdots \\
& \cdots \oplus\left(W\left(\overline{\mathbb{F}}_{p}\right) / p\right)^{\oplus\left(j_{\epsilon, i}-j_{\epsilon-1, i}\right)} \tag{7.8.1}
\end{align*}
$$

for nondecreasing sequences $0 \leq j_{1, i} \leq j_{2, i} \leq \cdots \leq j_{\epsilon, i} \leq n$ with $i=1,2$; in other words, $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{i}^{\circ} / \tilde{\mathcal{E}}_{i}$ has divisible sequence $\left(j_{1, i}, \ldots, j_{\epsilon, i}\right)$. Without loss of generality, we assume that $j_{1,1}$ and $j_{1,2}$ are not both zero. The essential part of the proof consists of checking the sequence of inequalities

$$
\begin{equation*}
0 \leq j_{1,1}<j_{1,2}<j_{2,1}<j_{2,2}<\cdots<j_{\epsilon, 1}<j_{\epsilon, 2} \leq n \tag{7.8.2}
\end{equation*}
$$

We first prove (7.8.2) with all strict inequalities replaced by nonstrict ones. Indeed, the obvious inclusion $F\left(\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{i}^{\circ}\right) \subseteq \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{3-i}^{\circ}$ implies that

$$
\begin{aligned}
& F\left(\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} / \tilde{\mathcal{E}}_{1}\right)=F\left(\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}\right) / \tilde{\mathcal{E}}_{2} \subseteq \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} / \tilde{\mathcal{E}}_{2}, \quad \text { and } \\
& F\left(\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} / \tilde{\mathcal{E}}_{2}\right)=F\left(\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ}\right) / p \tilde{\mathcal{E}}_{1} \subseteq \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} / p \tilde{\mathcal{E}}_{1}
\end{aligned}
$$

By (7.8.1), the first inclusion embeds a torsion $W\left(\bar{F}_{p}\right)$-module with divisible sequence $\left(j_{1,1}, \ldots, j_{\epsilon, 1}\right)$ into a torsion $W\left(\overline{\mathbb{F}}_{p}\right)$-module with divisible sequence $\left(j_{1,2}, \ldots, j_{\epsilon, 2}\right)$. The Claim above implies that $j_{\alpha, 1} \leq j_{\alpha, 2}$ for all $\alpha=1, \ldots, \epsilon$. Similarly, by (7.8.1), the second inclusion embeds a torsion $W\left(\overline{\mathbb{F}}_{p}\right)$-module with divisible sequence $\left(j_{1,2}, \ldots, j_{\epsilon, 2}\right)$ into a torsion $W\left(\overline{\mathbb{F}}_{p}\right)$-module with divisible sequence $\left(j_{1,1}, \ldots, j_{\epsilon, 1}, n\right)$. The Claim above implies that $j_{\alpha, 2} \leq j_{\alpha+1,1}$ for all $\alpha=1, \ldots, \epsilon-1$, and $j_{\epsilon, 2} \leq n$.

We now use the construction of $\mathbb{L}_{\mathbb{Q}}$ to show the strict inequalities in (7.8.2). Suppose first that $j_{\alpha, 1}=j_{\alpha, 2}$ for some $\alpha=1, \ldots, \epsilon$. Then it follows that the maps

$$
\begin{equation*}
F, \quad V:\left(p^{\epsilon-\alpha} \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \frac{1}{p} \tilde{\mathcal{E}}_{1}^{\circ}\right)+\tilde{\mathcal{E}}_{1}^{\circ} \rightarrow\left(p^{\epsilon-\alpha} \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \cap \frac{1}{p} \tilde{\mathcal{E}}_{2}^{\circ}\right)+\tilde{\mathcal{E}}_{2}^{\circ} \tag{7.8.3}
\end{equation*}
$$

are both isomorphisms (due to an easy length computation as $\left.\tilde{\mathcal{E}}_{2}^{\circ}=F\left(\tilde{\mathcal{E}}_{1}^{\circ}\right)=V\left(\tilde{\mathcal{E}}_{1}^{\circ}\right)\right)$. By the definition of $\mathbb{L}_{\mathbb{Q}}$ and $\tilde{\mathcal{E}}_{1}^{\circ}$, we must have

$$
\left(\left(p^{\epsilon-\alpha} \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \frac{1}{p} \tilde{\mathcal{E}}_{1}^{\circ}\right)+\tilde{\mathcal{E}}_{1}^{\circ}\right)^{F=V} \subseteq \mathbb{L}_{\mathbb{Q}} \cap \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \subseteq \tilde{\mathcal{E}}_{1}^{\circ}
$$

But this is absurd because the isomorphisms (7.8.3) implies by Hilbert's Theorem 90 that the left hand side above generates the source of (7.8.3), which is clearly not contained in $\tilde{\mathcal{E}}_{1}^{\circ}$.

Similarly, suppose that $j_{\alpha, 2}=j_{\alpha+1,1}$ for some $\alpha=1, \ldots, \epsilon-1$. Then the following morphisms are isomorphisms

$$
\begin{equation*}
F, V:\left(p^{\epsilon-\alpha} \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \cap \frac{1}{p} \tilde{\mathcal{E}}_{2}^{\circ}\right)+\tilde{\mathcal{E}}_{2}^{\circ} \rightarrow\left(p^{\epsilon-\alpha} \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \tilde{\mathcal{E}}_{1}^{\circ}\right)+p \tilde{\mathcal{E}}_{1}^{\circ} \tag{7.8.4}
\end{equation*}
$$

since $p \tilde{\mathcal{E}}_{1}^{\circ}=F\left(\tilde{\mathcal{E}}_{2}^{\circ}\right)=V\left(\tilde{\mathcal{E}}_{2}^{\circ}\right)$ and for length reasons. By the definition of $\mathbb{L}_{\mathbb{Q}}$ and $\tilde{\mathcal{E}}_{1}^{\circ}$,

$$
\left(\left(p^{\epsilon-\alpha} \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \tilde{\mathcal{E}}_{1}^{\circ}\right)+p \tilde{\mathcal{E}}_{1}^{\circ}\right)^{F^{-1}=V^{-1}} \subseteq \mathbb{L}_{\mathbb{Q}} \cap p \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \subseteq p \tilde{\mathcal{E}}_{1}^{\circ}
$$

(Note that $\epsilon-\alpha \geq 1$ now.) But this is absurd because the isomorphisms (7.8.4) imply by Hilbert's Theorem 90 that the left hand side above generates the target of (7.8.4), which is clearly not contained in $p \tilde{\mathcal{E}}_{1}^{\circ}$.

Summing up, we have proved the strict inequalities (7.8.2). So the $j_{\alpha, i}$ define a $\boldsymbol{j}$ as in the beginning of Section 7.1. We now construct a point of $Y_{\boldsymbol{j}}$ which maps to the point $z \in \mathrm{Sh}_{r, s}$. Put

$$
\begin{equation*}
\tilde{\mathcal{E}}_{\alpha, 1}:=\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha}} \tilde{\mathcal{E}}_{1} \quad \text { and } \quad \tilde{\mathcal{E}}_{\alpha, 2}:=\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha}} \tilde{\mathcal{E}}_{2} \tag{7.8.5}
\end{equation*}
$$

Using the exact construction in Section 7.2, we get the sequence of isogenies of abelian varieties

$$
A_{\epsilon} \underset{\psi_{\epsilon}}{\stackrel{\phi_{\epsilon}}{\rightleftarrows}} A_{\epsilon-1} \stackrel{\phi_{\epsilon-1}}{\underset{\psi_{\epsilon-1}}{\rightleftarrows}} \cdots \stackrel{\phi_{1}}{\underset{\psi_{1}}{\rightleftarrows}} A_{0}=\mathcal{A}_{z},
$$

such that $A_{\alpha}$ together with the induced polarization $\lambda_{\alpha}$ and the tame level structure $\eta_{\alpha}$ gives an $\overline{\mathbb{F}}_{p}$-point of $\mathrm{Sh}_{r_{\alpha}, s_{\alpha}}$, and $\tilde{\mathcal{D}}\left(A_{\alpha}\right)_{i}^{\circ}=\tilde{\mathcal{E}}_{\alpha, i}$ for all $\alpha$ and $i=1,2$.

Conditions (2)-(5) of Definition 7.4 easily follow from the description of the quotients $\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{i}^{\circ} / \tilde{\mathcal{E}}_{i}$ in (7.8.1). Condition (6) of Definition 7.4 is equivalent to

$$
p \tilde{\mathcal{D}}\left(A_{\alpha-1}\right)_{2}^{\circ} \subseteq V\left(\tilde{\mathcal{D}}\left(A_{\alpha}\right)_{1}^{\circ}\right)
$$

By the construction of these Dieudonné modules in (7.8.5), this is equivalent to

$$
p\left(\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha+1}} \tilde{\mathcal{E}}_{2}\right) \subseteq V\left(\tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha}} \tilde{\mathcal{E}}_{1}\right)
$$

But this follows from $p \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \subseteq V \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}$ and $\tilde{\mathcal{E}}_{2}=V \tilde{\mathcal{E}}_{1}$. Condition (7) of Definition 7.4 is equivalent to $\omega_{A_{\alpha}^{\vee} / \widetilde{F}_{p}, 1}^{\circ} \cap \operatorname{Ker}\left(\phi_{\alpha, *, 1}^{\mathrm{dR}}\right)$ having dimension $r_{\epsilon}$, which is zero in our case. Translating it into the language of Dieudonné modules, this is equivalent to

$$
V \tilde{\mathcal{D}}\left(A_{\alpha}\right)_{2}^{\circ} \cap p \tilde{\mathcal{D}}\left(A_{\alpha-1}\right)_{1}^{\circ}=p \tilde{\mathcal{D}}\left(A_{\alpha}\right)_{1}^{\circ} .
$$

By the construction of these Dieudonné module in (7.8.5), this is equivalent to

$$
\left(V \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha}} V \tilde{\mathcal{E}}_{2}\right) \cap\left(p \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha}} \tilde{\mathcal{E}}_{1}\right)=p \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha-1}} \tilde{\mathcal{E}}_{1}
$$

which follows from observing that $V \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \supseteq p \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}$ and $V \tilde{\mathcal{E}}_{2}=p \tilde{\mathcal{E}}_{1}$. Condition (8) of Definition 7.4 is equivalent to $\omega_{A_{\alpha-1}^{\vee} / \mathbb{F}_{p}, 1}^{\circ} \subseteq \operatorname{Ker}\left(\psi_{\alpha, *, 1}^{\mathrm{dR}}\right)$ (note that $r_{\epsilon}=0$ in our case). Translating it into the language of Dieudonné modules and using (7.8.5), this is equivalent to

$$
\begin{aligned}
V \tilde{\mathcal{D}}\left(A_{\alpha-1}\right)_{2}^{\circ} & \subseteq \tilde{\mathcal{D}}\left(A_{\alpha}\right)_{1}^{\circ}, \text { or equivalently, } \\
V \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha+1}} V \tilde{\mathcal{E}}_{2} & \subseteq \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ} \cap \frac{1}{p^{\epsilon-\alpha}} \tilde{\mathcal{E}}_{1}
\end{aligned}
$$

which follows from observing that $V \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{2}^{\circ} \subseteq \tilde{\mathcal{D}}\left(\mathcal{A}_{z}\right)_{1}^{\circ}$ and $V \tilde{\mathcal{E}}_{2}=p \tilde{\mathcal{E}}_{1}$. This concludes the proof.

Conjecture 7.9. The varieties $Y_{j}$ together with the natural morphisms to $\mathrm{Sh}_{r-\delta, s+\delta}$ and $\mathrm{Sh}_{r, s}$ satisfy condition (3) of Conjecture 2.12. Moreover, the union of the images of $Y_{j}$ in $\mathrm{Sh}_{r, s}$ is the closure of the locus where the Newton polygon of the universal abelian variety has slopes 0 and 1 each with multiplicity $2(r-\delta) n$, and slope $\frac{1}{2}$ with multiplicity $2(n-2 r+2 \delta) n$.

This conjecture in the case of $r=\delta=1$ was proved in Theorem 4.18.

## Appendix A: An explicit formula in the local spherical Hecke algebra for $\mathbf{G L}_{\boldsymbol{n}}$

In this appendix, let $F$ be a local field with ring of integers $\mathcal{O}, \varpi \in \mathcal{O}$ be a uniformizer, $\mathbb{F}=\mathcal{O} / \varpi \mathcal{O}$ and $q=\# \mathbb{F}$. Fix an integer $n \geq 1$. We consider the spherical Hecke algebra $\mathscr{H}_{K}=\mathbb{Z}\left[K \backslash \mathrm{GL}_{n}(F) / K\right]$ with $K=\mathrm{GL}_{n}(\mathcal{O})$. Here, the
product of two double cosets $u=K x K$ and $v=K y K$ in $\mathscr{H}_{K}$ is defined as

$$
\begin{equation*}
u \cdot v=\sum_{w} m(u, v ; w) w,{ }^{19} \tag{A.0.1}
\end{equation*}
$$

where the sum runs through all the double cosets $w=K z K$ contained in $K x K y K$, and the coefficient $m(u, v ; w) \in \mathbb{Z}$ is determined as follows: If $K x K=\coprod_{i \in I} x_{i} K$ and $K y K=\coprod_{j \in J} y_{j} K$, then

$$
m(u, v ; w)=\#\left\{(i, j) \in I \times J \mid x_{i} y_{j} K=z K \text { for a fixed element } z \text { in } w\right\} . \text { (A.0.2) }
$$

By the theory of elementary divisors, all double cosets $K x K$ are of the form
$T\left(a_{1}, \ldots, a_{n}\right):=K \operatorname{Diag}\left(\varpi^{a_{1}}, \ldots, \varpi^{a_{n}}\right) K \quad$ for $a_{i} \in \mathbb{Z}$ with $a_{1} \geq a_{2} \geq \cdots \geq a_{n}$.
They form a $\mathbb{Z}$-basis of $\mathscr{H}_{K}$. We put

$$
\begin{array}{rlr}
T^{(r)} & =T(\underbrace{1, \ldots, 1}_{r}, \underbrace{0, \ldots, 0}_{n-r}) & \text { for } 0 \leq r \leq n, \\
R^{(r, s)} & =T(\underbrace{2, \ldots, 2}_{r}, \underbrace{1, \ldots, 1}_{s-r}, \underbrace{0, \ldots, 0}_{n-s} & \text { for } 0 \leq r \leq s \leq n .
\end{array}
$$

In particular, $R^{(0, s)}=T^{(s)}$ and $T^{(0)}=[K]$.
Because of the lack of references, we include a proof of the following:
Proposition A.1. For $1 \leq r \leq n$, let

$$
\begin{equation*}
\binom{n}{r}_{q}=\frac{\left(q^{n}-1\right)\left(q^{n-1}-1\right) \cdots\left(q^{n-r+1}-1\right)}{(q-1)\left(q^{2}-1\right) \cdots\left(q^{r}-1\right)} \tag{A.1.1}
\end{equation*}
$$

be the Gaussian binomial coefficients, and put $\binom{n}{0}_{q}=1$. Then for $0 \leq r \leq s \leq n$,

$$
T^{(r)} T^{(s)}=\sum_{i=0}^{\min \{r, n-s\}}\binom{s-r+2 i}{i}_{q} R^{(r-i, s+i)}
$$

Proof. We fix a set of representatives $\tilde{\mathbb{F}} \subseteq \mathcal{O}$ of $\mathbb{F}=\mathcal{O} / \varpi \mathcal{O}$ which contains 0 . Then we have $T^{(r)}=\coprod_{x \in \mathcal{S}(n, r)} x K$, where $\mathcal{S}(n, r)$ is the set of $n \times n$ matrices $x=\left(x_{i, j}\right)_{1 \leq i, j \leq n}$ such that

- $r$ of the diagonal entries are $\varpi$ and the remaining $n-r$ ones are 1 ;
- if $i \neq j$, then $x_{i, j}=0$ unless $i>j, x_{i, i}=1$ and $x_{j, j}=\varpi$, in which case $x_{i, j}$ can take any values in $\tilde{\mathbb{F}}$.

[^14]For instance, the set $\mathcal{S}(3,2)$ consists of matrices:

$$
\left(\begin{array}{ccc}
1 & 0 & 0 \\
x_{2,1} & \varpi & 0 \\
x_{3,1} & 0 & \varpi
\end{array}\right), \quad\left(\begin{array}{ccc}
\varpi & 0 & 0 \\
0 & 1 & 0 \\
0 & x_{3,2} & \varpi
\end{array}\right), \quad\left(\begin{array}{ccc}
\varpi & 0 & 0 \\
0 & \varpi & 0 \\
0 & 0 & 1
\end{array}\right),
$$

with $x_{2,1}, x_{3,1}, x_{3,2} \in \tilde{\mathbb{F}}$. We have a similar decomposition $T^{(s)}=\bigsqcup_{y \in \mathcal{S}(n, s)} y K$. We write $T^{(r)} T^{(s)}$ as a linear combination of $T\left(a_{1}, \ldots, a_{n}\right)$ with $a_{i} \in \mathbb{Z}$ and $a_{1} \geq \cdots \geq a_{n}$. By looking at the diagonal entries of $x y$, we see easily that only $R^{(r-i, s+i)}$ with $0 \leq i \leq \min \{r, n-s\}$ have nonzero coefficients, namely, we have

$$
T^{(r)} T^{(s)}=\sum_{i=0}^{\min \{r, n-s\}} C^{(r, s)}(n, i) R^{(r-i, s+i)} \quad \text { for some } C^{(r, s)}(n, i) \in \mathbb{Z}
$$

By (A.0.1), $C^{(r, s)}(n, i)$ is the number of pairs $(x, y) \in \mathcal{S}(n, r) \times \mathcal{S}(n, s)$ such that

$$
x y K=\operatorname{Diag}(\underbrace{\varpi^{2}, \ldots, \varpi^{2}}_{r-i}, \underbrace{\varpi, \ldots, \varpi}_{s-r+2 i}, \underbrace{1, \ldots, 1}_{n-s-i}) K .
$$

In this case, $x$ and $y$ must be of the form

$$
x=\left(\begin{array}{ccc}
\varpi I_{r-i} & 0 & 0 \\
0 & A & 0 \\
0 & 0 & I_{n-s-i}
\end{array}\right), \quad y=\left(\begin{array}{ccc}
\varpi I_{r-i} & 0 & 0 \\
0 & B & 0 \\
0 & 0 & I_{n-s-i}
\end{array}\right),
$$

where $I_{k}$ denotes the $k \times k$ identity matrix, and $A \in \mathcal{S}(s-r+2 i, i)$ and $B \in$ $\mathcal{S}(s-r+2 i, s-r+i)$ satisfy $A B \cdot \mathrm{GL}_{s-r+2 i}(\mathcal{O})=\varpi I_{s-r+2 i} \mathrm{GL}_{s-r+2 i}(\mathcal{O})$. By (A.0.1), we see that $C^{(r, s)}(n, i)=C^{(i, s-r+i)}(s-r+2 i, i)$. Therefore, one is reduced to proving the following lemma, which is a special case of our proposition.

Lemma A.2. Under the notation and hypothesis of Proposition A.1, assume moreover that $n=r+s$. Then the coefficient of $R^{(0, n)}$ in the product $T^{(r)} T^{(s)}$ is $\binom{n}{r}_{q}$.

Proof. We induct on $n \geq 1$. The case $n=1$ is trivial. We assume thus $n>1$, and that the statement is true when $n$ is replaced by $n-1$. The case of $r=0$ being trivial, we may assume that $r \geq 1$. We say a pair $(x, y) \in \mathcal{S}(n, r) \times \mathcal{S}(n, n-r)$ is admissible if $x y K=\varpi I_{n} K$. We have to show that the number of admissible pairs is equal to $\binom{n}{r}_{q}$. Let $(x, y)$ be an admissible pair. Denote by $I$ (resp. by $J$ ) the set integers $1 \leq i \leq n$ such that $x_{i, i}=\varpi$ (resp. $\left.y_{i, i}=\varpi\right)$. Note that $(x, y)$ being admissible implies that $J=\{1, \ldots, n\} \backslash I$.

Assume first that $x_{1,1}=1$. Then $x$ and $y$ must be of the form $x=\left(\begin{array}{ll}1 & 0 \\ * & A\end{array}\right)$ and $y=\left(\begin{array}{cc}\infty & 0 \\ 0 & B\end{array}\right)$ where $(A, B) \in \mathcal{S}(n-1, r) \times \mathcal{S}(n-1, n-1-r)$ admissible. Note that $x y K=\varpi I_{n} K$ always hold. We have $x_{i, 1}=0$ for $i \notin I$, and $x_{i, 1}$ can take any values in $\mathbb{F}$ for $i \in I$. Therefore, the number of admissible pairs $(x, y)$ with $x_{1,1}=1$ is
equal to $q^{\# I}=q^{r}$ times that of the admissible $(A, B)$. The latter is equal to $\binom{n-1}{r}_{q}$ by the induction hypothesis.

Consider now the case $x_{1,1}=\varpi$. One can write $x=\left(\begin{array}{cc}\varpi & 0 \\ 0 & A\end{array}\right)$, and $y=\left(\begin{array}{cc}1 & 0 \\ * & B\end{array}\right)$ with $(A, B) \in \mathcal{S}(n-1, r-1) \times \mathcal{S}(n-1, n-r)$ admissible. Put $z=x y$. Then an easy computation shows that $z_{j, 1}=y_{j, 1}$ if $j \in J$, and $z_{j, 1}=0$ if $j \notin J$. Hence, $x y K=\varpi I_{n} K$ forces that $y_{j, 1}=0$ for all $j>1$. Therefore, the number of admissible $(x, y)$ in this case is equal to that of the admissible $(A, B)$, which is $\binom{n-1}{r-1}_{q}$ by the induction hypothesis. The lemma now follows immediately from the equality

$$
\binom{n}{r}_{q}=q^{r}\binom{n-1}{r}_{q}+\binom{n-1}{r-1}_{q}
$$

## Appendix B: A determinant formula

In this appendix, we prove the following:
Theorem B.1. Let $\alpha_{1}, \ldots, \alpha_{n}$ be $n$ indeterminates. For $i=1, \ldots, n$, let $s_{i}$ denote the $i$-th elementary symmetric polynomial in the $\alpha$, and $s_{0}=1$ by convention. Let $q$ be another indeterminate. We put $q_{r}=q^{r-1}+q^{r-3}+\cdots+q^{1-r}$. Consider the matrix $M_{n}(q)=\left(m_{i, j}\right)$ given as follows:

$$
m_{i, j}=\left\{\begin{array}{ll}
\sum_{\delta=0}^{\min \{i-1, n-j\}} \\
\sum_{\delta=0}^{\min \{j-1, n-i\}} & q_{n+i-j-2 \delta} s_{j-i+\delta} s_{n-\delta}
\end{array} \quad \text { if } i \leq j, ~ q_{n+j-i-2 \delta} s_{\delta} s_{n+j-i+\delta} \quad \text { if } i>j\right.
$$

Then we have

$$
\operatorname{det}\left(M_{n}(q)\right)=\alpha_{1} \cdots \alpha_{n} \prod_{i \neq j}\left(q \alpha_{i}-\frac{1}{q} \alpha_{j}\right)
$$

Proof. Let $N_{n}(q)$ be the resultant matrix of the polynomials $f(x)=\prod_{i=1}^{n}\left(x+q^{-1} \alpha_{i}\right)$ and $g(x)=\prod_{i=1}^{n}\left(x+q \alpha_{i}\right)$, that is, $N_{n}(q)$ is the $2 n \times 2 n$ matrix given by

$$
N_{n}(q)=\left(\begin{array}{cccccccc}
s_{0} & q^{-1} s_{1} & q^{-2} s_{2} & \cdots & q^{1-n} s_{n-1} & q^{-n} s_{n} & 0 & \cdots \\
0 & s_{0} & q^{-1} s_{1} & \cdots & q^{2-n} s_{n-2} & q^{1-n} s_{n-1} & q^{-n} s_{n} & \cdots \\
0 & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots & \vdots & \vdots & \ddots
\end{array}\right) \vdots
$$

It is well known that $\operatorname{det}\left(N_{n}(q)\right)=\prod_{i, j}\left(-q^{-1} \alpha_{i}+q \alpha_{j}\right)$. Thus it suffices to show that $\operatorname{det}\left(N_{n}(q)\right)=\left(q-q^{-1}\right)^{n} \operatorname{det}\left(M_{n}(q)\right)$.

We first make the following row operations on $N_{n}(q)$ : subtract row $i$ from row $n+i$ for all $i=1, \ldots, n$. We obtain a matrix whose first column is all 0 except
the first entry being 1 ; moreover, one can take out a factor $\left(q-q^{-1}\right)$ from row $n+1, \ldots, 2 n$. Let $N_{n}^{\prime}(q)$ be the right lower $(2 n-1) \times(2 n-1)$ submatrix of the remaining matrix. Then we have

$$
N_{n}^{\prime}(q)=\left(\begin{array}{ccccccccc}
s_{0} & q^{-1} s_{1} & q^{-2} s_{2} & \cdots & q^{1-n} s_{n-1} & q^{-n} s_{n} & 0 & \cdots & 0 \\
0 & s_{0} & q^{-1} s_{1} & \cdots & q^{2-n} s_{n-2} & q^{1-n} s_{n-1} & q^{-n} s_{n} & \cdots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \cdots & s_{0} & q^{-1} s_{1} & q^{-2} s_{2} & \cdots & q^{-n} s_{n} \\
q_{1} s_{1} & q_{2} s_{2} & q_{3} s_{3} & \cdots & q_{n-1} s_{n-1} & q_{n} s_{n} & 0 & \cdots & 0 \\
0 & q_{1} s_{1} & q_{2} s_{2} & \cdots & q_{n-2} s_{n-2} & q_{n-1} s_{n-1} & q_{n} s_{n} & \cdots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \cdots & 0 & q_{1} s_{1} & q_{2} s_{2} & \cdots & q_{n} s_{n}
\end{array}\right)
$$

with $\operatorname{det}\left(N_{n}(q)\right)=\left(q-q^{-1}\right)^{n} \operatorname{det}\left(N_{n}^{\prime}(q)\right)$. Thus we are reduced to proving that $\operatorname{det}\left(N_{n}^{\prime}(q)\right)=\operatorname{det}\left(M_{n}(q)\right)$. Consider the $(2 n-1) \times(2 n-1)$ matrix $R=\left(\begin{array}{cc}I_{n-1} & 0 \\ C & D\end{array}\right)$ with the lower $n \times(2 n-1)$ submatrix given by

$$
\left(\begin{array}{lll}
C & D
\end{array}\right)=\left(\begin{array}{ccccccccc}
-q_{1} s_{1}-q_{2} s_{2} & \cdots & -q_{n-1} s_{n-1} & 1 & q^{-1} s_{1} & q^{-2} s_{2} & \cdots & q^{2-n} s_{n-2} & q^{1-n} s_{n-1} \\
0 & -q_{1} s_{1} & \cdots & -q_{n-2} s_{n-2} & 0 & 1 & q^{-1} s_{1} & \cdots & q^{3-n} s_{n-3}
\end{array} q^{2-n} s_{n-2}\right)\left(\begin{array}{ccccccccc} 
\\
\vdots & \vdots & \ddots & \vdots & \vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & -q_{1} s_{1} & 0 & 0 & 0 & \cdots & 1 \\
0 & 0 & \cdots & 0 & 0 & 0 & 0 & \cdots & 0
\end{array}\right.
$$

By a careful computation, one verifies without difficulty that $R N_{n}^{\prime}(q)=\left(\begin{array}{cc}U & * \\ 0 & M_{n}(q)\end{array}\right)$, where $U$ is an $(n-1) \times(n-1)$-upper triangular matrix with all diagonal entries equal to 1 . Note that $\operatorname{det}(R)=\operatorname{det}(D)=\operatorname{det}(U)=1$, and it follows immediately that $\operatorname{det}\left(N_{n}^{\prime}(q)\right)=\operatorname{det}\left(M_{n}(q)\right)$.

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# Complex conjugation and Shimura varieties 

Don Blasius and Lucio Guerberoff


#### Abstract

In this paper we study the action of complex conjugation on Shimura varieties and the problem of descending Shimura varieties to the maximal totally real field of the reflex field. We prove the existence of such a descent for many Shimura varieties whose associated adjoint group has certain factors of type $A$ or $D$. This includes a large family of Shimura varieties of abelian type. Our considerations and constructions are carried out purely at the level of Shimura data and group theory.


## 1. Introduction

The goal of this paper is to analyze some aspects of complex conjugation acting on Shimura varieties. This topic has been studied for a long time by several authors, notably Shimura, Deligne, Langlands, Milne, Shih, and more recently Taylor. In general, given a Shimura variety $\operatorname{Sh}(G, X)$ defined by a Shimura datum ( $G, X$ ), and any automorphism $\alpha$ of $\mathbb{C}$, Langlands [1979] conjectured that the conjugate variety $\alpha \operatorname{Sh}(G, X)=\operatorname{Sh}(G, X) \times_{\mathbb{C}, \alpha} \mathbb{C}$ can be realized as a Shimura variety $\operatorname{Sh}\left({ }^{\alpha} G,{ }^{\alpha} X\right)$ for a very explicit pair $\left({ }^{\alpha} G,{ }^{\alpha} X\right)$. This has been proved by Milne [1983] (see also [Borovoĭ 1983; 1987; Milne 1999]). The case of $\alpha=c$ (complex conjugation) has, among other properties, the particularity that the pair $\left({ }^{c} G,{ }^{c} X\right)$ is very concrete. Namely, it can be identified with ( $G, \bar{X}$ ), where $\bar{X}$ is obtained by composing the elements of $x$ with complex conjugation on the Deligne torus $\mathbb{S}$. This simple description is hard to find in the literature, and hence, we include a proof of how it is deduced from the general constructions.

Assuming a few standard extra conditions on the Shimura datum $(G, X)$, the reflex field $E$ can be seen to be either totally real or a CM field. The Shimura variety has a canonical model $\operatorname{Sh}(G, X)_{E}$ over $E$, and the Hecke operators are defined over $E$ as well. In this paper we investigate descent of these varieties to the maximal totally real subfield $E^{+}$of $E$. The existence of such descent can be seen as a nice generalization of the useful fact that the field obtained by adjoining to $\mathbb{Q}$ the $j$-invariant of an order in an imaginary quadratic field has a real embedding. From now on, assume that $E$ is CM. We show in many cases that $\operatorname{Sh}(G, X)$ has

[^15]a model over $E^{+}$. Although the Hecke operators are not defined over $E^{+}$, they can nevertheless be characterized. The general framework for constructing such models comes from the construction of descent data arising from automorphisms of $G$ of order 2 taking $X$ to $\bar{X}$. Using the classification of (adjoint) Shimura data in terms of special nodes on Dynkin diagrams, our aim is to construct an involution of $G$ that induces the opposition involution on the based root datum (or the Dynkin diagram). The construction we make follows from the classification of semisimple groups. The groups $G$ with which we work are, roughly speaking, those for which the simple factors of $G^{\text {ad }}$ are of classical type $A$ or $D$, and satisfy an extra condition on the hermitian or skew-hermitian space defining them (see Definitions 4.1.4 and 4.2.1). For example, a factor of type $A$ is attached to a hermitian space over a central division algebra $D$ over a CM field $K$ endowed with an involution of the second kind $J$. We show that, if there exists an opposition involution on these groups, then $D$ must be either $K$ or a quaternion division algebra, and the involution $J$ is easily described. We carry out the construction of involutions if we assume the aforementioned extra condition, which in this case amounts to the existence of a basis of the underlying vector space such that the matrix of the hermitian form is diagonal with entries in $K$. In the quaternion algebra case, we can write $D=D_{0} \otimes_{F} K$, where $F$ is the maximal totally real subfield of $K$, and $D_{0}$ is a quaternion division algebra over $F$. We assume furthermore in this case that, if $D_{0, v}$ is not split for an embedding $v: F \hookrightarrow \mathbb{R}$, then the corresponding factor of $G_{\mathbb{R}}^{\text {ad }}$ is compact. If $D=K$, the conditions in Definition 4.1.4 are automatically satisfied. For factors of type $D$, there is a similar scenario, although we only restrict to groups of type $D^{\sharp H}$ as in the Appendix of [Milne and Shih 1981]. This encompasses a large family of Shimura varieties of abelian type. Under these assumptions, the existence of the involution on the group $G$ follows from a concrete construction of involutions on each of the simple factors of $G^{\text {ad }}$, which are explicitly given in terms of simple algebras.

To give a flavor of the type of involutions constructed in the paper, suppose that $\mathrm{SU}(V, h)$ is a simple factor of type $A$, corresponding to a hermitian space ( $V, h$ ) of dimension $n$ over a CM field $K$. Let $F$ be the maximal totally real subfield of $K$, and let $\iota$ be the nontrivial automorphism of $K / F$. We take an orthogonal basis $\left\{v_{1}, \ldots, v_{n}\right\}$ of $V$, and we let $I: V \rightarrow V$ be the $\iota$-semilinear map obtained by applying $\iota$ to the coordinates of elements of $V$ with respect to the given basis. Then the map $\theta: \mathrm{SU}(V, h) \rightarrow \mathrm{SU}(V, h)$ given by

$$
\theta(g)=I g I
$$

is an opposition involution.
We stress here that our methods are group-theoretic and we work purely at the level of Shimura data, in the sense that we do not directly make use of a moduli interpretation. However, the methods rely on Langlands conjugation of Shimura
varieties, which in turn is proved using the moduli interpretation in terms of abelian varieties [Milne 1990, §II.9]. An interesting question would be to consider factors of type $E_{6}$, which is the only other type apart from $A$ or $D^{\mapsto H}$ that contributes to the reflex field being CM instead of totally real. We plan to investigate this question in the future.

Let us describe the organization of the paper and outline the main argument. In Section 2, we start by recalling the general formalism of conjugation of Shimura varieties by an arbitrary automorphism of $\mathbb{C}$, we study the special case of complex conjugation explicitly, and we prove in this case that the conjugate Shimura datum is ( $G, \bar{X}$ ), where $\bar{X}$ is the complex conjugate conjugacy class of $X$. We show (Theorem 2.3.1) that, if $(G, X)$ is a Shimura datum and $\theta: G \rightarrow G$ is an involution such that $\theta(X)=\bar{X}$, then $\theta$ induces an isomorphism of algebraic varieties from the complex conjugate $c \operatorname{Sh}(G, X)$ to $\operatorname{Sh}(G, X)$, defined over the reflex field $E$, that constitutes a descent datum from $E$ to $E^{+}$.

In Section 3, we recall some basic facts about root data and opposition involutions, and in Proposition 3.4.8, we lay the ground for the prototype of involutions $\theta: G \rightarrow G$ that we will construct. Roughly speaking, suppose that $T \subset G$ is a maximal torus of $G$, and $x \in X$ factors through $T_{\mathbb{R}}$. If $\theta: G \rightarrow G$ is an involution that preserves $T_{\mathbb{R}}$ and induces complex conjugation on the group of characters $X^{*}(T)$, then $\theta(x)=\bar{x}$ and thus $\theta(X)=\bar{X}$. This is basically the type of involution that we will construct, with some slight changes. Since we will make use of the explicit classification of semisimple groups, we need to work with either $G^{\text {der }}$ or $G^{\text {ad }}$. We let $G_{i}$ be the almost simple factors of $G^{\text {der }}$, and $\widetilde{G}_{i}$ be their simply connected covers, so that $\widetilde{G}_{i}=\operatorname{Res}_{F_{i} / \mathbb{Q}} H_{i}$, for certain groups $H_{i}$ which are absolutely almost simple, simply connected, over a totally real field $F_{i}$. We recall the classification of these groups in Section 4, where we also construct opposition involutions on them preserving specific maximal tori $S_{i}$ and inducing complex conjugation on their characters (for noncompact places $v$ of $F_{i}$ ). We only do this for groups of type $A$ or $D^{\mathfrak{H}}$. These, together with type $E_{6}$, are the only ones that give a CM reflex field, as opposed to totally real. Furthermore, as noted above, we impose some extra conditions in order to construct the involutions. From the tori $S_{i}$, we get maximal tori $T^{\prime} \subset G^{\text {der }}$ and $T \subset G$, and an opposition involution $\theta^{\prime}: G^{\mathrm{der}} \rightarrow G^{\mathrm{der}}$ preserving $T^{\prime}$. As shown in Proposition 3.4.8, $\theta^{\prime}$ extends uniquely to an involution on $G$. To show that $\theta(X)=\bar{X}$, we need to relate in some way the choice of our tori $S_{i}$, which is a priori unrelated to the Shimura datum, to the conjugacy class $X$. In Section 5, we show that there always exists $x \in X$ such that $x^{\text {ad }}$ factors through the image of $T_{\mathbb{R}}$ in $G_{\mathbb{R}}^{\text {ad. }}$. This is all we need for Proposition 3.4.8. In Theorem 5.2.2, we state the existence of descent datum for Shimura varieties defined by groups $(G, X)$ such that the simple factors of $G^{\text {ad }}$ are of the type described in Section 4. We call these strongly of type $\left(A D^{\sharp}\right)$. Finally, we also note that involutions inducing the desired
descent datum on $\operatorname{Sh}(G, X)$ can be constructed whenever $G$ is adjoint and there exists an opposition involution $\theta: G \rightarrow G$. This is always the case if $G$ is quasisplit, for example.

The existence of the involutions constructed in this paper should have interesting applications, which will be explored in the future, for example, in the setting of integral models and the zeta function problem, and periods of automorphic forms.

Notation and conventions. We fix an algebraic closure $\mathbb{C}$ of the real numbers $\mathbb{R}$ and a choice of $i=\sqrt{-1}$, and we let $\overline{\mathbb{Q}}$ denote the algebraic closure of $\mathbb{Q}$ in $\mathbb{C}$. We let $c \in \operatorname{Gal}(\mathbb{C} / \mathbb{R})$ denote complex conjugation on $\mathbb{C}$, and we use the same letter to denote its restriction to $\overline{\mathbb{Q}}$. Sometimes we also write $c(z)=\bar{z}$ for $z \in \mathbb{C}$.

Let $k$ be a field. By a variety over $k$ we will mean a geometrically reduced scheme of finite type over $k$. We let $\mathbb{G}_{\mathrm{m}, k}$ denote the usual multiplicative group over $k$. For any algebraic group $G$ over $k$, we let $\operatorname{Lie}(G)$ denote its Lie algebra. For us, a reductive group will always be connected. If $G$ is reductive, we let $G^{\text {ad }}$ (resp. $G^{\text {der }}$ ) denote its adjoint group $G / Z(G)$ (resp. its derived subgroup), where $Z(G)$ is the center of $G$. We let $G^{\mathrm{ab}}=G / G^{\mathrm{der}}$ (a torus). If $T \subset G$ is a torus, we denote by $T^{\text {ad }}$ the image of $T$ under the projection $G \rightarrow G^{\text {ad }}$. For any commutative group scheme $G$, we denote by $\operatorname{inv}_{G}: G \rightarrow G$ the map $g \mapsto g^{-1}$.

We denote by $\mathbb{A}\left(\right.$ resp. $\left.\mathbb{A}_{f}\right)$ the ring of adèles of $\mathbb{Q}$ (resp. finite adèles). A CM field $K$ is a totally imaginary quadratic extension of a totally real field $F$.

We let $\mathbb{S}=R_{\mathbb{C} / \mathbb{R}} \mathbb{G}_{\mathrm{m}, \mathbb{C}}$. We denote by $c=c_{\mathbb{S}}$ the algebraic automorphism of $\mathbb{S}$ induced by complex conjugation. For any $\mathbb{R}$-algebra $A$, this is $c \otimes_{\mathbb{R}} \mathrm{id}_{A}$ : $\left(\mathbb{C} \otimes_{\mathbb{R}} A\right)^{\times} \rightarrow\left(\mathbb{C} \otimes_{\mathbb{R}} A\right)^{\times}$on the points of $\mathbb{S}(A)$. This is often denoted by $z \mapsto \bar{z}$, and on complex points it should not be confused with the other complex conjugation $\mathrm{id}_{\mathbb{C}} \otimes c$ on $\mathbb{S}(\mathbb{C})=\left(\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C}\right)^{\times}$on the second coordinate.

An involution of a group is an automorphism of order 2, whereas an involution of a ring is an antiautomorphism of order 2. This should not cause any confusion.

We will denote by $\mathbb{H}$ the nonsplit quaternion algebra over $\mathbb{R}$, identified with the set of matrices of the form

$$
\left(\begin{array}{rr}
x & y \\
-\bar{y} & \bar{x}
\end{array}\right)
$$

in $M_{2}(\mathbb{C})$.

## 2. Shimura varieties, conjugation, and descent

We will first review some basic facts about Shimura varieties and conjugation by an automorphism of $\mathbb{C}$, specializing to the case of complex conjugation. Then we set up our descent problem, describe some general considerations about reflex fields and Dynkin diagrams, and explain how to construct descent data based on involutions of a Shimura datum.
2.1. Shimura varieties. A Shimura datum $(G, X)$ will be understood in the sense of Deligne's axioms [1979, (2.1.1.1-3)]. We will assume moreover that the connected component $Z^{0}$ of the center $Z$ of $G$ splits over a CM field. For a compact open subgroup $K \subset G\left(\mathbb{A}_{f}\right)$, we put $\operatorname{Sh}_{K}(G, X)(\mathbb{C})=G(\mathbb{Q}) \backslash X \times G\left(\mathbb{A}_{f}\right) / K$. For $K$ sufficiently small (which we assume from now on), this complex analytic space is smooth and is equal to the complex points of a complex quasiprojective variety $\operatorname{Sh}_{K}(G, X)_{\mathbb{C}}$. Let $E=E(G, X) \subset \mathbb{C}$ be the reflex field of $(G, X)$; under our hypotheses, this is contained in a CM field, and thus, it is either a CM field or a totally real field. In any case, we let $E^{+}$be the maximal totally real subfield of $E$. The variety $\operatorname{Sh}_{K}(G, X)_{\mathbb{C}}$ admits a canonical model over $E$, denoted by $\operatorname{Sh}_{K}(G, X)_{E}$. We use the same notation for the pro-objects $\operatorname{Sh}(G, X)(\mathbb{C}), \operatorname{Sh}(G, X)_{\mathbb{C}}$, and $\operatorname{Sh}(G, X)_{E}$. We denote by $w_{X}: \mathbb{G}_{\mathrm{m}, \mathbb{R}} \rightarrow G_{\mathbb{R}}$ the composition of $x \in X$ with the weight morphism $w: \mathbb{G}_{\mathrm{m}, \mathbb{R}} \rightarrow \mathbb{S}$, for some (or any) $x \in X$, and call it the weight morphism of $(G, X)$. For $x \in X$, we let $\mu_{x}: \mathbb{G}_{\mathrm{m}, \mathbb{C}} \rightarrow G_{\mathbb{C}}$ be the map given by $\mu_{x}(z)=x_{\mathbb{C}}(z, 1)$, under the identification of $\mathbb{S}_{\mathbb{C}} \cong \mathbb{G}_{\mathrm{m}, \mathbb{C}} \times \mathbb{G}_{\mathrm{m}, \mathbb{C}}$ given by $(z \otimes a) \mapsto(z a, \bar{z} a)$.

We will fix the following notation once and for all. Let $p: G \rightarrow G^{\text {ad }}$ be the projection onto $G^{\text {ad }}$. The natural isogeny $Z^{0} \times G^{\text {der }} \rightarrow G$ and the projection $G \rightarrow G^{\mathrm{ab}}$ define an isogeny $Z^{0} \rightarrow G^{\mathrm{ab}}$. Let $G_{1}, \ldots, G_{r}$ be the almost simple factors of $G^{\text {der }}$ over $\mathbb{Q}$, and let $\widetilde{G}_{i} \rightarrow G_{i}$ be their simply connected covers. We can write $\widetilde{G}_{i}=\operatorname{Res}_{F_{i} / \mathbb{Q}} H_{i}$, where the fields $F_{i}$ are totally real and the groups $H_{i}$ are simply connected, absolutely almost simple over $F_{i}$. For each embedding $v \in I_{i}=\operatorname{Hom}\left(F_{i}, \mathbb{C}\right)$, we have groups $H_{i, v}=H_{i} \otimes_{F_{i}, v} \mathbb{R}$, and for a fixed $i=1, \ldots, r$, all these groups have the same Dynkin type $D_{i}$, which will be called the Dynkin type of $\widetilde{G}_{i}$ (or of $G_{i}$ or $H_{i}$ ). We let $I_{i, c}=\left\{v \in I_{i}: H_{i, v}^{\text {ad }}(\mathbb{R})\right.$ is compact $\}$ and we let $I_{i, n c}$ be its complement in $I_{i}$, which must be nonempty if $H_{i}$ is nontrivial. We also have that $G^{\text {ad }}$ is the direct product of the $G_{i}^{\text {ad }}=\operatorname{Res}_{F_{i} / \mathbb{Q}} H_{i}^{\text {ad }}$, and $G_{\mathbb{R}}^{\text {ad }}$ is the direct product of the $H_{i, v}^{\text {ad }}$ for $i=1, \ldots, r$ and $v \in I_{i}$. Let $X^{\text {ad }}$ be the $G^{\text {ad }}(\mathbb{R})$-conjugacy class containing $p_{\mathbb{R}}(X)$, and write $X^{\text {ad }}=\prod_{i, v} X_{i, v}$ with $X_{i, v}$ an $H_{i, v}^{\text {ad }}(\mathbb{R})$-conjugacy class of morphisms $\mathbb{S} \rightarrow H_{i, v}^{\text {ad }}$. For each $i$ and each $v \in I_{i, n c}$, there is a special node $s_{i, v}$ in the Dynkin diagram $D_{i, v}$ of $H_{i, v}$ attached to $X_{i, v}$, which uniquely determines $X_{i, v}$ as a conjugacy class with target $H_{i, v}^{\text {ad }}$ (in the sense that if $Y$ is an $H_{i, v}^{\text {ad }}(\mathbb{R})$-conjugacy class satisfying Deligne's axioms, for which its associated special node is $s_{i, v}$, then $Y=X_{i, v}$ [Deligne 1979, §1.2.6]).
2.2. Conjugation. For the general properties of conjugation of Shimura varieties, we mainly follow [Milne 1990; Milne and Shih 1982b; Deligne 1982; Milne and Shih 1982a]; see also [Langlands 1979]. Let ( $G, X$ ) be a Shimura datum. A special pair $(T, x)$ consists of a maximal torus $T \subset G$ and a point $x \in X$ factoring through $T_{\mathbb{R}}$. Fix $x \in X$ a special point, and let $\sigma \in \operatorname{Aut}(\mathbb{C})$. We denote by $\left({ }^{\sigma, x} G,{ }^{\sigma, x} X\right)$ the conjugate Shimura datum. We recall its construction below. By Theorem II.4.2 of
[Milne 1990], there exists a unique isomorphism

$$
\varphi_{\sigma, x}: \sigma \operatorname{Sh}(G, X)_{\mathbb{C}}=\operatorname{Sh}(G, X)_{\mathbb{C}} \times_{\mathbb{C}, \sigma} \mathbb{C} \simeq \operatorname{Sh}\left({ }^{\sigma, x} G,{ }^{\sigma, x} X\right)_{\mathbb{C}}
$$

satisfying certain conditions. Choosing a different special point gives canonically isomorphic results [Milne 1990, Proposition II.4.3]. The reflex field of $\left({ }^{\sigma, x} G,{ }^{\sigma, x} G\right)$ is $\sigma(E)$, and $\varphi_{\sigma, x}$ identifies $\sigma \operatorname{Sh}(G, X)_{E}=\operatorname{Sh}(G, X)_{E} \times_{E, \sigma} \sigma(E)$ with the canonical model of $\operatorname{Sh}\left({ }^{\sigma, x} G,{ }^{\sigma, x} G\right)_{\mathbb{C}}$ over $\sigma(E)$ [Milne 1990, Theorem II.5.5]. In particular, if $\sigma(E)=E$, then $\varphi_{\sigma, x}$ defines an isomorphism

$$
\varphi_{\sigma, x}: \sigma \operatorname{Sh}(G, X)_{E} \simeq \operatorname{Sh}\left({ }^{\sigma, x} G,{ }^{\sigma, x} X\right)_{E}
$$

over $E$. All of this also works at finite level: if $K \subset G\left(\mathbb{A}_{f}\right)$ is compact open, $\varphi_{\sigma, x}$ sends $\sigma \operatorname{Sh}_{K}(G, X)_{\mathbb{C}}$ to $\operatorname{Sh}_{\sigma, x}\left({ }^{\sigma, x} G,{ }^{\sigma, x} X\right)_{\mathbb{C}}$ (same thing replacing $\mathbb{C}$ by $E$ and $\sigma(E)$ ), where ${ }^{\sigma, x} K \subset{ }^{\sigma, x} G\left(\mathbb{A}_{f}\right)$ is explicit (see below).

We are interested mainly in the case $\sigma=c$, but nevertheless it will be useful to recall the general construction of $\left({ }^{\sigma, x} G,{ }^{\sigma, x} X\right)$. Let $\mathfrak{S}$ be the (connected) Serre group. This can be defined as the group of automorphisms of the forgetful fiber functor from the Tannakian category of $\mathrm{CM} \mathbb{Q}$-Hodge structures to the category of finite-dimensional $\mathbb{Q}$-vector spaces. (Here a $\mathbb{Q}$-Hodge structure is a $\mathbb{Q}$-vector space $V$ such that $V \otimes \mathbb{C}$ is endowed with a Hodge structure; the structure is CM if the algebra of elements of $\operatorname{End}(V)$ which induce morphisms of Hodge structure contains a commutative semisimple subalgebra of dimension $\operatorname{dim}_{\mathbb{Q}}(V)$.) Let $\mathfrak{T}$ denote the Taniyama group, defined here as the group of automorphisms of the Betti fiber functor in Deligne's Tannakian category of CM motives for absolute Hodge cycles over $\mathbb{Q}$; this is the Tannakian category generated by Artin motives and by the cohomology of abelian varieties over $\mathbb{Q}$ which are potentially CM. These are proalgebraic groups, and there is a natural exact sequence

$$
1 \rightarrow \mathfrak{S} \rightarrow \mathfrak{T} \xrightarrow{\pi} \operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q}) \rightarrow 1
$$

where the second arrow corresponds to the functor taking a CM motive $M$ to its CM Hodge structure $H_{B}(M)$, and $\pi$ corresponds to the natural inclusion of the category of Artin motives into the category of CM motives. The group Gal( $\overline{\mathbb{Q}} / \mathbb{Q})$ is to be considered as the proalgebraic group given by the inverse limit of the finite constant groups $\operatorname{Gal}(L / \mathbb{Q})$, for $L \subset \mathbb{C}$ a finite Galois extension of $\mathbb{Q}$. There is a continuous section of $\pi$ over $\mathbb{A}_{f}$ denoted by

$$
\operatorname{sp}: \operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q}) \rightarrow \mathfrak{T}\left(\mathbb{A}_{f}\right)
$$

For a motive $M, \operatorname{sp}(\sigma)$ corresponds to the automorphism of $H_{B}(M) \otimes_{\mathbb{Q}} \mathbb{A}_{f}$ obtained from the Galois action of $\sigma$ on étale cohomology using the comparison isomorphism.

Suppose that $\sigma=c$ is complex conjugation. Then $\operatorname{sp}(c) \in \mathfrak{T}\left(\mathbb{A}_{f}\right)$ can be described as follows, as explained in [Deligne 1982, Lemme 5]. Suppose that $M$ is a CM motive over $\mathbb{Q}$, realized as the cohomology of an algebraic variety $X$ over $\mathbb{Q}$. The action of complex conjugation on $X(\mathbb{C})$ induces an involution $F$ on the Betti realization $H_{B}(M)=H^{i}(X(\mathbb{C}), \mathbb{Q})$. The automorphism $\operatorname{sp}(c): H_{B}(M) \otimes \mathbb{A}_{f} \rightarrow H_{B}(M) \otimes \mathbb{A}_{f}$ is then equal to $F \otimes_{\mathbb{Q}} \mathrm{id}_{\mathbb{A}_{f}}$. This implies, in particular, that $\operatorname{sp}(c) \in \mathfrak{T}(\mathbb{Q})$.

For any $\sigma \in \operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q})$, we let ${ }^{\sigma} \mathfrak{S}=\pi^{-1}(\sigma)$. There is a cocharacter $\mu_{\text {can }}$ : $\mathbb{G}_{\mathrm{m}, \mathbb{C}} \rightarrow \mathfrak{S}_{\mathbb{C}}$, which in Tannakian terms gives rise to the Hodge cocharacter of the Hodge structures on $H_{B}(M) \otimes_{\mathbb{Q}} \mathbb{C}$.

Let $G$ be any algebraic group over $\mathbb{Q}$ and $\rho: \mathfrak{S} \rightarrow G^{\text {ad }}$ be a homomorphism, inducing an action of $\mathfrak{S}$ on $G$ by group automorphisms (conjugation). Let ${ }^{\sigma, \rho} G=$ ${ }^{\sigma} \mathfrak{S} \times{ }_{\mathfrak{S}, \rho} G$ be the group obtained by twisting $G$ by the torsor ${ }^{\sigma} \mathfrak{S}$. Thus, ${ }^{\sigma, \rho} G$ is the fpqc sheaf associated with the presheaf sending a $\mathbb{Q}$-algebra $R$ to the group ${ }^{\sigma} \mathfrak{S}(R) \times{ }_{\mathfrak{S}(R), \rho} G(R)$, which is the quotient of ${ }^{\sigma} \mathfrak{S}(R) \times G(R)$ by the right action $(s, g) s_{1}=\left(s s_{1}, s_{1}^{-1} g\right)$ of $\mathfrak{S}(R)$. The class of an element $(s, g)$ in this quotient will be denoted by $s \cdot g$.

Lemma 2.2.1. Keep the notation and assumptions as above, with $\sigma=c$. There exists a natural isomorphism ${ }^{c, \rho} G \rightarrow G$.

Proof. As explained above, $\operatorname{sp}(c) \in{ }^{c} \mathfrak{S}(\mathbb{Q})$, so ${ }^{c} \mathfrak{S}$ is trivialized over $\mathbb{Q}$. In particular, the map $\operatorname{sp}(c)_{R} \cdot g \mapsto g$ (for $g \in G(R)$ ) defines a group isomorphism between the presheaves defining ${ }^{c, \rho} G$ and $G$. A fortiori, this defines an isomorphism ${ }^{c, \rho} G \rightarrow G$.

Remark 2.2.2. If $H \subset G$ is a subgroup on which $\mathfrak{S}$ acts trivially, then ${ }^{c, \rho} H$ is canonically isomorphic to $H$ (this is true for any $\sigma$ ). This identification is compatible with that of Lemma 2.2.1.

Remark 2.2.3. In [Milne 1990, §II.4], an isomorphism $G\left(\mathbb{A}_{f}\right) \rightarrow^{c, \rho} G\left(\mathbb{A}_{f}\right)$ is constructed, which is denoted by $g \mapsto{ }^{c} g$. When identifying ${ }^{c, \rho} G$ with $G$ using Lemma 2.2.1, this becomes the identity map. A similar remark applies to the isomorphism $g \mapsto{ }^{c} g$ between $G_{\mathbb{C}}$ and ${ }^{c, \rho} G_{\mathbb{C}}$ defined in [Milne 1990, §III.1] (note that the element $z_{\infty}(c)$ defined in [op. cit.] is equal to $\left.\operatorname{sp}(c)_{\mathbb{C}}\right)$.

Suppose that $(G, X)$ is a Shimura datum as before, and $(T, x)$ is a special pair. The map $\mu_{x}$ factors through $T_{\mathbb{C}}$, and there exists a unique homomorphism

$$
\rho_{x}^{\mathrm{ad}}: \mathfrak{S} \rightarrow G^{\mathrm{ad}}
$$

such that $\left(\rho_{x}^{\text {ad }}\right)_{\mathbb{C}} \circ \mu_{\text {can }}=\mu_{x}^{\text {ad }}$. For $\sigma \in \operatorname{Aut}(\mathbb{C})$, the group ${ }^{\sigma, x} G$ is defined to be $\sigma, \rho_{x}^{\text {ad }} G$ in the previous notation (where we take the restriction of $\sigma$ to $\overline{\mathbb{Q}}$ ). Since the cocharacter $\sigma\left(\mu_{x}\right)$ of $T={ }^{\sigma, \rho_{x}} T$ commutes with its complex conjugate, it is the

Hodge cocharacter associated with a map $\mathbb{S} \rightarrow{ }^{\sigma, x} G_{\mathbb{R}}$ which we denote by

$$
{ }^{\sigma} x: \mathbb{S} \rightarrow{ }^{\sigma, x} G_{\mathbb{R}}
$$

Finally, ${ }^{\sigma, x} X$ is defined to be the ${ }^{\sigma, x} G(\mathbb{R})$-conjugacy class of ${ }^{\sigma} x$.
Assume now that $\sigma=c$. By Lemma 2.2.1, we can identify ${ }^{c, x} G$ with $G$, and hence, we can see ${ }^{c} x: \mathbb{S} \rightarrow^{c, x} G_{\mathbb{R}}$ as a map

$$
{ }^{c} x: \mathbb{S} \rightarrow G_{\mathbb{R}}
$$

with target $G_{\mathbb{R}}$. Let $c=c_{\mathbb{S}}: \mathbb{S} \rightarrow \mathbb{S}$ be complex conjugation on $\mathbb{S}$. For any $h: S \rightarrow G_{\mathbb{R}}$, let $\bar{h}=h \circ c$.

Lemma 2.2.4. In the notation above, we have that

$$
{ }^{c} x=\bar{x}
$$

Proof. It is enough to show that $\bar{x}(\mathbb{C})=\left({ }^{c} x\right)(\mathbb{C}): \mathbb{S}(\mathbb{C}) \rightarrow G(\mathbb{C})$. Recall that we are identifying $\mathbb{S}(\mathbb{C})=\left(\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C}\right)^{\times}$with $\mathbb{C}^{\times} \times \mathbb{C}^{\times}$via the map $(z \otimes a) \mapsto(z a, \bar{z} a)$. Then $c=c_{\mathbb{S}}: \mathbb{S}(\mathbb{C}) \rightarrow \mathbb{S}(\mathbb{C})$, which is given by $c(z \otimes a)=\bar{z} \otimes a$, becomes the map $(a, b) \mapsto(b, a)$. This is an algebraic automorphism of $\mathbb{S}$. There is another complex conjugation, which will be denoted by $c^{\prime}$ here, on the complex points of $\mathbb{S}$. Namely,

$$
c^{\prime}: S(\mathbb{C}) \rightarrow \mathbb{S}(\mathbb{C})
$$

which is induced by complex conjugation on $\mathbb{C}$. It is given by $c^{\prime}(z \otimes a)=z \otimes \bar{a}$. Then, as a map on $\mathbb{C}^{\times} \times \mathbb{C}^{\times}$, it is given by $(a, b) \mapsto(\bar{b}, \bar{a})$.

Recall that $\mu_{x}(a)=x(\mathbb{C})(a, 1)$ for $a \in \mathbb{C}^{\times}$. For readability purposes, we use the notation $x_{\mathbb{C}}(a, b)$ instead of $x(\mathbb{C})(a, b)$ in what follows. Then, for $a, b \in \mathbb{C}^{\times}$, we have that

$$
\begin{equation*}
\bar{x}_{\mathbb{C}}(a, b)=x_{\mathbb{C}}(b, a)=x_{\mathbb{C}}(b, 1) x_{\mathbb{C}}(1, a) \tag{2.2.5}
\end{equation*}
$$

Similarly,

$$
\begin{equation*}
\left({ }^{c} x\right)_{\mathbb{C}}(a, b)=\mu_{c_{x}}(a) \mu_{\bar{c}_{x}}(b) \tag{2.2.6}
\end{equation*}
$$

Now, by definition, $\mu_{c_{x}}=c\left(\mu_{x}\right)$, where $c$ now denotes the action on cocharacters. If we let $g \mapsto \bar{g}$ denote the map on $G(\mathbb{C})$ induced by $c: \mathbb{C} \rightarrow \mathbb{C}$, then

$$
\begin{equation*}
\mu_{c_{x}}(a)=\overline{\mu_{x}(\bar{a})}=\overline{x_{\mathbb{C}}(\bar{a}, 1)} \tag{2.2.7}
\end{equation*}
$$

for $a \in \mathbb{C}^{\times}$. Since $x$ is defined over $\mathbb{R}$, it commutes with the maps on complex points induced by $c: \mathbb{C} \rightarrow \mathbb{C}$. That is, the diagram

is commutative. From this and (2.2.7), it follows that

$$
\begin{equation*}
\mu_{c_{x}}(a)=x_{\mathbb{C}}(1, a) . \tag{2.2.8}
\end{equation*}
$$

Similarly,

$$
\mu_{\bar{c}_{x}}(b)=\bar{c}_{\mathbb{C}}(b, 1)=\left({ }^{c} x\right)_{\mathbb{C}}(1, b)=\mu_{c\left({ }_{c} x\right)}(b) .
$$

But $\mu_{c\left({ }_{( }{ }_{x}\right)}=c\left(\mu_{c_{x}}\right)=\mu_{x}$, so this equals $x_{\mathbb{C}}(b, 1)$. The proof finishes by combining (2.2.5), (2.2.6), and (2.2.8).

By the last lemma, we can identify ${ }^{c, x} X$ with the $G(\mathbb{R})$-conjugacy class of $\bar{x}$. The map $h \mapsto \bar{h}$ defines an antiholomorphic isomorphism between $X$ and ${ }^{c, x} X$. This does not depend on $x$, and from now on we let

$$
\bar{X}=\{\bar{h}: h \in X\} .
$$

Thus, the pair ( ${ }^{c, x} G,{ }^{c, x} X$ ) becomes naturally identified with the pair $(G, \bar{X})$. The isomorphism $\varphi_{c, x}$ becomes, under this identification, an isomorphism $\varphi$ : $\operatorname{Sh}_{K}(G, X)_{E} \times_{E, c} E \rightarrow \operatorname{Sh}_{K}(G, \bar{X})_{E}$. On complex points, it defines an antiholomorphic isomorphism between $\operatorname{Sh}_{K}(G, X)(\mathbb{C})$ and $\operatorname{Sh}_{K}(G, \bar{X})(\mathbb{C})$, which we denote by $\phi$. For $[h, g] \in \operatorname{Sh}_{K}(G, X)(\mathbb{C})$, we have that $\phi([h, g])=[\bar{h}, g] \in \operatorname{Sh}_{K}(G, \bar{X})(\mathbb{C})$.

For example, suppose that $E \subset \mathbb{R}$. Then there is an antiholomorphic involution on $\mathrm{Sh}_{K}(G, X)(\mathbb{C})$ defined by complex conjugation acting on $\mathbb{C}$. It follows from the theory of canonical models that this involution takes the form $[h, g] \mapsto[\eta(h), g]$, where $\eta: X \rightarrow X$ is an antiholomorphic involution of the form $\eta(g \cdot x)=(g n) \cdot x$ for some $n \in N(\mathbb{R})$ (here $N$ is the normalizer in $G$ of $T$ ). See [Milne 1990, §II.7] for details. In fact, the theory implies that there exists $n \in N(\mathbb{R})$ such that ${ }^{c} x=n \cdot x$, and thus, $\bar{X}=X$. Then the map $\eta$ becomes what we called $\phi$; that is, $\eta(h)=\bar{h}$ for any $h \in X$.
2.3. Involutions of Shimura data and descent. Fix a Shimura datum $(G, X)$, with reflex field $E$. For an involution $\theta: G \rightarrow G$, let $\theta(X)$ be the $G(\mathbb{R})$-conjugacy class $\{\theta(h): h \in X\}$, where $\theta(h)=\theta_{\mathbb{R}} \circ h$. Since we want to consider involutions $\theta$ that send $X$ to $\bar{X} \neq X$, from now on, we will focus on the case where $E$ is a CM field (if $E$ is totally real, the identity map on $G$ takes $X$ to $\bar{X}$ ). Let $E^{+} \subset E$ be the maximal totally real subfield, and let $\iota \in \operatorname{Gal}\left(E / E^{+}\right)$be the nontrivial automorphism, i.e., the restriction of complex conjugation $c$ to $E$.

Suppose that $\theta$ is an involution of $G$ such that $\theta(X)=\bar{X}$. For a compact open subgroup $K \subset G\left(\mathbb{A}_{f}\right)$, denote by ${ }^{\theta} K=\theta(K) \subset G\left(\mathbb{A}_{f}\right)$. Then $\theta$ induces an isomorphism of algebraic varieties $\operatorname{Sh}(\theta): \operatorname{Sh}_{K}(G, X)_{E} \rightarrow \operatorname{Sh}_{\theta_{K}}(G, \bar{X})_{E}$. On complex points, this takes $[h, g]$ to $\left[\theta_{\mathbb{R}} \circ h, \theta(g)\right]$. Suppose that ${ }^{\theta} K=K$. Then $\operatorname{Sh}(\theta)^{-1} \circ \varphi$ defines an isomorphism $\psi: \iota\left(\operatorname{Sh}_{K}(G, X)_{E}\right)=\operatorname{Sh}_{K}(G, X)_{E} \times_{E, \iota} E \rightarrow$ $\operatorname{Sh}_{K}(G, X)_{E}$.

Let $V$ be an arbitrary scheme over $E$. Recall that an $E / E^{+}$-descent datum is a pair of isomorphisms $\psi_{\mathrm{id}}: \operatorname{id}(V)=V \times_{E, \text { id }} E \rightarrow V$ and $\psi_{\iota}: \iota V=V \times_{E, \iota} E \rightarrow V$ of schemes over $E$ satisfying the cocycle condition

$$
\psi_{\sigma} \circ \sigma\left(\psi_{\tau}\right)=\psi_{\sigma \tau}
$$

for all $\sigma, \tau \in \operatorname{Gal}\left(E / E^{+}\right)$, using the natural identification $\sigma(\tau(V))=(\sigma \tau) V$. Then necessarily $\psi_{\text {id }}$ is the first projection $\operatorname{id}(V) \rightarrow V$, and thus, to give a descent datum amounts to give an isomorphism $\psi=\psi_{l}: \iota(V) \rightarrow V$ such that $\psi \circ \iota(\psi): \iota(\iota(V)) \rightarrow V$ is equal to the identity map, when identifying $\iota(\iota(V))=V$. By definition, such a descent datum is effective if there exists a scheme $V_{0}$ over $E^{+}$and an isomorphism $m: V \rightarrow V_{0, E}=V_{0} \times_{E^{+}} E$ such that $m \circ \psi=\iota(m)$, after identifying $\iota\left(V_{0, E}\right)=V_{0, E}$. If $V$ is a quasiprojective algebraic variety, then any descent datum for $V$ is effective. This was first proved by Weil [1956]. For a modern reference, see [Bosch et al. 1990, §6.2].
Theorem 2.3.1. The map $\psi: \iota\left(\operatorname{Sh}_{K}(G, X)_{E}\right) \rightarrow \operatorname{Sh}_{K}(G, X)_{E}$ obtained as above from an involution $\theta: G \rightarrow G$ such that $\theta(X)=\bar{X}$ and ${ }^{\theta} K=K$ is an effective $E / E^{+}$-descent datum on the Shimura variety $\operatorname{Sh}_{K}(G, X)_{E}$. Hence, there exists a quasiprojective, smooth, algebraic variety $\operatorname{Sh}_{K}(G, X)_{E^{+}}$over $E^{+}$, and an isomorphism $m: \operatorname{Sh}_{K}(G, X)_{E} \rightarrow \operatorname{Sh}_{K}(G, X)_{E^{+}} \times_{E^{+}} E$ such that $m \circ \psi=\iota(m)$.
Proof. Let $V=\operatorname{Sh}_{K}(G, X)_{E}$ and $\bar{V}=\operatorname{Sh}_{K}(G, \bar{X})_{E}$, and let $n: V \rightarrow \iota(\iota(V))$ be the natural isomorphism. We need to check that $\psi \circ \iota(\psi) \circ n=\mathrm{id}_{V}$, and for this it is enough to see that both morphisms are equal on the set of complex points $V(\mathbb{C})$. Let $c_{V}: V(\mathbb{C}) \rightarrow(\iota V)(\mathbb{C})$ be the bijection that sends $x: \operatorname{Spec}(\mathbb{C}) \rightarrow V$ to $p_{\iota, V}^{-1} \circ x \circ \operatorname{Spec}(c)$, where $p_{\iota, V}: \iota V \rightarrow V$ is the first projection, and define $c_{l V}:(\iota V)(\mathbb{C}) \rightarrow(\iota(\iota V))(\mathbb{C})$ similarly. Then we have that $n(\mathbb{C})=c_{l V} \circ c_{V}, \iota(\psi)(\mathbb{C})=$ $c_{V} \circ \psi(\mathbb{C}) \circ c_{l V}^{-1}$, and $\psi$ satisfies that $\psi(\mathbb{C}) \circ c_{V}=\operatorname{Sh}(\theta)^{-1}(\mathbb{C}) \circ \phi$. Recall that $\phi: V(\mathbb{C}) \rightarrow \bar{V}(\mathbb{C})$ sends $[h, g]$ to $[\bar{h}, g]$. Putting all this together, we get that

$$
(\psi \circ \iota(\psi) \circ n)(\mathbb{C})=\operatorname{Sh}(\theta)^{-1}(\mathbb{C}) \circ \phi \circ \operatorname{Sh}(\theta)^{-1}(\mathbb{C}) \circ \phi
$$

and thus,

$$
\left.(\psi \circ \iota(\psi) \circ n)(\mathbb{C})([h, g])=\left[\theta^{-1}\left(\overline{\theta^{-1}(\bar{h}}\right)\right), \theta^{-2}(g)\right]
$$

But for any $y \in \bar{X}, \theta^{-1}(y)=\theta_{\mathbb{R}}^{-1} \circ y$, and so
$\theta^{-1}\left(\overline{\theta^{-1}(\bar{h})}\right)=\theta^{-1}\left(\overline{\theta_{\mathbb{R}}^{-1} \circ \bar{h}}\right)=\theta^{-1}\left(\overline{\theta_{\mathbb{R}}^{-1} \circ h \circ c}\right)=\theta^{-1}\left(\theta_{\mathbb{R}}^{-1} \circ h\right)=\theta_{\mathbb{R}}^{-2} \circ h=\theta^{-2}(h)$, and thus, $(\psi \circ \iota(\psi) \circ n)(\mathbb{C})([h, g])=[h, g]$, using the fact that $\theta^{2}=$ id. Finally, since $\mathrm{Sh}_{K}(G, X)_{E}$ is quasiprojective, the descent datum just constructed is effective.

Remark 2.3.2. The model of Theorem 2.3.1 depends on the descent datum, which in turns depends on the particular involution $\theta$.

We note that, by the nature of the descent datum, Hecke operators do not descend to the model $\operatorname{Sh}_{K}(G, X)_{E^{+}}$. Given $q \in G\left(\mathbb{A}_{f}\right)$, the Hecke operator $T_{q}$ is a morphism of algebraic varieties $T_{q}: \mathrm{Sh}_{K}(G, X)_{E} \rightarrow \mathrm{Sh}_{q^{-1} K q}(G, X)_{E}$, which on complex points is given by $T_{q}([h, g])=[h, g q]$. Then $T_{\theta(q)} \circ \operatorname{Sh}(\theta)=\operatorname{Sh}(\theta) \circ T_{q}$ : $\operatorname{Sh}_{K}(G, X)_{E} \rightarrow \operatorname{Sh}_{\theta(q)^{-1 \theta} K \theta(q)}(G, \bar{X})_{E}$. The Hecke operator $T_{q}$ descends to a map $\mathrm{Sh}_{K}(G, X)_{E^{+}} \rightarrow \mathrm{Sh}_{q^{-1} K q}(G, X)_{E^{+}}$if and only if $T_{\theta(q)}=T_{q}$.

In the following sections we will construct several examples of involutions $\theta$ as above, and explain a general framework for such constructions.

## 3. Opposition involutions

In this section we recall some basic facts about opposition involutions and prove a few results that will be needed in the forthcoming sections. For the basic facts regarding root data, see [Springer 1979].
3.1. Root data. Let $\Psi=\left(X, \Phi, X^{\vee}, \Phi^{\vee}\right)$ be a root datum with $\Phi \neq \varnothing$. Let $Q$ be the subgroup of $X$ generated by $\Phi$, and $V=Q \otimes_{Z} \mathbb{Q}$. Let $W=W(\Phi)$ be the Weyl group of the root system $\Phi$ in $V$. This can be naturally identified with the Weyl group of $\Phi^{\vee}$ and with the subgroup of $\operatorname{Aut}_{\mathbb{Z}}(X)$ generated by the reflections $s_{\alpha}$ for $\alpha \in \Phi$. Choose a basis $\Delta$, and consider the associated based root datum $\Psi_{0}=\left(X, \Phi, \Delta, X^{\vee}, \Phi^{\vee}, \Delta^{\vee}\right)$.

There is an obvious notion of isomorphism of root data (resp. based root data) $\Psi \rightarrow \Psi^{\prime}$ (resp. $\Psi_{0} \rightarrow \Psi_{0}^{\prime}$ ). It amounts to giving a Z-linear isomorphism $f: X \rightarrow X^{\prime}$ such that $f(\Phi)=\Phi^{\prime}$ and ${ }^{t} f\left(f(\alpha)^{\vee}\right)=\alpha^{\vee}$ for all $\alpha \in \Phi$ (resp. and $f(\Delta)=\Delta^{\prime}$ ). Here ${ }^{t} f$ denotes the transpose with respect to the root data pairings. We denote by $\operatorname{Aut}(\Psi)\left(\operatorname{resp} . \operatorname{Aut}\left(\Psi_{0}\right)\right)$ the group of automorphisms of $\Psi\left(\right.$ resp. $\left.\Psi_{0}\right)$. Each $s_{\alpha}$ can be seen as an automorphism of $\Psi$, and thus, there is a natural inclusion $W \subset \operatorname{Aut}(\Psi)$. We also denote by $-1 \in \operatorname{Aut}(\Psi)$ the automorphism that sends $x \in X$ to $-x \in X$.

Assume from now on that $\Phi$ is reduced. If $\Delta$ is a basis, let $w_{0}$ be the longest element of $W$ with respect to it. Then $w_{0}(\Delta)=-\Delta$, and thus, $-w_{0}=-1 \circ w_{0} \in$ $\operatorname{Aut}\left(\Psi_{0}\right)$. We call $\star=-w_{0}$ the opposition involution of $\Psi_{0}$ (since $w_{0}^{2}=1$ it is indeed an involution). We denote the action of $\star$ on elements $x$ (which can be characters of $T$, nodes of the Dynkin diagram, etc.) by $x \mapsto x^{\star}$. When $\Phi=\varnothing$, in which case $\Psi$ is called toral, we directly define $\star=-1 \in \operatorname{Aut}_{\mathbb{Z}}(X)$.

Remark 3.1.1. An isogeny (in particular, an isomorphism) of based root data will commute with the corresponding opposition involutions. In particular, $\star$ is a central element of Aut $\Psi_{0}$.

Remark 3.1.2. Let $X_{0} \subset X$ denote the subgroup of $X$ orthogonal to $\Phi^{\vee}$. The root datum $\Psi$ is called semisimple when $X_{0}=0$. If this is not the case, then there exists a nonzero $x \in X_{0}$, which hence must be invariant under $W$. In particular,
$x^{\star}=-x \neq x$, so $\star$ cannot be the identity map if the root datum is not semisimple. In the same vein, if the root datum is toral, then $\star \neq 1$ unless $\Psi$ is trivial (that is, also semisimple).

Suppose now that $k$ is an algebraically closed field of characteristic 0 , and let $G$ be a reductive group over $k$. Let $T \subset G$ be a maximal torus, and $\Psi=\Psi(G, T)$ be the associated root datum, so that $X=X^{*}(T)$. Let $B \supset T$ be a Borel subgroup, and let $\Psi_{0}=\Psi_{0}(G, T, B)$ be the corresponding based root datum. Let Aut $(G)$ be the group of automorphisms of $G$, and $\operatorname{Inn}(G) \subset \operatorname{Aut}(G)$ be the subgroup of inner automorphisms (that is, defined by elements in $G(k)$ ). Thus, $\operatorname{Inn}(G) \simeq G^{\text {ad }}(k) \simeq$ $G(k) / Z(k)$, where $Z$ is the center of $G$. Then there is a split exact sequence

$$
\begin{equation*}
1 \rightarrow \operatorname{Inn}(G) \rightarrow \operatorname{Aut}(G) \rightarrow \operatorname{Aut} \Psi_{0} \rightarrow 1 \tag{3.1.3}
\end{equation*}
$$

where, for $f \in \operatorname{Aut}(G)$, the third arrow sends $f$ to the automorphism of $\Psi_{0}$ induced by $f^{\prime} \in \operatorname{Aut}(G, T, B)$, where $f^{\prime}=\operatorname{int}(g) \circ f$ for any element $g \in G(k)$ such that $\operatorname{int}(g) f(B, T)=(B, T)$. We define an opposition involution of $G$ (with respect to $(B, T))$ to be any element $\theta \in \operatorname{Aut}(G)$ of order 1 or 2 that induces the opposition involution $\star$ in Aut $\Psi_{0}$. Note that this definition does not require $\theta$ to preserve $T$ or $B$. If $\theta^{\prime}$ is another such involution, then $\theta^{\prime}=\operatorname{int}(g) \circ \theta$ for some $g \in G(k)$. If $\theta$ is an opposition involution for $(B, T)$ and ( $B^{\prime}, T^{\prime}$ ) is another Borel pair, then it is also an opposition involution for $\left(B^{\prime}, T^{\prime}\right)$. The exact sequence (3.1.3) is split by the choice of a pinning. More precisely, let $\Delta \subset \Phi$ be the set of simple roots corresponding to $B$. For each $\alpha \in \Delta$, let $U_{\alpha} \in G$ be the root group of $\alpha$ [Springer 1979, §2.3], and let $u_{\alpha} \in U_{\alpha}$ be a nontrivial element. The pinning is the datum $\left\{u_{\alpha}\right\}_{\alpha \in \Delta}$ with respect to $(B, T)$, and a splitting $\operatorname{Aut} \Psi_{0} \rightarrow \operatorname{Aut}(G)$ of (3.1.3) associated with this pinning is given by an isomorphism $\operatorname{Aut} \Psi_{0} \simeq \operatorname{Aut}\left(G, T, B,\left\{u_{\alpha}\right\}_{\alpha \in \Delta}\right)$; two such splittings differ by an automorphism $\operatorname{int}(t)$ for some $t \in T(k)$. In particular, after choosing a pinning, we can take $\theta \in \operatorname{Aut}(G)$ to be the image of $\star$ under the splitting and this will be an opposition involution, which proves their existence. Note that we are actually showing that there are opposition involutions in $\operatorname{Aut}(G)$ which preserve $T$ and $B$ (and a fixed pinning).

Let $k$ be any field of characteristic 0 , and $\bar{k}$ be an algebraic closure of $k$. Let $\Gamma=\operatorname{Aut}(\bar{k} / k)$. Let $G$ be a reductive group over $k, T \subset G$ a maximal torus, and $B \supset T_{\bar{k}}$ a Borel subgroup of $G_{\bar{k}}$. Let $\Psi=\Psi\left(G_{\bar{k}}, T_{\bar{k}}\right)$ and $\Psi_{0}=\Psi_{0}\left(G_{\bar{k}}, T_{\bar{k}}, B\right)$. There is a natural action of $\Gamma$ on $X$, denoted by $\chi \mapsto^{\gamma} \chi$, where

$$
\gamma \chi(t)=\gamma\left(\chi\left(\gamma^{-1}(t)\right)\right)
$$

for $\gamma \in \Gamma$ and $t \in T(\bar{k})$. We call it the usual action of $\Gamma$ on $X$. It defines an action of $\Gamma$ on $\Psi$. Let $\gamma \in \Gamma$. Then we define a second action $\mu_{G}(\gamma)$ on $X$, the *-action, given by $\mu_{G}(\gamma)(\chi)(t)={ }^{\gamma} \chi\left(n^{-1} t n\right)$ for $t \in T(\bar{k})$, where $n \in G(\bar{k})$ is
an element such that $\operatorname{int}(n)$ sends the Borel pair $\left(\gamma(B), \gamma\left(T_{\bar{k}}\right)\right)$ to $\left(B, T_{\bar{k}}\right)$. For example, if $B$ is a Borel defined over $k$, then we can take $n=1$ and the $*$-action is just the usual action $\chi \mapsto^{\gamma} \chi$. Going back to the general case, this gives a morphism $\mu_{G}: \Gamma \rightarrow$ Aut $\Psi_{0}$, and it induces an action of $\Gamma$ on Aut $\Psi_{0}$ by taking $\rho \mapsto \mu_{G}(\gamma) \circ \rho \circ \mu_{G}(\gamma)^{-1}$ for $\rho \in$ Aut $\Psi_{0}$. There is also an action of $\Gamma$ on $\operatorname{Aut}\left(G_{\bar{k}}\right)$ given by $\gamma \cdot f=\left(\mathrm{id}_{G} \times \operatorname{Spec}(k) \operatorname{Spec}\left(\gamma^{-1}\right)\right) \circ f \circ\left(\mathrm{id}_{G} \times \operatorname{Spec}(k) \operatorname{Spec}(\gamma)\right)$, which on $G(\bar{k})$-points is simply $g \mapsto \gamma\left(f\left(\gamma^{-1}(g)\right)\right)$. It preserves the subgroup $\operatorname{Inn}\left(G_{\bar{k}}\right)=G(\bar{k}) / Z(\bar{k})$, where it acts as usual. The exact sequence (3.1.3) becomes

$$
\begin{equation*}
1 \rightarrow \operatorname{Inn}\left(G_{\bar{k}}\right) \rightarrow \operatorname{Aut}\left(G_{\bar{k}}\right) \rightarrow \operatorname{Aut} \Psi_{0} \rightarrow 1 \tag{3.1.4}
\end{equation*}
$$

and is $\Gamma$-equivariant. We define an opposition involution of $G$ to be an automorphism $\theta \in \operatorname{Aut}(G)$ of order 1 or 2 such that $\theta_{\bar{k}}$ is an opposition involution on $G_{\bar{k}}$.

There may not be a $\Gamma$-equivariant splitting of (3.1.4), so it may not always be possible to construct in this way an opposition involution of $G$. However, if $G$ is quasisplit and $B$ is a Borel subgroup defined over $k$, it can be shown [Demazure $1965 / 66, \S 3.10$ ] that there exists a $\Gamma$-equivariant splitting. Since $\star \in \operatorname{Aut} \Psi_{0}$ is central, it commutes with $\mu_{G}(\gamma)$ for any $\gamma \in \Gamma$, and thus, it is a $\Gamma$-invariant element in the last group of (3.1.4). Thus, for quasisplit reductive groups over $k$, there always exist opposition involutions on $G$ over $k$, but the condition of $G$ being quasisplit is far from necessary. There are many nonquasisplit cases where the opposition involution is trivial (see below), and so obviously defined over $k$. There are many nontrivial examples as well, as we will see later.

Remark 3.1.5. If $G=T$ is a torus, then there exists one and only one opposition involution $\theta \in \operatorname{Aut}(G)$, namely $\theta=\operatorname{inv}_{G}$.
Lemma 3.1.6. If $\theta$ is an opposition involution of $G$, then $\theta_{Z}: Z \rightarrow Z$ is equal to $\operatorname{inv}_{Z}$.
Proof. It is enough to see that both maps induce the same map on $X^{*}(Z)$, that is, that $\theta_{Z}^{*}: X^{*}(Z) \rightarrow X^{*}(Z)$ is multiplication by -1 , and thus, we can assume that $k=\bar{k}$. Let $(B, T)$ be a Borel pair. Then $Z \subset T$. Let $\chi \in X^{*}(Z)$. Then there exists $\mu \in X^{*}(T)$ such that $\left.\mu\right|_{Z}=\chi$. We claim that $\theta_{Z}^{*}(\chi)=\left.\left(\mu^{\star}\right)\right|_{Z}$. Indeed, for $z \in Z(k)$, $\theta_{Z}^{*}(\chi)(z)=\chi(\theta(z))$, whereas $\left.\left(\mu^{\star}\right)\right|_{Z}(z)=\Psi_{0}(\theta)(\mu)(z)=\mu((\operatorname{int}(g) \circ \theta)(z))=$ $\mu(\theta(z))$ (where $g \in G(k)$ sends $\theta(B, T)$ to $(B, T)$ ), which shows that $\theta_{Z}^{*}=\left.\left(\mu^{\star}\right)\right|_{Z}$.

On the other hand, if $n_{0} \in N_{G}(T)(k)$ represents $w_{0} \in W=N_{G}(T)(k) / T(k)$, then for $z \in Z(k), \mu^{\star}(z)=\mu\left(n_{0}^{-1} z^{-1} n_{0}\right)=\mu\left(z^{-1}\right)=\mu^{-1}(z)$ because $z \in Z(k)$. Thus, $\theta_{Z}^{*}(\chi)=-\chi$, as desired, where we have switched back to the additive notation for the group $X^{*}(Z)$.
Remark 3.1.7. The last lemma shows in particular that if the identity map is an opposition involution, then $Z$ is killed by 2 . Then $Z^{0}$ must be trivial; that is, $G$ must be semisimple (see also Remark 3.1.2).
3.2. Dynkin diagrams and special nodes. Let $\Psi_{0}$ be a based root datum with $\Phi \neq \varnothing$ and reduced, and let $\mathcal{D}$ be its Dynkin diagram. Then the opposition involution $\star$ acts on $\mathcal{D}$. We include for reference the list of connected Dynkin diagrams and their opposition involutions; see [Bourbaki 2002] for notation of nodes and more details. We also list the special nodes of each diagram (see [Deligne 1979, §1.2.5], for the definition of special node). Also, note that if $\Psi$ is semisimple, then $\star$ is trivial on $\Psi_{0}$ if and only if it is trivial on $\mathcal{D}$. For a Shimura datum $(G, X)$, the only factors of $G^{\text {ad }}$ that contribute to a CM reflex field are the ones of type $A_{l}(l \geq 2), D_{l}(l \geq 5 \mathrm{odd})$, or $E_{6}$. This follows from the list below and Proposition 2.3.6 of [Deligne 1979]:

- $\mathcal{D}=A_{l}(l \geq 1)$.
$\alpha_{i}^{\star}=\alpha_{l+1-i}($ so $\star$ is trivial if $l=1)$.
All nodes $\alpha_{i}$ are special.
- $\mathcal{D}=B_{l}(l \geq 2)$ or $C_{l}(l \geq 3)$.
$\star$ is trivial.
There is only one special node: $\alpha_{1}$ in the $B_{l}$ case, and $\alpha_{l}$ in the $C_{l}$ case.
- $\mathcal{D}=D_{l}(l \geq 4)$.

If $l$ is even, $\star$ is trivial.
If $l$ is odd, $\alpha_{i}^{\star}=\alpha_{i}$ for $i<l-1$, and $\alpha_{l-1}^{\star}=\alpha_{l}$.
The special nodes are $\alpha_{1}, \alpha_{l-1}$, and $\alpha_{l}$.

- $\mathcal{D}=E_{6}$.
$\alpha_{1}^{\star}=\alpha_{6}, \alpha_{2}^{\star}=\alpha_{2}, \alpha_{3}^{\star}=\alpha_{5}$, and $\alpha_{4}^{\star}=\alpha_{4}$.
The special nodes are $\alpha_{1}$ and $\alpha_{6}$.
- $\mathcal{D}=E_{7}, E_{8}, F_{4}$, or $G_{2}$.
$\star$ is trivial.
Only $E_{7}$ has a special node, which is $\alpha_{7}$.
3.3. Multiplicative groups of CM type. From now on let $k=\mathbb{Q}$ and $\Gamma=\operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q})$. Let $T_{1}$ and $T_{2}$ be algebraic groups over $\mathbb{Q}$ of multiplicative type, not necessarily connected. Then there is a natural bijection $\operatorname{Hom}\left(T_{1}, T_{2}\right) \simeq \operatorname{Hom}_{\Gamma}\left(X_{2}, X_{1}\right)$, where Aut ${ }_{\Gamma}$ means $\Gamma$-equivariant morphisms for the natural Galois structures on $X_{i}=X^{*}\left(T_{i}\right)$. In particular, for $T$ over $\mathbb{Q}$ of multiplicative type, there is a natural isomorphism $\operatorname{Aut}(T) \simeq \operatorname{Aut}_{\Gamma}(X)$, with $X=X^{*}(T)$. We let $c_{T}^{*}: X \rightarrow X$ be the map $c_{T}^{*}(\chi)={ }^{c} \chi$. We say $T$ splits over an extension $K \subset \overline{\mathbb{Q}}$ of $\mathbb{Q}$ if $\operatorname{Aut}(\overline{\mathbb{Q}} / K)$ acts trivially on $X^{*}(T)$.

Lemma 3.3.1. If $T$ is a group of multiplicative type that splits over a CM field, then $c_{T}^{*} \in \operatorname{Aut}_{\Gamma}(X)$.

Proof. Suppose that $T$ splits over $K \subset \overline{\mathbb{Q}}$, a CM field. Let $\chi \in X$. Then ${ }^{\gamma} \chi=\chi$ for any $\gamma \in \operatorname{Aut}(\overline{\mathbb{Q}} / K)$, and thus, ${ }^{\gamma_{1}} \chi={ }^{\gamma_{2}} \chi$ if $\gamma_{1}, \gamma_{2} \in \Gamma$ have the same restriction
to $K$. For any $\gamma \in \Gamma, \gamma c$ and $c \gamma$ have the same restriction to $K$, and so $c_{T}^{*}(\gamma \chi)=$ ${ }^{c}\left({ }^{\gamma} \chi\right)={ }^{c \gamma} \chi={ }^{\gamma c} \chi={ }^{\gamma}\left(c_{T}^{*}(\chi)\right)$.

Under the assumptions of the last lemma, we let $c_{T}: T \rightarrow T$ denote the unique involution inducing $c_{T}^{*}$ on $X$. If $T_{1}$ and $T_{2}$ are groups of multiplicative type which are split over a CM field, and $f: T_{1} \rightarrow T_{2}$ is a morphism, then $f \circ c_{T_{1}}=c_{T_{2}} \circ f$, because both maps induce the same morphism $X_{2} \rightarrow X_{1}$.

Suppose now that $T$ is a group of multiplicative type over $\mathbb{R}$. Using the same procedure, there exists a unique involution $c_{T}: T \rightarrow T$ inducing complex conjugation on characters. If $T$ is defined over $\mathbb{Q}$ and split over a CM field, these definitions are compatible with base change from $\mathbb{Q}$ to $\mathbb{R}$.

Example 3.3.2. For $T=\mathbb{S}$ over $\mathbb{R}$, the map $c_{\mathbb{S}}$ is given by $c_{\mathbb{S}}(z \otimes a)=\bar{z} \otimes a$ for an $\mathbb{R}$-algebra $A$ and $z \otimes a \in\left(\mathbb{C} \otimes_{\mathbb{R}} A\right)^{\times}$.

Remark 3.3.3. If $T$ is an anisotropic $\mathbb{R}$-torus (that is, if $T(\mathbb{R})$ is compact), then it is easy to see that ${ }^{c} \chi=-\chi$ for any $\chi \in X$ and thus $c_{T}=\operatorname{inv}_{T}$ is the opposition involution on $T$.
3.4. Involutions taking $\boldsymbol{X}$ to $\overline{\boldsymbol{X}}$. Let $(G, X)$ be a Shimura datum. Recall that we are assuming that $Z^{0}$ splits over a CM field, and hence, we have the conjugation involution $c_{Z^{0}}: Z^{0} \rightarrow Z^{0}$.

Remark 3.4.1. Let $x \in X$. From the fact that $\operatorname{int}(x(i)): G_{\mathbb{R}}^{\text {ad }} \rightarrow G_{\mathbb{R}}^{\text {ad }}$ is a Cartan involution, it follows that $G_{\mathbb{R}}^{\text {ad }}$ is an inner form of an anisotropic group $H$ over $\mathbb{R}$ (that is, $H(\mathbb{R})$ is compact). A similar statement holds for $G_{\mathbb{R}}^{\mathrm{der}}$ (the element $x(i)$ may not belong to $G^{\operatorname{der}}(\mathbb{R})$; however, over $\mathbb{C}, \operatorname{int}(x(i))$ can be replaced by $\operatorname{int}\left(x(i)^{\prime}\right)$ for some $\left.x(i)^{\prime} \in\left(T \cap G^{\text {der }}\right)(\mathbb{C})\right)$. The next lemma is well known.

Lemma 3.4.2. Let $G$ be a reductive group over $\mathbb{R}$, and assume that it is an inner form of a group $H$ over $\mathbb{R}$ which is anisotropic. Assume furthermore that $T \subset G$ is a maximal torus, and the inner automorphism of $G_{\mathbb{C}}$ defining a cocycle for $H$ is given by $\operatorname{int}\left(t_{0}\right)$ for some $t_{0} \in T(\mathbb{C})$. Then the following hold.
(i) $c_{T}=\operatorname{inv}_{T}$.
(ii) For a Borel subgroup $B \supset T_{\mathbb{C}}$, the opposition involution acting on $\Psi_{0}(G, T, B)$ is given by the $*$-action of $c$.
(iii) The subgroup $c(B) \subset G_{\mathbb{C}}$ is the opposite Borel subgroup of $B$; that is,

$$
c(B) \cap B=T_{\mathbb{C}} .
$$

Proof. By hypothesis, we can choose an isomorphism $\phi: G_{\mathbb{C}} \rightarrow H_{\mathbb{C}}$ such that $f: G_{\mathbb{C}} \rightarrow G_{\mathbb{C}}$ defined by $f(g)=\phi^{-1} \overline{(\phi(\bar{g}))}$ is an inner automorphism of the form $\operatorname{int}\left(t_{0}\right)$, with $t_{0} \in T(\mathbb{C})$. Then there exists a maximal torus $T_{H} \subset H$ such
that $T_{H, \mathbb{C}}=\phi\left(T_{\mathbb{C}}\right)$, and we let $B_{H}=\phi(B)$. Since $T_{H}(\mathbb{R})$ is compact, this implies that $T(\mathbb{R})$ is compact, so (i) follows from Remark 3.3.3. Let $r=\Psi_{0}(\phi)$ : $\Psi_{0}\left(H, T_{H}, B_{H}\right) \rightarrow \Psi_{0}(G, T, B)$ be the induced isomorphism. It is $\operatorname{Gal}(\mathbb{C} / \mathbb{R})$ equivariant for the $*$-actions, as follows from the fact that the forms are inner, and it commutes with $\star$, so it is enough to prove part (ii) when $G$ itself is anisotropic, which is well known. For (iii), the fact that $f$ preserves $T_{\mathbb{C}}$ and $B$ again allows us to reduce to the case of $G$ anisotropic, in which case the statement is well known.
Remark 3.4.3. In the last lemma, if the group is quasisplit and $B$ is a Borel subgroup defined over $\mathbb{R}$, the inner automorphism will not usually belong to $T(\mathbb{C})$; otherwise, we would have $B=T$. There are quasisplit semisimple groups with $B \neq T$ which are inner forms of anisotropic groups, for example $\mathrm{SU}(n, n)$. In this case, the Cartan involution coming from a certain Shimura datum and special pair will preserve the maximal torus and a Borel subgroup containing it, but not a rational Borel subgroup.

Remark 3.4.4. Suppose that $(G, X)$ is a Shimura datum, and let $(T, x)$ be a special pair. Then $G_{\mathbb{R}}^{\text {der }}$ satisfies all the hypotheses of the previous lemma. Here the inner automorphism defining the cocycle is $\operatorname{int}\left(x(i)^{\prime}\right)$ as before. Alternatively, we can work with the adjoint group $G_{\mathbb{R}}^{\text {ad }}$ and $x(i)$.
Remark 3.4.5. Suppose that $\theta: G \rightarrow G$ is an involution such that there exists a special pair $(T, x)$ with the property that $\theta$ preserves $T$ and induces $c_{T_{\mathbb{R}}}$ on $T_{\mathbb{R}}$. Then $\theta_{\mathbb{R}}(x)=c_{T_{\mathbb{R}}} \circ x=x \circ c_{\mathbb{S}}=\bar{x}$, and thus, $\theta(X)=\bar{X}$.
Lemma 3.4.6. Let $G$ be a reductive group over $\mathbb{R}$, and $T \subset G$ a maximal torus. If $\theta: G \rightarrow G$ is an involution such that $\theta(T) \subset T$ and $\left.\theta\right|_{T}=c_{T}$, then $\theta(B)=c(B) \subset G_{\mathbb{C}}$ for any Borel subgroup $B \supset T_{\mathbb{C}}$.
Proof. Let $R \subset X=X^{*}(T)$ denote the set of roots of ( $G_{\mathbb{C}}, T_{\mathbb{C}}$ ). Let $R^{+}$denote the set of positive roots with respect to $B$. Then $\theta(B)$ is the Borel subgroup whose Lie algebra is $\operatorname{Lie}\left(T_{\mathbb{C}}\right) \oplus \bigoplus_{\alpha \in R^{+}} \operatorname{Lie}\left(G_{\mathbb{C}}\right)_{\alpha \circ \theta}$. Since $\alpha \circ \theta={ }^{c} \alpha$, it follows that this is the Lie algebra of $c(B)$, and since both $\theta(B)$ and $c(B)$ are connected, this proves the lemma.

The construction of involutions taking $X$ to $\bar{X}$ that we will perform will be based on involutions $\theta$ which will roughly be as in Remark 3.4.5. By the following proposition, we need to look for opposition involutions on semisimple groups.
Proposition 3.4.7. Let $(G, X)$ be a Shimura datum, and let $\theta: G \rightarrow G$ be an involution of $G$, such that there exists a special pair $(T, x)$ with the property that $\theta$ preserves $T$ and induces $c_{T_{\mathbb{R}}}$ on $T_{\mathbb{R}}$. Then $\theta^{\mathrm{der}}: G^{\mathrm{der}} \rightarrow G^{\mathrm{der}}$ is an opposition involution, and $\theta_{0}=\left.\theta\right|_{Z^{0}}: Z^{0} \rightarrow Z^{0}$ is equal to $c_{Z^{0}}$.
Proof. Suppose that $\theta$ is an involution with $(T, x)$ as in the statement. To see that $\theta_{0}=c_{Z^{0}}$, it is enough to see that $\theta_{\mathbb{R}, 0}=c_{Z_{\mathbb{R}}^{0}}$. Since $Z_{\mathbb{R}}^{0} \subset T_{\mathbb{R}}$ and $\left.\theta_{\mathbb{R}}\right|_{T_{\mathbb{R}}}=c_{T_{\mathbb{R}}}$,
it follows that $\theta_{0, \mathbb{R}}=c_{Z_{\mathbb{R}}^{0}}$. Let $T^{\prime}=T \cap G^{\text {der }}$, let $B \subset G_{\mathbb{C}}$ be a Borel subgroup containing $T_{\mathbb{C}}$, and $B^{\prime}=B \cap G_{\mathbb{C}}^{\text {der }} \supset T_{\mathbb{C}}^{\prime}$. Let $\Psi_{0}^{\prime}=\Psi_{0}\left(G^{\text {der }}, T^{\prime}, B^{\prime}\right)$, and let $r=\Psi_{0}\left(\theta^{\mathrm{der}}\right): \Psi_{0}^{\prime} \rightarrow \Psi_{0}^{\prime}$ be the induced isomorphism. It is given by $r(\chi)=\chi \circ \operatorname{int}(q) \circ$ $\left.\theta^{\text {der }}\right|_{T_{\mathbb{C}}^{\prime}}$ for $\chi \in X^{\prime}=X^{*}\left(T^{\prime}\right)$, where $q \in G^{\operatorname{der}}(\mathbb{C})$ is such that $\operatorname{int}(q) \theta^{\operatorname{der}}\left(T_{\mathbb{C}}^{\prime}, B^{\prime}\right)=$ $\left(T_{\mathbb{C}}^{\prime}, B^{\prime}\right)$. On the other hand, by Lemma 3.4.2(ii), $\star: \Psi_{0}^{\prime} \rightarrow \Psi_{0}^{\prime}$ is given by $\chi^{\star}=$ ${ }^{c} \chi \circ \operatorname{int}\left(a^{-1}\right)$, where $a \in G^{\operatorname{der}}(\mathbb{C})$ is such that $\operatorname{int}(a) c\left(T_{\mathbb{C}}^{\prime}, B^{\prime}\right)=\left(T_{\mathbb{C}}^{\prime}, B^{\prime}\right)$. By Lemma 3.4.6, we can take $a=q$. Finally, the hypothesis that $\left.\theta^{\text {der }}\right|_{T^{\prime}}=c_{T^{\prime}}$ implies that $\chi^{\star}=\chi \circ \theta \circ \operatorname{int}\left(q^{-1}\right)$. Thus, to see that $r(\chi)=\chi^{\star}$, it is enough to see that $\theta^{\operatorname{der}} \circ \operatorname{int}\left(q^{-1}\right)$ and $\theta^{\operatorname{der}} \circ \operatorname{int}(\varphi(q))$ induce the same automorphism of $T_{\mathbb{C}}^{\prime}$, and this follows from the fact that both elements $\theta^{-1}(q)$ and $q^{-1}$ conjugate the Borel pair ( $T_{\mathbb{C}}^{\prime}, B^{\prime}$ ) to the same Borel pair.

The following proposition is a partial converse and the main result of this section. Since our construction will be explicit using the classification of semisimple groups, we need to work with either the derived group or the adjoint group. The idea is to construct an involution on $G$ taking $X$ to $\bar{X}$ by extending an opposition involution on $G^{\text {der }}$. Ideally we would want the involution to be as in Remark 3.4.5, but it is enough to consider a weaker hypothesis, as stated in the proposition. Recall the notation from Section 2.1. Suppose that for each $i, S_{i} \subset H_{i}$ is a maximal torus, and let $\widetilde{T}_{i}=\operatorname{Res}_{F_{i} / \mathbb{Q}} S_{i} \subset \widetilde{G}_{i}, T_{i} \subset G_{i}$ its image in $G_{i}, T^{\prime} \subset G^{\text {der }}$ the image of their product, and $T=Z^{0} T^{\prime}$. Note that $T_{\mathbb{R}}^{\text {ad }}=T_{\mathbb{R}}^{\prime \text { ad }}=\prod_{i, v} S_{i, v}^{\text {ad }}$, where $S_{i, v} \subset H_{i, v}$ and $S_{i, v}^{\text {ad }}$ is its image in $H_{i, v}^{\text {ad }}$.
Proposition 3.4.8. Suppose that $\theta_{i}: H_{i} \rightarrow H_{i}$ is an opposition involution for each $i$. Suppose moreover that $\theta_{i}\left(S_{i}\right)=S_{i}$ and $\left.\theta_{i, v}^{\text {ad }}\right|_{S_{i, v}} ^{\text {ad }}=c_{S_{i, v}^{\text {ad }}}$ for every $i$ and $v \in I_{i, n c}$. Finally, assume that there exists $x \in X$ such that $x^{\text {ad }}$ factors through $T_{\mathbb{R}}$. ${ }^{\text {dd }}$. Then there exists an involution $\theta: G \rightarrow G$ such that $\theta(X)=\bar{X}$.

Proof. For each $i$, the involution $\operatorname{Res}_{F_{i} / \mathbb{Q}} \theta_{i}$ defines an opposition involution of $\widetilde{G}_{i}$. Moreover, the kernel $K_{i}$ of the projection $\widetilde{G}_{i} \rightarrow G_{i}$ is contained in the center of $\widetilde{G}_{i}$. By Lemma 3.1.6, $\operatorname{Res}_{F_{i} / \mathbb{Q}} \theta_{i}$ induces $x \mapsto x^{-1}$ on the center. In particular, it preserves $K_{i}$ and induces an opposition involution on $G_{i}$. Similarly, the product of these involutions defines an opposition involution $\theta^{\prime}: G^{\text {der }} \rightarrow G^{\text {der }}$. Let $q$ : $Z^{0} \times G^{\text {der }} \rightarrow G$ be the natural isogeny. We can look at the product involution $\theta^{\prime} \times c_{Z^{0}}: G^{\text {der }} \times Z^{0} \rightarrow G^{\text {der }} \times Z^{0}$. We claim that this preserves $\operatorname{ker}(q)$, and thus, it induces an involution on $G$. To show this, we can work with $\mathbb{C}$-points. The kernel consists of pairs $(g, z)$ such that $z g=1$, so we need to check that if $(g, z)$ is such a pair, then $\theta^{\prime}(g) c_{Z^{0}}(z)=1$. The element $g=z^{-1}$ belongs to $Z^{0} \cap G^{\text {der }} \subset Z_{G^{\text {der }}}$. The maps $c_{Z^{0}}: Z^{0} \rightarrow Z^{0}$ and $c_{Z_{G^{\text {der }}}}: Z_{G^{\text {der }}} \rightarrow Z_{G^{\text {der }}}$ are equal on $Z^{0} \cap G^{\text {der }}$, and by part (i) of Lemma 3.4.2, $c_{Z_{G^{\text {der }}}}=\operatorname{inv}_{Z_{G^{\text {der }}}}$, so $c_{Z^{0}}(z)=z^{-1}$. On the other hand, by Lemma 3.1.6, $\theta^{\prime}$ induces $\operatorname{inv}_{Z_{G^{\text {der }}}}$ on $Z_{G^{\text {der }}}$, and so $\theta^{\prime}(g)=g^{-1}=z$. This proves that there exists a (unique) involution $\theta: G \rightarrow G$ such that $\theta^{\mathrm{der}}=\theta^{\prime}$ and $\theta_{0}=c_{Z^{0}}$.

We also have that $\theta$ preserves $T$ and $\theta_{\mathbb{R}}^{\text {ad }}=\prod_{i, v} \theta_{i, v}^{\text {ad }}$. Now, we know that there exists $x \in X$ such that $x^{\text {ad }}$ factors through $T_{\mathbb{R}}^{\text {ad }}$. Let $y=\theta_{\mathbb{R}}(x)$. Then $y^{\text {ad }}=\theta_{\mathbb{R}}^{\text {ad }}\left(x^{\text {ad }}\right)$. For $v \in I_{i, n c}$, we have $\theta_{i, v}^{\text {ad }}\left(x_{i, v}\right)=\overline{x_{i, v}}$ because $\theta_{i, v}^{\text {ad }}$ induces $c_{S_{i, v}^{\text {ad }}}$ and $x_{i, v}$ factors through $S_{i, v}^{\text {ad }}$. For $v \in I_{i, c}, x_{i, v}=1$. Thus, $y^{\text {ad }}=\bar{x}^{\text {ad }}$. Also, since $\theta_{\mathbb{R}}$ induces $c_{Z_{\mathbb{R}}^{0}}$ on $Z_{\mathbb{R}}^{0}$, and $q: Z^{0} \rightarrow G^{\mathrm{ab}}$ is an isogeny, it follows that $\theta_{\mathbb{R}}$ induces $c_{G_{\mathbb{R}}}^{\mathrm{ab}}$ on $G_{\mathbb{R}}^{\mathrm{ab}}$. From this it follows that $y$ and $\bar{x}$ have the same projections to $G_{\mathbb{R}}^{\text {ad }}$ and to $G_{\mathbb{R}}^{\text {ab }}$, and thus, $y=\bar{x}$ (see for instance the proof of Proposition 5.7 of [Milne 2005]). Since $y=\theta_{\mathbb{R}}(x)$, this shows that $\theta(X)=\bar{X}$.

## 4. Involutions on classical semisimple groups

In this section, we make use of several results regarding the classification of semisimple algebraic groups over totally real fields. For notation and terminology regarding algebras with involutions and their associated groups, we freely follow our main reference [Knus et al. 1998]. We are only interested in the explicit classification of groups of type $A$ and $D$ in order to construct our desired involutions on certain Shimura varieties. Furthermore, not all the groups in the general classification appear in the theory of Shimura varieties, so we are only interested in classifying the groups $H_{i}$ (in the notation of Section 2.1) of type $A_{l}(l \geq 2)$ or $D_{l}(l \geq 4$ odd) that can occur. Furthermore, in accordance with the previous section, we are also interested in constructing, whenever possible, opposition involutions on these groups.

The following construction regarding quaternion algebras will be used often in the following. Suppose that $D$ is a quaternion division algebra over a number field $K$. Let $\lambda \in D^{\times}$be a pure quaternion (that is, such that $\sigma(\lambda)=-\lambda$, where $\sigma: D \rightarrow D$ is the canonical involution), and choose another pure quaternion $\mu \in D^{\times}$such that $\lambda \mu=-\mu \lambda$. Then $\{1, \lambda, \mu, \lambda \mu\}$ is a standard basis of $D$. If we let $L=K(\lambda)$, then $L$ is a maximal subfield of $D$ (a quadratic extension of $K$ ). We have an isomorphism of $L$-algebras $\phi: D \otimes_{K} L \rightarrow M_{2}(L)$ defined by

$$
\phi(\lambda \otimes 1)=\left(\begin{array}{rr}
\lambda & 0 \\
0 & -\lambda
\end{array}\right)
$$

and

$$
\phi(\mu \otimes 1)=\left(\begin{array}{cc}
0 & \mu^{2} \\
1 & 0
\end{array}\right)
$$

Then the isomorphism $\phi$ sends $L \otimes_{K} L$ to the subalgebra of diagonal matrices in $M_{2}(L)$.

Throughout this section, let $F$ be a totally real field and $H$ be an absolutely almost simple, simply connected algebraic group over $F$. We let $\mathcal{D}$ be the Dynkin diagram of $H_{\bar{F}}$ (where $\bar{F}$ is some algebraic closure of $F$ ). We let $I=\operatorname{Hom}(F, \mathbb{R})$, $I_{c}=\left\{v \in I: H_{v}^{\text {ad }}(\mathbb{R})\right.$ is compact $\}$, and let $I_{n c}$ be its complement in $I$.
4.1. Groups of type $\boldsymbol{A}_{\boldsymbol{l}}(\boldsymbol{l} \geq \mathbf{2})$. Suppose that $\mathcal{D}=A_{l}$ with $l \geq 2$. Then there exists a quadratic étale extension $K / F$ (so $K / F$ is a quadratic extension of fields, or $K=F \times F$ ), and a central simple algebra $B$ over $K$, of degree $l+1$, endowed with an involution $\tau: B \rightarrow B$ of the second kind (that is, inducing $\iota$ on $K$, where $\iota$ is the nontrivial automorphism of $K$ which fixes $F$ ) such that $H=\mathbf{S U}(B, \tau)$ [Knus et al. 1998, Theorem 26.9]. If $H$ is one of the $H_{i}$ as above, then $K$ must be a field. Indeed, if otherwise, then $H \simeq \mathbf{S L}_{1}(A)$ for some central simple algebra $A$ over $F$ of degree $l+1$. For each $v \in \operatorname{Hom}(F, \mathbb{R})$, we have

$$
A_{v}=A \otimes_{F, v} \mathbb{R} \simeq M_{l+1}(\mathbb{R}) \quad \text { or } \quad A_{v} \simeq M_{(l+1) / 2}(\mathbb{H})
$$

In both cases, it follows that $H_{v}$ is an inner form of $\mathrm{SL}_{l+1, \mathbb{R}}$, so the $*$-action of $c$ is trivial (a condition that does not depend on the Borel pair), and thus, it cannot be the opposition involution because $l \geq 2$. From this and Lemma 3.4.2 it follows that $H$ cannot occur as one of the factors $H_{i}$. Thus, we have proved that $K$ must be a field. Moreover, a similar argument implies that $K$ must be totally imaginary, that is, $K / F$ is a CM extension. The adjoint group $H^{\text {ad }}$ is $\mathbf{P G U}(B, \tau)$.

We can then write $B=\operatorname{End}_{D}(V)$ for some central division algebra $D$ over $K$, endowed with an involution $J: D \rightarrow D$ of the second kind, whose action we denote by $d \mapsto d^{J}$, and a finite-dimensional right $D$-vector space $V$. There is a nondegenerate hermitian form $h: V \times V \rightarrow D$ inducing the involution $\tau: B \rightarrow B$. The pair $(V, h)$ is called a hermitian space over $D$.

Suppose that $\theta: H \rightarrow H$ is an opposition involution. There is a natural isomorphism between $\operatorname{Aut}(H)$ and the group of $F$-algebra automorphisms of $B$ that commute with $\tau$ [Knus et al. 1998, Theorem 26.9], and thus, there exists such an automorphism $\gamma: B \rightarrow B$ of order 2, inducing $\theta$. If $\left.\gamma\right|_{K}$ is the identity map on $K$, then $\gamma=\operatorname{int}\left(b_{0}\right)$ for some $b_{0} \in B^{\times}$by the Skolem-Noether theorem, and $b_{0}$ is moreover a similitude for $\tau$. The induced map $\theta: H \rightarrow H$ would thus be an inner automorphism, inducing the identity map on the Dynkin diagram, but the opposition involution on $A_{l}$ is nontrivial for $l \geq 2$. Hence, $\left.\gamma\right|_{K}$ must be $\iota$. Let $\bar{B}$ and $\bar{D}$ denote the $K$-algebras $B$ and $D$ with $\iota$-conjugate structure. Thus, $\gamma: B \rightarrow \bar{B}$ is a $K$-algebra isomorphism. We let $\operatorname{Br}(K)$ be the Brauer group of $K$ and $[B]=[D] \in \operatorname{Br}(K)$ be the class of $B$ in it. Then $[D]=[B]=[\bar{B}]=[\bar{D}]$, which implies that there must exist a ring automorphism $\alpha: D \rightarrow D$ inducing $\iota$ on $K$.
Proposition 4.1.1. Let $D$ be a central division algebra over a CM extension $K / F$ of number fields, endowed with an involution $J: D \rightarrow D$ of the second kind. Then the following are equivalent:
(a) $D=K$ or $D$ is a quaternion division algebra over $K$.
(b) The order of $[D] \in \operatorname{Br}(K)$ is 1 or 2 .
(c) There exists a ring automorphism $\alpha: D \rightarrow D$ inducing $\iota$ on $K$.

Moreover, in this case, $\alpha$ is unique up to composition with an inner automorphism of $D$. Furthermore, it can be chosen to have order 2 and such that $\alpha J=J \alpha$ is either $\operatorname{id}_{D}$ if $D=K$ or the canonical involution if $D$ is a quaternion division algebra.

Proof. The fact that (a) implies (b) in the quaternion algebra case follows from the existence of the canonical involution on $D$, which gives an isomorphism $D \rightarrow D^{\mathrm{op}}$, so $[D]=[D]^{-1}$. To see that (b) implies (a), use [Scharlau 1985, 10.2.3].

Now suppose that (a) is true. If $D=K$, then take $\alpha=\iota$. If $D$ is a quaternion division algebra, let $\sigma: D \rightarrow D$ be its canonical involution, and take $\alpha=J \sigma=\sigma J$ (they commute because $J \sigma J$ is a symplectic involution of the first kind on $D$, and hence equal to $\sigma$ ).

Finally, suppose that $\alpha: D \rightarrow D$ is as in (c). The involution $J: D \rightarrow D$ induces an isomorphism $D \rightarrow \bar{D}^{\text {op }}$, where $\bar{D}$ is the conjugate algebra $\lambda \cdot d=\iota(\lambda) d$ for $\lambda \in K$. Similarly, $\alpha$ induces an isomorphism $D \rightarrow \bar{D}$, and thus, in the end we have an isomorphism $D \rightarrow D^{\mathrm{op}}$, which implies that the order of $[D]$ is 1 or 2 .

The uniqueness of $\alpha$ up to inner automorphism follows because if $\beta$ is another such automorphism, then $\alpha \beta^{-1}: D \rightarrow D$ is a $K$-linear automorphism and hence inner by the Skolem-Noether theorem.
Remark 4.1.2. Suppose that $D$ is a quaternion division algebra. Under the conditions of the previous proposition, there exists a unique quaternion algebra $D_{0} \subset D$ over $F$ such that $D=D_{0} \otimes_{F} K$ and $J=\sigma_{0} \otimes_{F} \iota$, where $\sigma_{0}$ is the canonical involution of $D_{0}$ [Knus et al. 1998, §II.22]. Then the map $\alpha$ constructed in the proof is $\alpha=\mathrm{id}_{D_{0}} \otimes_{F} \iota$. We define the canonical conjugation $\alpha: D \rightarrow D$ (attached to $J$ or $D_{0}$ ) to be $\alpha=\mathrm{id}_{D_{0}} \otimes_{F} \iota$. If $D=K$, we also call $\alpha=\iota$ the canonical conjugation.

Thus, we have shown that if there exists $\theta: H \rightarrow H$ an opposition involution, then either $D=K$ (and $J=\imath$ ) or $D$ is a quaternion division algebra (and $J=\sigma_{0} \otimes_{F} \iota$ ). Conversely, suppose that $D=K$ or $D$ is a quaternion division algebra. We will construct suitable opposition involutions under an additional assumption.

Remark 4.1.3. Suppose that $D=K$ (and $D_{0}=F$ ) or $D$ is a quaternion division algebra. Let $I_{s} \subset \operatorname{Hom}(F, \mathbb{R})$ be the subset of places $v \in I=\operatorname{Hom}(F, \mathbb{R})$ such that $D_{0, v}=D_{0} \otimes_{F, v} \mathbb{R}$ is split, and let $I_{n s} \subset I$ be its complement. We let $I_{c} \subset I$ be the subset of places $v$ such that $H_{v}^{\text {ad }}(\mathbb{R})$ is compact, and $I_{n c}$ its complement. The group $H_{\mathbb{R}}$ can be written as a product of special unitary groups $\prod_{v \in I} \mathrm{SU}\left(p_{v}, q_{v}\right)$, and the compact places are exactly the places where $p_{v} q_{v}=0$.

Definition 4.1.4. We say that a hermitian space $(V, h)$ over $D$ (where $D=K$ or a quaternion division algebra) is strongly hermitian if there exists an $h$-orthogonal $D$-basis $\beta=\left\{v_{1}, \ldots, v_{n}\right\}$ of $V$ such that $h\left(v_{i}, v_{i}\right) \in K^{\times}$for all $i$. In the quaternion algebra case, we ask furthermore that $I_{n s} \subset I_{c}$.
Remark 4.1.5. A hermitian space over $D=K$ is always strongly hermitian.

Remark 4.1.6. The existence of the basis $\beta$ in Definition 4.1.4 is what allows us to explicitly construct an opposition involution $\theta: H \rightarrow H$. In the quaternion algebra case, this involution will define an involution $\theta_{v}: H_{v} \rightarrow H_{v}$, and we want this to induce complex conjugation on $S_{v}^{\text {ad }}$ when $v \in I_{n c}$. The involution that we construct does not satisfy this at places $v \in I_{n s}$ (see Remark 4.1.9). Since we only care about noncompact places, we make the assumption $I_{n s} \subset I_{c}$.

Suppose that ( $V, h$ ) is strongly hermitian, and let $\beta$ be a basis as in the definition. Let $I: V \rightarrow V$ be the $\alpha$-semilinear isomorphism obtained by applying $\alpha$ to the coordinates of elements of $V$ with respect to the basis $\beta$ (this map is inspired by the constructions of [Taylor 2012]). Then $h(I(x), I(y))=\alpha(h(x, y))$. Let $\theta: H \rightarrow H$ be given as $\theta_{A}(g)=I_{A} g I_{A}$ for an $F$-algebra $A$ and a $D \otimes_{F} A$-linear automorphism $g$ of $V \otimes_{F} A$. Let $L \subset D$ be a maximal subfield. More precisely, if $D=K$, then $L=K$, and if $D$ is a quaternion division algebra, take $L=K(\lambda)$, where $\lambda$ is a pure quaternion in $D_{0}$. Let $S=S_{L, \beta}$ be the subgroup of $H$ defined as follows. For an $F$-algebra $A, H(A) \subset \operatorname{Aut}_{D \otimes_{F} A}\left(V \otimes_{F} A\right)$, and we let
$S(A)=\left\{h \in H(A): h\left(v_{i} \otimes 1\right)=\left(v_{i} \otimes 1\right) \lambda_{i}\right.$ for some $\left.\lambda_{i} \in\left(L \otimes_{F} A\right)^{\times}(i=1, \ldots, n)\right\}$.
This is a maximal torus in $H$.
Proposition 4.1.7. With the above hypotheses, the following statements are true:
(a) The involution $\theta: H \rightarrow H$ is an opposition involution.
(b) We have $\theta(S)=S$ and for every $v \in I_{n c}, \theta_{v}: H_{v} \rightarrow H_{v}$ induces $c_{S_{v}}$ on $S_{v}$.

In particular, $\theta_{v}^{\text {ad }}: H_{v}^{\text {ad }} \rightarrow H_{v}^{\text {ad }}$ induces $c_{S_{v}^{\text {ad }}}$ on $S_{v}^{\text {ad }}$ for $v \in I_{n c}$.
Proof. For part (a), it suffices to see that $\theta_{L}: H_{L} \rightarrow H_{L}$ is an opposition involution. We can identify $H_{K}$ with $\mathrm{SL}_{V / D}$, where $\mathrm{SL}_{V / D}(A)$ consists, for a $K$-algebra $A$, of the $D \otimes_{K} A$-linear automorphisms of $V \otimes_{K} A$ with reduced norm 1 . Using the basis $\beta$, we can further identify $\mathrm{SL}_{V / D}(A) \cong \mathrm{SL}_{n}\left(D \otimes_{K} A\right)$. Let $Q \in \mathrm{GL}_{n}(K)$ be the matrix of $h$ with respect to $\beta$. Then it is easy to see that $\theta_{A}: \operatorname{SL}_{n}\left(D \otimes_{K} A\right) \rightarrow \mathrm{SL}_{n}\left(D \otimes_{K} A\right)$ is explicitly given by the formula

$$
\theta_{A}(X)=Q^{-1}\left({ }^{t} X^{-1}\right)^{\sigma} Q
$$

where $\sigma: D \rightarrow D$ is the canonical involution of $D$ if $D$ is a quaternion division algebra, and $\sigma=\mathrm{id}$ if $D=K$. Note that $Q$ is a diagonal matrix in $\mathrm{GL}_{n}(K)$.

If $D=K$, we denote by $\phi: D \otimes_{K} L \rightarrow L$ the unique obvious isomorphism. If $D$ is a quaternion division algebra, we take $\phi: D \otimes_{K} L \rightarrow M_{2}(L)$ to be an isomorphism of $L$-algebras taking $L \otimes_{K} L$ to the subalgebra of diagonal matrices in $M_{2}(L)$, as constructed above (we use for this the pure quaternion $\lambda \in D_{0}$ and another pure quaternion $\mu \in D_{0}$ such that $\lambda \mu=-\mu \lambda$ ). In particular, $\sigma$ preserves $L$. The identification $H_{K}(A) \cong \operatorname{SL}_{n}\left(D \otimes_{K} A\right)$ sends $S_{K}(A)$ to the subgroup of matrices
in $\mathrm{SL}_{n}\left(D \otimes_{K} A\right)$ which are diagonal and have entries in $L \otimes_{K} A$. Since $\sigma$ preserves $L$, it follows that $\theta_{K}$ sends the torus $S_{K}$ to itself. Moreover, if we now extend scalars to $L$, the map $\phi$ provides an isomorphism

$$
\begin{equation*}
H_{L} \cong \mathrm{SL}_{n m, L} \tag{4.1.8}
\end{equation*}
$$

where $S_{n m, L}$ is the usual group of $n m \times n m$ matrices of determinant 1 ; furthermore, the torus $S_{L}$ maps to the torus of diagonal matrices in $\mathrm{SL}_{n m, L}$ (so $S$ is indeed a maximal torus, as claimed). If $m=1$, then $\theta_{L}(X)=Q^{-1 t} X^{-1} Q$ for $X \in \mathrm{SL}_{n, L}$. Suppose that $m=2$. Write $Q=\operatorname{diag}\left(q_{1}, \ldots, q_{n}\right)$, and let $\widetilde{Q}=\operatorname{diag}\left(q_{1}, \ldots, q_{n}, q_{1}, \ldots, q_{n}\right) \in$ $\mathrm{GL}_{2 n}(K)$. Write matrices $X \in \mathrm{SL}_{2 n, L}$ as blocks

$$
X=\left(\begin{array}{ll}
A & B \\
C & D
\end{array}\right),
$$

with $A, B, C, D$ of size $n \times n$. Then $\theta_{L}: \mathrm{SL}_{2 n, L} \rightarrow \mathrm{SL}_{2 n, L}$ is explicitly given as

$$
\theta_{L}(X)=\widetilde{Q}^{-1}\left(\begin{array}{rr}
{ }^{t} D & -{ }^{t} B \\
-{ }^{t} C & { }^{t} A
\end{array}\right) \widetilde{Q}
$$

From this explicit expression of $\theta$ as an involution of $\mathrm{SL}_{n m, L}$, it is easy to see that it preserves the maximal torus $S_{L}$ of diagonal matrices and that it induces the opposition involution on the root datum.

For part (b), fix $v \in I_{n c}$. We need to check that if $\chi \in X=X^{*}\left(S_{v}\right)=\operatorname{Hom}\left(S_{v} \times_{\mathbb{R}} \mathbb{C}\right.$, $\mathbb{G}_{\mathrm{m}, \mathbb{C}}$ ), then $\chi \circ \theta_{v, \mathbb{C}}={ }^{c} \chi$. To compute ${ }^{c} \chi$, we need to compute how complex conjugation acts on $H_{v}(\mathbb{C})$. Choose once and for all an extension $\tau: L \hookrightarrow \mathbb{C}$ of $v$ to $L$. Using the embedding $\tau$ and the isomorphism (4.1.8), we can identify $H_{v} \times_{\mathbb{R}} \mathbb{C}=H_{L} \times_{L, \tau} \mathbb{C} \cong \mathrm{SL}_{n m, \mathbb{C}}$. Moreover, the action of $c$ on $H_{v}(\mathbb{C}) \cong \mathrm{SL}_{n m}(\mathbb{C})$ is explicitly given as follows. Let $Q_{v}=\operatorname{diag}\left(v\left(q_{1}\right), \ldots, v\left(q_{n}\right)\right) \in \mathrm{GL}_{n}(\mathbb{R})$ and $\widetilde{Q}_{v}=\operatorname{diag}\left(v\left(q_{1}\right), \ldots, v\left(q_{n}\right), v\left(q_{1}\right), \ldots, v\left(q_{n}\right)\right) \in \mathrm{GL}_{2 n}(\mathbb{R})$. Let

$$
\gamma=\left(\begin{array}{rr}
0 & I_{n} \\
-I_{n} & 0
\end{array}\right) .
$$

If $m=1$ and $X \in \mathrm{SL}_{n}(\mathbb{C})$, then $c(X)=Q_{v}^{-1} X^{*,-1} Q$. If $m=2$ and $X \in \mathrm{SL}_{2 n}(\mathbb{C})$, then $c(X)=Q_{v}^{-1} \gamma X^{*,-1} \gamma^{-1} Q_{v}$. The last case easily follows from (4.1.8) and the fact that $D_{0, v}$ is split. We can identify $X^{*}\left(S_{v}\right)$ in the standard way with $\mathbb{Z}^{n m} / L$, where $L=\{(k, k, \ldots, k): k \in \mathbb{Z}\}$. It then follows easily from our calculations of the action of $c$ that if $\chi \in X^{*}\left(S_{v}\right)$ is identified with the class of the tuple $\left(a_{1}, \ldots, a_{n}\right)$ in the case $m=1$, respectively the class of the tuple $\left(a_{1}, \ldots, a_{n}, b_{1}, \ldots, b_{n}\right)$ in the case $m=2$, then ${ }^{c} \chi$ is identified with $\left(-a_{1}, \ldots,-a_{n}\right)$ or with $\left(-b_{1}, \ldots,-b_{n},-a_{1}, \ldots,-a_{n}\right)$, respectively. This, together with our formulas for $\theta$, show that $\theta_{v}$ induces $c_{S_{v}}$ on $S_{v}$, which is what we wanted to prove.

Remark 4.1.9. When $D_{0, v}$ is not split, there is also an explicit formula for $c$ that involves a matrix $\gamma$ as above, but $\gamma$ turns out to be a diagonal matrix. So in this case $\theta_{v}$ does not induce $c_{S_{v}}$ on $S_{v}$. We only care about noncompact places, hence our assumption $I_{n s} \subset I_{c}$.

Remark 4.1.10. Keep the assumptions and notation as above. For each $v \in I_{n c}$, we will construct a map $y: \mathbb{S} \rightarrow H_{v}^{\text {ad }}$ satisfying Deligne's axioms [1979, §1.2.1] and factoring through $S_{v}^{\text {ad }}$. Namely, fix $\tau: L \hookrightarrow \mathbb{C}$ an extension of $v$ to $L$, and let $w=\left.\tau\right|_{K}$ (so $w=\tau$ when $D=K=L$ ). Let $D_{w}=D \otimes_{K, w} \mathbb{C}$ and $J_{w}: D_{w} \rightarrow D_{w}$ be defined by $J_{w}(d \otimes z)=J(d) \otimes \bar{z}$. The group $H_{v}(A)$ can be identified, using the basis $\beta$, with the group of matrices $X \in \mathrm{GL}_{n}\left(D_{w} \otimes_{\mathbb{R}} A\right)$ such that ${ }^{t} X^{J_{w}} Q X=Q$ and $\operatorname{Nrd}(X)=1$. If $m=1$, let $\phi_{\tau}: D_{w} \rightarrow \mathbb{C}$ be the unique isomorphism. If $m=2$, consider the $\mathbb{C}$-algebra isomorphism $\phi_{\tau}: D_{w} \rightarrow M_{2}(\mathbb{C})$ given by

$$
\phi_{\tau}\left(\lambda \otimes_{K, w} 1\right)=\left(\begin{array}{cc}
\tau(\lambda) 0 & \\
0 & -\tau(\lambda)
\end{array}\right), \quad \phi_{\tau}\left(\mu \otimes_{K, w} 1\right)=\left(\begin{array}{cc}
0 & v\left(\mu^{2}\right) \\
1 & 0
\end{array}\right) .
$$

As above, for any $\mathbb{R}$-algebra $A$, this induces an isomorphism $\mathrm{GL}_{n}\left(D_{w} \otimes_{\mathbb{R}} A\right) \cong$ $\mathrm{GL}_{m n}\left(\mathbb{C} \otimes_{\mathbb{R}} A\right)$ taking the subgroup of diagonal matrices with entries in $L_{w} \otimes_{\mathbb{R}} A$ (where $L_{w}=L \otimes_{K, w} \mathbb{C}$ ) to the subgroup of diagonal matrices in $\mathrm{GL}_{m n}\left(\mathbb{C} \otimes_{\mathbb{R}} A\right.$ ). Moreover, the corresponding involution $X \mapsto{ }^{t} X^{J_{w}}$ gets identified with $X \mapsto$ $\gamma X^{*} \gamma^{-1}$, where if $m=1, \gamma=I_{n}$, and if $m=2, \gamma$ is the hermitian matrix defined by

$$
\gamma=\left(\begin{array}{cc}
0 & i I_{n} \\
-i I_{n} & 0
\end{array}\right)
$$

if $v\left(\lambda^{2}\right)>0$, and

$$
\gamma=\left(\begin{array}{cc}
-v\left(\mu^{2}\right) I_{n} & 0 \\
0 & I_{n}
\end{array}\right)
$$

if $v\left(\lambda^{2}\right)<0$ (note that in this case, we must have $v\left(\mu^{2}\right)>0$ ). In this way, we can write

$$
H_{v}(A) \cong\left\{X \in \mathrm{GL}_{2 n}\left(\mathbb{C} \otimes_{\mathbb{R}} A\right):\left(\gamma X^{*} \gamma^{-1}\right) Q^{\prime} X=Q^{\prime}, \operatorname{det}(X)=1\right\}
$$

where $Q^{\prime}=Q_{v}$ if $m=1$ and $Q^{\prime}=\widetilde{Q}_{v}=\operatorname{diag}\left(v\left(q_{1}\right), \ldots, v\left(q_{n}\right), \ldots, v\left(q_{1}\right), \ldots\right.$, $\left.v\left(q_{n}\right)\right)$ if $m=2$. Thus, we can identify $H_{v}$ with the special unitary group $\mathrm{SU}\left(\gamma^{-1} Q^{\prime}\right)$ of the hermitian matrix $\gamma^{-1} Q^{\prime}$, and the maximal torus $S_{v}$ is the torus of diagonal matrices. Note that $H_{v}^{\text {ad }}$ is also the adjoint group of the similitude unitary group $\operatorname{GU}\left(\gamma^{-1} Q^{\prime}\right)$. We define $y^{\prime}: \mathbb{S} \rightarrow \operatorname{GU}\left(\gamma^{-1} Q^{\prime}\right)$ as follows. For an $\mathbb{R}$-algebra $A$ and $z \in \mathbb{S}(A)$, let

$$
y_{A}^{\prime}(z)=\left(\begin{array}{cc}
\operatorname{diag}\left(y_{A}^{\prime}(z)_{1}, \ldots, y_{A}^{\prime}(z)_{n}\right) & 0 \\
0 & \operatorname{diag}\left(y_{A}^{\prime}(z)_{1}, \ldots, y_{A}^{\prime}(z)_{n}\right)
\end{array}\right)
$$

where $y_{A}^{\prime}(z)_{i}=z$ if $v\left(q_{i}\right)>0$ and $y_{A}^{\prime}(z)_{i}=\bar{z}$ if $v\left(q_{i}\right)<0$. We let $y=y^{\text {ad }}: \mathbb{S} \rightarrow H_{v}^{\text {ad }}$.

Using the explicit computation of $\gamma^{-1} Q^{\prime}$ in each case, the group $\mathrm{GU}\left(\gamma^{-1} Q^{\prime}\right)$ is isomorphic to a similitude unitary group $\operatorname{GU}(p, q)$ of a certain signature $(p, q)$ (furthermore, if $m=2$, in our case where $D_{0, v}$ is split, the signature is always ( $n, n$ ), so the group $H_{v}$ is in fact quasisplit). It is then standard that $y^{\prime}$, and hence $y$, satisfy Deligne's axioms (see for instance the Appendix of [Milne and Shih 1981]).
4.2. Groups of type $D_{l}(l \geq 4$ odd $)$. Suppose that $\mathcal{D}=D_{l}$ with $l \geq 5$ odd. Then $H=\operatorname{Spin}(B, \tau)$, where $B$ is a central simple algebra over $F$ of degree $2 l$ and $\tau$ is an orthogonal involution [Knus et al. 1998, Theorem 26.15]. The adjoint group is $H^{\text {ad }}=\mathbf{P G O}^{+}(B, \tau)$. In order to avoid introducing spin groups, we will work in this section with $H^{\text {ad }}$. Since the map $\operatorname{Aut}(H) \rightarrow \operatorname{Aut}\left(H^{\text {ad }}\right)$ is an isomorphism, an opposition involution on $H^{\text {ad }}$ will uniquely lift to an opposition involution on $H$; moreover, suppose that $S \subset H$ is a maximal torus and the involution on $H^{\text {ad }}$ preserves $S^{\text {ad }}$ and induces $c_{S_{v}^{\text {ad }}}$ on $S_{v}^{\text {ad }}$ for every $v \in I_{n c}$. Then the lifted involution on $H$ preserves $S$ and also obviously induces $c_{S_{v}^{\text {ad }}}$ on $S_{v}^{\text {ad }}$ for every $v \in I_{n c}$. This will allow us to concentrate on $H^{\text {ad }}$ and avoid spin groups.

Since $F$ is a number field, it can be shown that $B=\operatorname{End}_{D}(\Lambda)$, where $D=F$ or a quaternion division algebra over $F$ [Scharlau 1985, §8.2.3], and $\Lambda$ is a right $D$-vector space of finite dimension $n$. Let $m=\operatorname{deg}_{F} D$. Moreover, the involution $\tau: B \rightarrow B$ must be attached to a nondegenerate $F$-bilinear form $q: \Lambda \times \Lambda \rightarrow D$. In the case $D=F$ (where $\operatorname{dim}_{F} \Lambda=2 l$ ), $q$ is a symmetric bilinear form. In the case that $D$ is a quaternion division algebra (where $\operatorname{dim}_{D} \Lambda=l$ ), $q$ is a skew-hermitian form with respect to the canonical involution $\sigma: D \rightarrow D$. We will only treat the case where $D$ is a quaternion division algebra. Let $I_{s} \subset I=\operatorname{Hom}(F, \mathbb{R})$ be the set of $v: F \rightarrow \mathbb{R}$ such that $D_{v}=D \otimes_{F, v} \mathbb{R}$ is split, and let $I_{n s}$ be its complement in $I$. For $v \in I_{s}$, the skew-hermitian form $q_{v}$ on $\Lambda_{v}$ defines a nondegenerate symmetric bilinear form $b_{v}$ over a real vector space $W_{v}$ of dimension $2 n$ [Scharlau 1985], and then we have that $I_{c} \subset I_{s}$ is the set of split places where $b_{v}$ is definite. As in the Appendix of [Milne and Shih 1981] (type $D^{\circledR}$ ), we will assume that $I_{c}=I_{s}$. We call the pair $(\Lambda, q)$ a skew-hermitian space over $D$. Note that $n=l$ is odd.

Let $\beta=\left\{v_{1}, \ldots, v_{n}\right\}$ be a $D$-basis of $\Lambda$, which is $q$-orthogonal. The group $H^{\text {ad }}=\mathbf{P G O}^{+}(\Lambda, q)$ can also be seen as the adjoint group of $G=\mathbf{S O}(\Lambda, q)$, where

$$
\begin{aligned}
& G(A)=\left\{g \in \operatorname{Aut}_{D \otimes_{F} A}\left(\Lambda_{A}\right): \operatorname{Nrd}(g)=1\right. \text { and } \\
& \left.\qquad q_{A}(g(x), g(y))=q_{A}(x, y) \text { for all } x, y \in \Lambda_{A}\right\}
\end{aligned}
$$

for an $F$-algebra $A$. Here $\Lambda_{A}=\Lambda \otimes_{F} A$ and $\operatorname{Nrd}$ is the reduced norm in $\operatorname{End}_{D}(\Lambda)$. We let $S^{\prime}=S_{\beta}^{\prime} \subset G$ be the subgroup of $G$ defined as follows. For every $i=1, \ldots, n$, let $q_{i}=q\left(v_{i}, v_{i}\right)$. This is a pure quaternion in $D$, and so $L_{i}=F\left(q_{i}\right)$ is a quadratic field extension of $F$. For an $F$-algebra $A$, let
$S^{\prime}(A)=\left\{g \in G(A): g\left(v_{i} \otimes 1\right)=\left(v_{i} \otimes 1\right) \lambda_{i}\right.$ for some $\left.\lambda_{i} \in\left(L_{i} \otimes_{F} A\right)^{\times}(i=1, \ldots, n)\right\}$.

Then $S^{\prime} \subset G$ is a maximal torus of $G$, and it defines maximal tori $S \subset H$ and $S^{\text {ad }}=S^{\prime \text { ad }} \subset H^{\text {ad }}$.

We will construct involutions on $H$ modeled after our constructions for the case of type $A_{l}$. For this we need to make an analogous extra assumption.

Definition 4.2.1. We say that the skew-hermitian space $(\Lambda, q)$ over $D$ is strongly skew-hermitian if there exists a $q$-orthogonal $D$-basis $\beta=\left\{v_{1}, \ldots, v_{n}\right\}$ of $\Lambda$ and an $F$-automorphism $\alpha: D \rightarrow D$ such that $q\left(v_{i}, v_{j}\right)=-\alpha\left(q\left(v_{j}, v_{i}\right)\right)$ and $\alpha^{2}=1$.

Remark 4.2.2. Any automorphism $\alpha: D \rightarrow D$ as above must be necessarily inner, of the form $\alpha(d)=r d r^{-1}$ for some $r \in D^{\times}$such that $r^{2} \in F^{\times}$. This implies that $r \sigma(r)^{-1} \in F^{\times}$as well (because $F$ is the set of elements of $D$ fixed by $\sigma$ ). Moreover, since $q\left(v_{i}, v_{i}\right) \in D^{\times}, r$ must be a pure quaternion in $D$. As in the previous case, the existence of the basis $\beta$ will allow us to construct an explicit involution. The map $\alpha$ plays the role of the canonical conjugation of case $A_{l}$.

Suppose that $(\Lambda, q)$ is strongly hermitian, and let $\beta$ and $\alpha=\operatorname{int}(r)$ be as in the definition. We then have $\alpha \sigma=\sigma \alpha$. Let $I: \Lambda \rightarrow \Lambda$ be the $\alpha$-semilinear automorphism obtained by applying $\alpha$ to the coefficients of elements of $\Lambda$ with respect to the basis $\beta$. Then $q(I(x), I(y))=-\alpha(q(x, y))$. Let $\theta: G \rightarrow G$ be defined by $\theta_{A}(g)=I_{A} g I_{A}$ for an $F$-algebra $A$ and a $\left(D \otimes_{F} A\right)$-linear automorphism $g$ of $\Lambda \otimes_{F} A$.

Let $L=F(r)$, where $r \in D$ is as above. This is again a quadratic extension of $F$ (and a maximal subfield of $D$ ). Let $S^{\prime}$ and $S^{\text {ad }}$ be the maximal tori of $G$ and $H^{\text {ad }}$ defined above using the basis $\beta$.

Proposition 4.2.3. With the above hypotheses, the following statements are true:
(a) The map $\theta: G \rightarrow G$ is an opposition involution (and hence so is $\theta^{\mathrm{ad}}$ ).
(b) We have $\theta\left(S^{\prime}\right)=S^{\prime}$ and for every $v \in I_{n c}, \theta_{v}: G_{v} \rightarrow G_{v}$ induces $c_{S_{v}^{\prime}}$ on $S_{v}^{\prime}$.

In particular, $\theta_{v}^{\text {ad }}: H_{v}^{\text {ad }} \rightarrow H_{v}^{\text {ad }}$ induces $c_{S_{v}^{\mathrm{ad}}}$ on $S_{v}^{\text {ad }}$ for $v \in I_{n c}$.
Proof. For part (a), it suffices to see that $\theta_{E}: G_{E} \rightarrow G_{E}$ is an opposition involution, for a convenient extension $E / F$. Using the basis $\beta$ and the isomorphism $\phi$ : $D \otimes_{F} L \rightarrow M_{2}(L)$ as constructed above, we can identify $G_{L}$ as follows. Implicit in the construction of $\phi$ is the choice of a pure quaternion $s \in D$ with $r s=-r s$, and we let $t=v\left(s^{2}\right) \in \mathbb{R}$. Let $q_{i}=q\left(v_{i}, v_{i}\right)$. Since $\sigma\left(q_{i}\right)=-q_{i}$ and $r q_{i} r^{-1}=-q_{i}$, we have

$$
\phi(r)=\left(\begin{array}{rr}
r & 0 \\
0 & -r
\end{array}\right) \in \mathrm{GL}_{2}(L)
$$

and

$$
\phi\left(q_{i}\right)=\left(\begin{array}{cc}
0 & b_{i} \\
c_{i} & 0
\end{array}\right)
$$

for some $b_{i}, c_{i} \in L$. The image in $M_{2}(L)$ under $\phi$ of $L_{i} \otimes_{F} L \subset D \otimes_{F} L$ consists of the matrices in $M_{2}(L)$ of the form

$$
\left(\begin{array}{cc}
x & y b_{i} \\
y c_{i} & x
\end{array}\right)
$$

for some $x, y \in L$. Thus, the induced isomorphism $\phi: M_{n}\left(D \otimes_{F} L\right) \rightarrow M_{2 n}(L)$ sends the subalgebra of diagonal matrices $L_{1} \otimes_{F} L \times \cdots \times L_{n} \otimes_{F} L$ to the set of matrices in $M_{2 n}(L)$ of the form

$$
X=\left(\begin{array}{cc}
\operatorname{diag}\left(x_{1}, \ldots, x_{n}\right) & \operatorname{diag}\left(y_{1} b_{1}, \ldots, y_{n} b_{n}\right)  \tag{4.2.4}\\
\operatorname{diag}\left(y_{1} c_{1}, \ldots, y_{n} c_{n}\right) & \operatorname{diag}\left(x_{1}, \ldots, x_{n}\right)
\end{array}\right)
$$

with $x_{i}, y_{i} \in L$. Let

$$
\widetilde{Q}=\left(\begin{array}{cc}
0 & \operatorname{diag}\left(b_{1}, \ldots, b_{n}\right) \\
\operatorname{diag}\left(c_{1}, \ldots, c_{n}\right) & 0
\end{array}\right) \in \operatorname{GL}_{2 n}(L)
$$

Then, for any $L$-algebra $R$, writing a matrix $X \in \mathrm{GL}_{2 n}(R)$ as $X=\left(\begin{array}{cc}A & B \\ C & D\end{array}\right)$, there is an isomorphism

$$
G(R) \cong\left\{X \in \mathrm{GL}_{2 n}(R):\left(\begin{array}{rr}
{ }^{t} D & -{ }^{t} B  \tag{4.2.5}\\
-{ }^{t} C & { }^{t} A
\end{array}\right) \widetilde{Q}\left(\begin{array}{ll}
A & B \\
C & D
\end{array}\right)=\widetilde{Q}, \operatorname{det}(X)=1\right\}
$$

that takes the subgroup $S^{\prime}$ to the subgroup of matrices of the form (4.2.4) in the right-hand side. Note that the equation is equivalent to ${ }^{t} X \widetilde{Q}^{\prime} X=\widetilde{Q}^{\prime}$, where

$$
\widetilde{Q}^{\prime}=\left(\begin{array}{cc}
\operatorname{diag}\left(c_{1}, \ldots, c_{n}\right) & 0 \\
0 & -\operatorname{diag}\left(b_{1}, \ldots, b_{n}\right)
\end{array}\right)
$$

(the matrix $\widetilde{Q}^{\prime}$ is the matrix of the associated bilinear form [Scharlau 1985, §10.3]). Moreover, if

$$
\gamma=\left(\begin{array}{cc}
r I_{n} & 0 \\
0 & -r I_{n}
\end{array}\right)
$$

then $\theta_{R}(X)=\gamma X \gamma^{-1}$ for $X \in G(R)$; in block matrix terms,

$$
\theta_{R}\left(\begin{array}{ll}
A & B \\
C & D
\end{array}\right)=\left(\begin{array}{rr}
A & -B \\
-C & D
\end{array}\right)
$$

It is clear then that $\theta$ preserves $S^{\prime}$.
Let $E / L$ be a field extension such that there exist elements $e_{i}, f_{i} \in E$ with $e_{i}^{2}=c_{i}$ and $f_{i}^{2}=b_{i}$ (for example, take $E=\mathbb{C}$ with a fixed embedding of $L$ ). For elements $a_{1}, \ldots, a_{n}$, let adiag $\left(a_{1}, \ldots, a_{n}\right)$ be the antidiagonal matrix whose $(i, n+1-i)$-th entry is $a_{i}$, and let $J_{n}=\operatorname{adiag}(1, \ldots, 1)$. Let

$$
\delta=\left(\begin{array}{cc}
\operatorname{adiag}\left(e_{n}, \ldots, e_{1}\right) & \operatorname{adiag}\left(-f_{n}, \ldots,-f_{1}\right) \\
\operatorname{diag}\left(e_{1} / 2, \ldots, e_{n} / 2\right) & \operatorname{diag}\left(f_{1} / 2, \ldots, f_{n} / 2\right)
\end{array}\right) \in \operatorname{GL}_{2 n}(E)
$$

Then the map $X \mapsto \delta X \delta^{-1}$ sends $G_{E}$ (viewed inside $\mathrm{GL}_{2 n, E}$ via (4.2.5)) to the special orthogonal group $\mathrm{SO}_{2 n}$ of the matrix $J_{2 n}$ over $E$. The maximal torus $S_{E}^{\prime}$ maps to the subgroup of diagonal matrices in $\mathrm{SO}_{2 n}$, and $\theta$ becomes conjugation by the matrix

$$
\delta \gamma \delta^{-1}=\left(\begin{array}{cc}
0 & 2 r J_{n} \\
(r / 2) J_{n} & 0
\end{array}\right)
$$

inside $\mathrm{GL}_{2 n}$. We identify in the usual way $X^{*}\left(S^{\prime}\right) \cong \mathbb{Z}^{n}$. As a Borel subgroup of $G_{E}$ we take the subgroup $B$ of upper-triangular matrices belonging to $G_{E}$. The map $\theta$ sends $B$ to the subgroup $B^{-}$of lower-triangular matrices. Let $J_{2 n}^{\prime}$ be the matrix obtained from $J_{2 n}$ by swapping the rows $n$ and $n+1$. Then it is easy to see that $J_{2 n}^{\prime} \in G(E)$ and sends $B^{-}$to $B$. It follows that $\Psi_{0}(\theta)(\chi)=\chi \circ \operatorname{int}\left(J_{2 n}^{\prime}\right) \circ \theta$ for $\chi \in X^{*}\left(S^{\prime}\right)$. If $\chi$ is parametrized by $\left(a_{1}, \ldots, a_{n}\right)$, then $\Psi_{0}(\theta)(\chi)$ is parametrized by $\left(a_{1}, \ldots, a_{n-1},-a_{n}\right)=\left(a_{1}, \ldots, a_{n}\right)^{\star}[$ Bourbaki 2002, Plate IV]. Thus, $\theta: G \rightarrow G$ is an opposition involution.

Let $v: F \hookrightarrow \mathbb{R}$, and let $\tau: L \hookrightarrow \mathbb{C}$ be an extension of $v$ to $L$. If $\tau=\bar{\tau}$, then $\tau(r) \in \mathbb{R}$. Thus, $\tau(r)^{2} \in \mathbb{R}_{>0}$, and this implies that $D_{v}$ is split, so $v \in I_{s}=I_{c}$. In part (b), we only care for $v \in I_{n c}$, so suppose from now on that $\tau \neq \bar{\tau}$, so that $\tau(r) \in i \mathbb{R}_{>0}$. By the same reasoning we have that $t=v\left(s^{2}\right)<0$. We use $\tau$ to identify $G_{\mathbb{C}} \cong \mathrm{SO}_{2 n}$ as above. We first work out the induced complex conjugation on $G(\mathbb{C}) \cong \mathrm{SO}_{2 n}(\mathbb{C})$. Using the isomorphisms $D \otimes_{F, v} \mathbb{C} \simeq\left(D \otimes_{F} L\right) \otimes_{L, \tau} \mathbb{C} \cong M_{2}(\mathbb{C})$ (the last one coming from $\phi$ ), it is easy to see that complex conjugation on $D \otimes_{F, v} \mathbb{C}$ corresponds to taking a matrix $X \in M_{2}(\mathbb{C})$ to

$$
\left(\begin{array}{ll}
t & 0 \\
0 & 1
\end{array}\right)\left(\begin{array}{ll}
\overline{X_{22}} & \overline{X_{21}} \\
\overline{X_{12}} & \overline{X_{11}}
\end{array}\right)\left(\begin{array}{cc}
t^{-1} & 0 \\
0 & 1
\end{array}\right),
$$

where $t=v\left(s^{2}\right)$ as above. Note that $q_{i} \in D \subset D \otimes_{F, v} \mathbb{C}$, so this implies that $t \overline{\tau\left(c_{i}\right)}=\tau\left(b_{i}\right)$ and thus

$$
\begin{equation*}
t e_{i} / \bar{f}_{i}=-f_{i} / \bar{e}_{i} \tag{4.2.6}
\end{equation*}
$$

It follows that the induced complex conjugation on $G(\mathbb{C})$, viewed inside $\mathrm{GL}_{2 n}(\mathbb{C})$ as in (4.2.5), is given by

$$
X=\left(\begin{array}{ll}
A & B \\
C & D
\end{array}\right) \mapsto c^{\prime}(X)=\left(\begin{array}{cc}
\bar{D} & t \bar{C} \\
t^{-1} \bar{B} & \bar{A}
\end{array}\right) .
$$

Finally, we apply conjugation by $\delta$ to identify $G_{\mathbb{C}}$ with $\mathrm{SO}_{2 n}$. We only need to consider the action of $c$ on diagonal matrices. Let $X=\operatorname{diag}\left(x_{1}, \ldots, x_{n}, x_{n}^{-1}, \ldots, x_{1}^{-1}\right) \in$ $\mathrm{SO}_{2 n}(\mathbb{C})$. Then $c(X)=\delta c^{\prime}(\delta)^{-1} c^{\prime}(X) c^{\prime}(\delta) \delta^{-1}$, and a long but easy direct calculation using (4.2.6) shows that

$$
\delta c^{\prime}(\delta)^{-1}=\left(\begin{array}{cc}
2 \operatorname{adiag}\left(e_{n} / \bar{f}_{n}, \ldots, e_{1} / \bar{f}_{1}\right) & 0 \\
0 & \frac{1}{2} \operatorname{adiag}\left(f_{1} / \bar{e}_{1}, \ldots, f_{n} / \bar{e}_{n}\right)
\end{array}\right),
$$

and thus,

$$
c\left(\operatorname{diag}\left(x_{1}, \ldots, x_{n}, x_{n}^{-1}, \ldots, x_{1}^{-1}\right)\right)=\operatorname{diag}\left(\left(\overline{x_{1}}\right)^{-1}, \ldots,\left(\overline{x_{n}}\right)^{-1}, \overline{x_{n}}, \ldots, \overline{x_{1}}\right)
$$

This implies that, if $\chi \in X^{*}\left(S^{\prime}\right)$ is parametrized by $\left(a_{1}, \ldots, a_{n}\right) \in \mathbb{Z}^{n}$, then ${ }^{c} \chi$ is parametrized by $\left(-a_{1}, \ldots,-a_{n}\right)$. This is also easily seen to be the parameter of $\chi \circ \theta$, which shows that $\theta_{v}$ induces $c_{S_{v}^{\prime}}$ on $S_{v}^{\prime}$.

Remark 4.2.7. Keep the assumptions and notation as above. For each $v \in I_{n c}$, we will construct a map $y: S \rightarrow H_{v}^{\text {ad }}$ satisfying Deligne's axioms [1979, §1.2.1] and factoring through $S_{v}^{\text {ad }}$. Recall that $t=v\left(s^{2}\right)$, and let $u=v\left(r^{2}\right)$. Since $v \in I_{n c}$, by our assumptions $D_{v}$ is not split. This implies that $u<0$ and $t<0$. Let $\psi: D_{v} \rightarrow \mathbb{H}$ be the isomorphism of $\mathbb{R}$-algebras sending $r \otimes 1$ to $\sqrt{-u} e_{2}$ and $s \otimes 1$ to $\sqrt{-t} e_{3}$. Here $e_{1}, e_{2}, e_{3}$, and $e_{4}$ are the following elements of $\mathbb{H}$ :

$$
e_{1}=I_{2}, \quad e_{2}=\left(\begin{array}{rr}
i & 0 \\
0 & -i
\end{array}\right), \quad e_{3}=\left(\begin{array}{rr}
0 & 1 \\
-1 & 0
\end{array}\right), \quad e_{4}=e_{2} e_{3}
$$

As above, we can write $\psi\left(q_{i}\right)=\left(\begin{array}{rr}0 & y_{i} \\ -\bar{y}_{i} & 0\end{array}\right)$ with $y_{i} \in \mathbb{C}^{\times}$. Let

$$
T=\left(\begin{array}{cc}
0 & \operatorname{diag}\left(y_{1}, \ldots, y_{n}\right) \\
-\operatorname{diag}\left(\overline{y_{1}}, \ldots, \overline{y_{n}}\right) & 0
\end{array}\right)
$$

We then have, for an $\mathbb{R}$-algebra $R$,

$$
G_{v}(R) \cong\left\{X=\left(\begin{array}{rr}
A & B  \tag{4.2.8}\\
-\bar{B} & \bar{A}
\end{array}\right) \in \mathrm{GL}_{2 n}\left(\mathbb{C} \otimes_{\mathbb{R}} R\right): X^{*} T X=T, \operatorname{det}(X)=1\right\}
$$

The maximal torus $S^{\prime}$ corresponds to the subgroup of matrices on the right-hand side where $A=\operatorname{diag}\left(a_{1}, \ldots, a_{n}\right)$ and $B=\operatorname{diag}\left(b_{1} y_{1}, \ldots, b_{n} y_{n}\right)$ with $a_{i}, b_{i} \in R$. We can actually see $H_{v}$ as the adjoint group of $G_{v}^{\prime}$, where $G_{v}^{\prime}(R)$ is isomorphic to

$$
\left\{X=\left(\begin{array}{rr}
A & B \\
-\bar{B} & \bar{A}
\end{array}\right) \in \mathrm{GL}_{2 n}\left(\mathbb{C} \otimes_{\mathbb{R}} R\right): X^{*} T X=v(X) T, \operatorname{det}(X)=v(X)^{n}\right\}
$$

We define $y^{\prime}: \mathbb{S} \rightarrow G_{v}^{\prime}$ by the formula

$$
y_{R}^{\prime}(z)=\left(\begin{array}{cc}
\operatorname{Re}(z) I_{n} & \operatorname{diag}\left(\frac{\operatorname{Im}(z)}{\left|y_{1}\right|} y_{1}, \ldots, \frac{\operatorname{Im}(z)}{\left|y_{n}\right|} y_{n}\right) \\
\operatorname{diag}\left(-\frac{\operatorname{Im}(z)}{\left|y_{1}\right|} \bar{y}_{i}, \ldots,-\frac{\operatorname{Im}(z)}{\left|y_{n}\right|} \overline{y_{n}}\right) & \operatorname{Re}(z) I_{n}
\end{array}\right)
$$

for $z \in \mathbb{S}(R)$. Conjugating by a suitable matrix $U \in \mathrm{GL}_{2 n}(\mathbb{C})$, we can write $G_{v}^{\prime} \cong \mathrm{GO}^{*}(2 n)$ and $y$ becomes the map in the Appendix of [Milne and Shih 1981], so it satisfies Deligne's axioms, and hence also does $y=y^{\prime \text { ad }}$.

## 5. Involutions on certain Shimura varieties

In this section we combine all our previous results to prove the existence of descent data on certain Shimura varieties $\operatorname{Sh}(G, X)$. As we said before, we only consider the case where the simple groups $H_{i}$ are of type $A$ or $D^{\sharp H}$. In the previous section, we constructed opposition involutions on some of these groups, preserving a certain maximal torus $S_{i}$ and inducing complex conjugation on its characters. Furthermore, we constructed maps $y_{i, v}: \mathbb{S} \rightarrow H_{i, v}^{\text {ad }}$ for every $v \in I_{i, n c}$ satisfying Deligne's axioms [1979, §1.2.1], factoring through $S_{i, v}^{\text {ad }}$. We now show that we can always find an element $x \in X$ such that $x_{i, v}$ factors through $S_{i, v}^{\text {ad }}$ for every $i$ and $v \in I_{i, n c}$. The existence of descent data will follow by combining this with Proposition 3.4.8.
5.1. Existence of particular elements $\boldsymbol{x} \in \boldsymbol{X}$. Let $H$ be an almost simple, simply connected group over $\mathbb{R}$ (to play the role of one of the noncompact $H_{i, v}$ ). Suppose that there exist morphisms $y: S \rightarrow H^{\text {ad }}$ satisfying Deligne's axioms [1979, §1.2.1]; in particular, $H$ is absolutely almost simple. Let $D$ be the Dynkin diagram of $H_{\mathbb{C}}$ associated with a choice of maximal torus and Borel subgroup. To each $H^{\text {ad }}(\mathbb{R})$ conjugacy class $Y$ of morphisms $y$ as above, we can attach a special node $s_{Y} \in D$, and $s_{Y}=s_{Y^{\prime}}$ if and only if $Y=Y^{\prime}$.

Lemma 5.1.1. Under the above conditions, there exist at most two $H(\mathbb{R})$-conjugacy classes $Y$ of morphisms satisfying Deligne's axioms [1979, §1.2.1]. Moreover, given such a conjugacy class $Y$, any morphism satisfying these axioms must belong to either $Y$ or $Y^{-1}$.

Proof. Suppose first that $D$ is not of type $A_{l}$. This case is easy because there are not too many special nodes. Indeed, assume first that $H(\mathbb{R})$ is connected, and fix $Y$ one of the conjugacy classes. Then $s_{Y^{-1}}=s_{Y}^{\star} \neq s_{Y}$ [Deligne 1979, §1.2.8], and hence, $Y^{-1}$ and $Y$ are two distinct conjugacy classes. Suppose that $Z$ is a third conjugacy class, that is, $s_{Z}$ is neither equal to $s_{Y}$ nor to $s_{Y}^{\star}$. Again by [Deligne 1979, §1.2.8], $s_{Z} \neq s_{Z}^{\star}$, and thus, we have four distinct special nodes $s_{Y}, s_{Y}^{\star}, s_{Z}$, and $s_{Z}^{\star}$. There is no connected Dynkin diagram with four special nodes which is not of type $A_{l}$, and thus, this is a contradiction. If $H(\mathbb{R})$ is not connected, then $s_{Y}=s_{Y}^{\star}$. If $Z$ is another conjugacy class, then again by [op. cit.] we must have $s_{Z}=s_{Z}^{\star}$. But for any connected Dynkin diagram, there is at most one special node which is fixed under the opposition involution, and thus $Z=Y$.

Suppose now that $H$ is of type $A_{l}$ with $l \geq 2$, so $H=\mathrm{SU}(p, q)$ for some nonzero pair of integers $p, q$ such that $p+q=l+1$. The isomorphism $\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C} \simeq \mathbb{C} \times \mathbb{C}$ given by $z \otimes a \mapsto(z a, \bar{z} a)$ induces by projection on the first coordinate an isomorphism $H_{\mathbb{C}} \simeq \mathrm{SL}_{l+1, \mathbb{C}}$; fix the usual Borel pair here to define the Dynkin diagram. Define a morphism

$$
y_{0}: S \rightarrow H^{\mathrm{ad}}=\operatorname{PGU}(p, q)
$$

with $y_{0}(z)$ being the class of the matrix

$$
\left(\begin{array}{cc}
z I_{p} & 0 \\
0 & \bar{z} I_{q}
\end{array}\right)
$$

Then $y_{0}$ satisfies Deligne's axioms [1979, §1.2.1], and the special node $s_{0}$ attached to its $H^{\text {ad }}(\mathbb{R})$-conjugacy class $Y_{0}$ is $\alpha_{p}$. From the conjugate map $\bar{y}_{0}=y_{0}^{-1}$ we get the special node $\alpha_{q}$ associated with $Y_{0}^{-1}$. If $Y$ is another conjugacy class, say with special node $\alpha_{t}$, then there would be an isomorphism $\operatorname{PGU}(p, q) \cong \operatorname{PGU}(t, l+1-t)$ sending $Y_{0}$ or $Y_{0}^{-1}$ to $Y$. In particular, $t=p$ or $t=q$, and we conclude that there are at most two possible conjugacy classes of morphisms satisfying Deligne's axioms for the fixed form $\operatorname{PGU}(p, q)$ of $\mathrm{PGL}_{l+1, \mathbb{C}}$ (and there are exactly two in all cases except when $p=q$, when there is only one).

Going back to our general Shimura datum ( $G, X$ ), for each $i$, let $S_{i} \subset H_{i}$ be a maximal torus, $\widetilde{T}_{i}=\operatorname{Res}_{F_{i} / \mathbb{Q}} S_{i} \subset \widetilde{G}_{i}, T_{i} \subset G_{i}$ its image in $G_{i}, T^{\prime} \subset G^{\text {der }}$ the image of their product, and $T=Z^{0} T^{\prime}$. Note that $T_{\mathbb{R}}^{\text {ad }}=T_{\mathbb{R}}^{\prime \text { ad }}=\prod_{i, v} S_{i, v}^{\text {ad }}$, where $S_{i, v} \subset H_{i, v}$ and $S_{i, v}^{\mathrm{ad}}$ is its image in $H_{i, v}^{\text {ad }}$.

Lemma 5.1.2. Suppose that $T \subset G$ is the maximal torus defined above. Suppose that for each $v \in I_{i, n c}$, there exists a morphism $y_{i, v}: \mathbb{S} \rightarrow H_{i, v}^{\text {ad }}$ satisfying axioms [Deligne 1979, §1.2.1] and factoring through $S_{i, v}^{\mathrm{ad}}$. Then there exists an element $x \in X$ such that $x^{\text {ad }}$ factors through $T_{\mathbb{R}}^{\mathrm{ad}}$.

Proof. Let $z \in X$ be an arbitrary element. The previous lemma implies that $z_{i, v}$ is $H_{i, v}^{\text {ad }}(\mathbb{R})$-conjugate to a map $y_{i, v}: \mathbb{S} \rightarrow S_{i, v}^{\text {ad }}$. Thus, we can write $z_{i, v}=u_{i, v} \cdot y_{i, v}$ for $u_{i, v} \in H_{i, v}^{\text {ad }}(\mathbb{R})$. We claim that, after possibly changing the $y_{i, v}$, we can arrange for $u_{i, v}$ to be in $H_{i, v}^{\text {ad }}(\mathbb{R})^{+}$. Indeed, if $u_{i, v}$ is not in that connected component, then in particular $H_{i, v}^{\text {ad }}(\mathbb{R})$ is not connected, and thus, there is only one conjugacy class in question, with two connected components, one containing $z_{i, v}$ and the other one containing $y_{i, v}$. Thus, we only need to replace $y_{i, v}$ with $y_{i, v}^{-1}$, which also factors through $S_{i, v}^{\text {ad }}$. For $v \in I_{i, c}$, let $u_{i, v}=1$. It follows that $u=\left(u_{i, v}\right) \in G^{\text {ad }}(\mathbb{R})^{+}$, and thus, there exists $g \in G(\mathbb{R})$ lifting $u$. Let $x=g^{-1} \cdot z \in X$, so that $x^{\text {ad }}=\left(y_{i, v}\right)$, which factors through $T_{\mathbb{R}}^{\text {ad }}$ as desired.
5.2. The main theorem. In this subsection we put all the ingredients together to obtain the main theorem on the existence of involutions of $G$ taking $X$ to $\bar{X}$.

Definition 5.2.1. The Shimura datum $(G, X)$ is said to be strongly of type $\left(A D^{\sharp \Perp}\right)$ if each of the groups $H_{i}$ is either of type $A_{l}$ with $l \geq 2$ and attached to a strongly hermitian space (as in Definition 4.1.4), or of type $D_{l}$ with $l \geq 5$ odd and attached to a strongly skew-hermitian space (as in Definition 4.2.1).

For example, a Shimura variety defined by a similitude unitary group attached to a hermitian space over a CM field is strongly of type $\left(A D^{\sharp H}\right)$. Note however that the definition only restricts the semisimple part of $G$.

Theorem 5.2.2. Suppose that $(G, X)$ is strongly of type $\left(A D^{\mathbb{H}}\right)$. Then there exists an involution $\theta: G \rightarrow G$ such that $\theta(X)=\bar{X}$, and hence, there exists a model of $\operatorname{Sh}(G, X)$ over $E^{+}$as in Theorem 2.3.1.

Proof. In Subsections 4.1 and 4.2, we constructed for every $i$, an opposition involution $\theta_{i}: H_{i} \rightarrow H_{i}$ and a maximal torus $S_{i} \subset H_{i}$ such that $\theta_{i}\left(S_{i}\right)=S_{i}$ and $\theta_{i, v}^{\text {ad }}$ induces $c_{S_{i, v}^{\text {ad }}}$ for every $v \in I_{i, n c}$. Moreover, by Remarks 4.1.10 and 4.2.7, for every $i$ and $v \in I_{i, n c}$, there is a map $y_{i, v}: \mathbb{S} \rightarrow H_{i, v}^{\text {ad }}$ satisfying Deligne's axioms [1979, §1.2.1], and factoring through $S_{i, v}^{\text {ad }}$. The result then follows by combining Proposition 3.4.8 and Lemma 5.1.2.

Remark 5.2.3. The conclusion of the previous theorem holds in other cases as well. For instance, if $G$ is adjoint and there exists an opposition involution $\theta: G \rightarrow G$ (which is always the case if $G$ is also quasisplit, for example), then by the adjointness of $G$, we conclude that $\theta(X)=\bar{X}$. On the other hand, the cases that we considered in this paper are concretely given by simple algebras, and thus are intimately related to moduli problems, even though we do not use the moduli interpretation explicitly. It is an interesting problem to consider factors of other types, for instance of type $E_{6}$, and analyze whether it is possible to construct opposition involutions with the desired properties in some of these cases. We plan to investigate this problem in the future.

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# A subspace theorem for subvarieties 

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#### Abstract

We establish a height inequality, in terms of an (ample) line bundle, for a sum of subschemes located in $\ell$-subgeneral position in an algebraic variety, which extends a result of McKinnon and Roth (2015). The inequality obtained in this paper connects the result of McKinnon and Roth (the case when the subschemes are points) and the results of Corvaja and Zannier (2004), Evertse and Ferretti (2008), Ru (2017), and Ru and Vojta (2016) (the case when the subschemes are divisors). Furthermore, our approach gives an alternative short and simpler proof of McKinnon and Roth's result.


## 1. Introduction and statements

McKinnon and M. Roth [2015] introduced the approximation constant $\alpha_{x}(L)$ to an algebraic point $x$ on an algebraic variety $V$ with an ample line bundle $L$. The invariant $\alpha_{x}(L)$ measures how well $x$ can be approximated by rational points on $V$ with respect to the height function associated to $L$. They showed that $\alpha_{x}(L)$ is closely related to the Seshadri constant $\epsilon_{x}(L)$ measuring the local positivity of $L$ at $x$. They also showed that the invariant $\alpha_{x}(L)$ can be computed through another invariant $\beta_{x}(L)$ in the height inequality they established (see Theorem 5.1 and Theorem 6.1 in [McKinnon and Roth 2015]). By computing the Seshadri constant $\epsilon_{x}(L)$ for the case of $V=\mathbb{P}^{1}$, their result recovers Roth's theorem, so the height inequality they established can be viewed as a generalization of this theorem to arbitrary projective varieties.

In this paper, we provide a simpler proof of the above results. Furthermore, we extend the results from the points of a projective variety to subschemes. The generalized result in terms of subschemes connects, as well as gives a clearer explanation to, the above mentioned result of McKinnon and Roth with the recent Diophantine approximation results in terms of the divisors obtained in [Corvaja and Zannier 2004; Evertse and Ferretti 2008; Levin 2014; Ru and Vojta 2016; Ru 2017].

We now state our result. Let $V$ be a projective variety defined over a number field $k$.

[^16]Definition 1.1. Let $L$ be a line bundle over $V$ with $h^{0}(V, N L) \geq 1$ for $N$ big enough. Let $Y$ be a proper closed subscheme of $V$ and $\pi: \widetilde{V} \rightarrow V$ be the blow-up along $Y$, and $E$ be the exceptional divisor. We define

$$
\beta_{L, Y}:=\liminf _{N \rightarrow \infty} \frac{\sum_{m=1}^{\infty} h^{0}\left(\tilde{V}, N \pi^{*} L-m E\right)}{N \cdot h^{0}(V, N L)}
$$

Remark 1.2. (a) If $Y$ is an effective Cartier divisor, then the blow-up is an isomorphism. Without loss of generality, we let $\pi$ be the identity map, $\widetilde{V}=V$ and $E=Y$.
(b) Let $D$ be an effective divisor on $V$, we define $\beta_{D, Y}:=\beta_{\mathcal{O}(D), Y}$, where $\mathcal{O}(D)$ is the line sheaf associated to $D$.
(c) In the case when $L$ is big, the $\lim _{N \rightarrow \infty}$ in the definition above exists. Indeed (see [McKinnon and Roth 2015, pp. 544-545]), we have

$$
\beta_{L, Y}=\int_{0}^{\gamma_{\mathrm{eff}}} \frac{\operatorname{Vol}\left(L_{\gamma}\right)}{\operatorname{Vol}(L)} d \gamma
$$

where $L_{\gamma}:=\pi^{*} L-\gamma E$ and $\gamma_{\text {eff }}=\sup \left\{\gamma \geq 0 \mid L_{\gamma}\right.$ is effective $\}$.
Definition 1.3. We say that the closed subschemes $Y_{1}, \ldots, Y_{q}$ of a projective variety $V$ are in $\ell$-subgeneral position if, for any $x \in V$, there are at most $\ell$ subschemes among $Y_{1}, \ldots, Y_{q}$ which contain $x$.
Remark 1.4. In the case that $Y_{1}=y_{1}, \ldots, Y_{q}=y_{q}$ are points (as in [McKinnon and Roth 2015]), the condition that $y_{1}, \ldots, y_{q}$ are distinct implies that $Y_{1}, \ldots, Y_{q}$ are in 1 -subgeneral position (i.e., with $\ell=1$ ).

We establish the following result.
Main Theorem. Let $k$ be a number field and $M_{k}$ be the set of places on $k$. Let $S \subset M_{k}$ be a finite subset containing all archimedean places. Let $V$ be a projective variety defined over $k$ and $Y_{1}, \ldots, Y_{q}$ be closed subschemes of $V$ defined over $k$ in $\ell$-subgeneral position. For any $v \in S$, choose a local Weil function $\lambda_{Y_{j}, v}$ for each $Y_{j}, 1 \leq j \leq q$. Let L be a big line bundle. Then for any $\epsilon>0$

$$
\begin{equation*}
\sum_{v \in S} \sum_{i=1}^{q} \lambda_{Y_{i}, v}(x) \leq \ell\left(\max _{1 \leq i \leq q}\left\{\beta_{L, Y_{i}}^{-1}\right\}+\epsilon\right) h_{L}(x) \tag{1-1}
\end{equation*}
$$

holds for all $x$ outside a proper Zariski-closed subset $Z$ of $V(k)$.
The following corollary of our main theorem recovers the main result of [McKinnon and Roth 2015]. The proof will be given in Section 3.
Corollary $\mathbf{1 . 5}$ [McKinnon and Roth 2015, Theorem 6.1]. Let $V$ be a projective variety over $k$. Then for any ample line bundle $L$ and any $x \in V(\bar{k})$ either

- $\alpha_{x}(L) \geq \beta_{L, x}$ or
- there exists a proper subvariety $Z \subset V$, irreducible over $\bar{k}$, with $x \in Z(\bar{k})$ so that $\alpha_{x, V}(L)=\alpha_{x, Z}(L \mid Z)$, i.e., " $\alpha_{x}(L)$ is computed on a proper subvariety of $V^{\prime \prime}$,
where $\alpha_{x}(L)$ is the approximation constant defined in [McKinnon and Roth 2015, Definition 2.9], and $\beta_{L, x}$ is defined in Definition 1.1 (with $Y$ taken as a point $x$ ).

We will show in Lemma 2.2 that for any line bundle $L, x \in V$

$$
\begin{equation*}
\beta_{L, x} \geq \frac{n}{n+1} \epsilon_{x}(L) \tag{1-2}
\end{equation*}
$$

where $n=\operatorname{dim} V$. We note that the Seshadri constant $\epsilon_{x}(L)$ does not decrease when restricting to a subvariety [McKinnon and Roth 2015, Proposition 3.4], so we can use induction to further get, from Corollary 1.5 and (1-2), the following result.
Corollary 1.6 [McKinnon and Roth 2015, Theorem 6.2, alternative statement]. Let $V$ be a projective variety over $k$. Let $L$ be any ample line bundle and choose any $x \in V(\bar{k})$. Then for any $\delta>0$, there are only finitely many solutions $y \in V(k)$ to

$$
d_{v}(x, y)<H_{L}(y)^{-\left((n+1) /\left(n \epsilon_{x}(L)\right)+\delta\right)}
$$

In the case when $V=\mathbb{P}^{n}$ and $L=\mathcal{O}_{\mathbb{P}^{n}}(1)$, we have $\epsilon_{x}(L)=1$ for all $x \in \mathbb{P}^{n}$ (see [McKinnon and Roth 2015, Lemma 3.3]). Therefore the above result generalizes the theorem of Roth.

We now turn to another extreme case when the subschemes $Y_{1}, \ldots, Y_{q}$ are effective Cartier divisors $D_{1}, \ldots, D_{q}$. Let $D:=D_{1}+\cdots+D_{q}$. Assume that each $D_{j}$ is linearly equivalent to a fixed ample divisor $A$. Then we have the following relation of height functions $h_{D}=q h_{A}+O(1)$. On the other hand, by the Riemann-Roch theorem, with $n=\operatorname{dim} V$,

$$
h^{0}(N D)=h^{0}(q N A)=\frac{(q N)^{n} A^{n}}{n!}+o\left(N^{n}\right)
$$

and

$$
h^{0}\left(N D-m D_{j}\right)=h^{0}((q N-m) A)=\frac{(q N-m)^{n} A^{n}}{n!}+o\left(N^{n}\right)
$$

Thus

$$
\sum_{m \geq 1} h^{0}\left(N D-m D_{j}\right)=\frac{A^{n}}{n!} \sum_{l=0}^{q N-1} l^{n}+o\left(N^{n+1}\right)=\frac{A^{n}(q N-1)^{n+1}}{(n+1)!}+o\left(N^{n+1}\right)
$$

Hence

$$
\beta_{D, D_{j}}=\lim _{N \rightarrow \infty} \frac{\frac{A^{n}(q N-1)^{n+1}}{(n+1)!}+o\left(N^{n+1}\right)}{N \frac{(q N)^{n} A^{n}}{n!}+o\left(N^{n+1}\right)}=\frac{q}{n+1}
$$

Thus the Main Theorem, together with the above computation, implies the following result of Chen, Ru, and Yan [2012] (see also [Corvaja and Zannier 2006]).

Theorem 1.7. Let $k$ be a number field and $M_{k}$ the set of places on $k$. Let $S \subset M_{k}$ be a finite subset containing all archimedean places. Let $V$ be a projective variety of dimension $n$ defined over $k$. Let $D_{1}, \ldots, D_{q}$ be effective Cartier divisors in $\ell$-subgeneral position on $V$. Assume that each $D_{j}, 1 \leq j \leq q$, is linearly equivalent to a fixed ample divisor $A$. For any $v \in S$, choose a Weil function $\lambda_{D_{j}, v}$ for each $D_{j}, 1 \leq j \leq q$. Then for any $\epsilon>0$

$$
\begin{equation*}
\sum_{v \in S} \sum_{i=1}^{q} \lambda_{D_{i}, v}(x) \leq \ell(n+1+\epsilon) h_{A}(x) \tag{1-3}
\end{equation*}
$$

holds for all $x$ outside a proper Zariski-closed subset $Z$ of $V(k)$. In particular, if $D_{1}, \ldots, D_{q}$ are in general position on $V$, then the inequality

$$
\begin{equation*}
\sum_{v \in S} \sum_{i=1}^{q} \lambda_{D_{i}, v}(x) \leq n(n+1+\epsilon) h_{A}(x) \tag{1-4}
\end{equation*}
$$

holds for all but finitely many $x \in V(k)$.
In the general case when $D_{1}, \ldots, D_{q}$ are only assumed to be big and nef, we can also compute $\beta_{D, D_{j}}$. The details will be carried out in the next section.

We note that recently the first named author and P. Vojta [2016] obtained the following sharp result in the case when $D_{1}, \ldots, D_{q}$ are in general position and when $V$ is Cohen-Macaulay (for example when $V$ is smooth).

Theorem 1.8 [ Ru and Vojta 2016]. Let $k$ be a number field and $M_{k}$ be the set of places on $k$. Let $S \subset M_{k}$ be a finite subset containing all archimedean places. Let $V$ be a projective variety defined over $k$. Assume that $V$ is Cohen-Macaulay. Let $D_{1}, \ldots, D_{q}$ be effective Cartier divisors in general position on $V$. For any $v \in S$, choose a Weil function $\lambda_{D_{j}, v}$ for each $D_{j}, 1 \leq j \leq q$. Let $L$ be a line bundle on $V$ with $h^{0}(V, N L) \geq 1$ for $N$ big enough. Then for any $\epsilon>0$

$$
\begin{equation*}
\sum_{v \in S} \sum_{i=1}^{q} \lambda_{D_{i}, v}(x) \leq\left(\max _{1 \leq i \leq q}\left\{\beta_{L, D_{i}}^{-1}\right\}+\epsilon\right) h_{L}(x) \tag{1-5}
\end{equation*}
$$

holds for all $x$ outside a proper Zariski-closed subset $Z$ of $V(k)$.
Theorem 1.8, together with the above computation, recovers the result of [Evertse and Ferretti 2002; 2008] in the case when $V$ is smooth.

## 2. Computation of the constant $\boldsymbol{\beta}_{L, Y}$

We first compute the constant $\beta_{L, y}$, i.e., we let $Y=y$ be a point in $V(k)$. The following lemma is a reformulation of Lemma 4.1 in [McKinnon and Roth 2015].

Lemma 2.1. Let $V$ be a projective variety and $x$ be a point in $V$. Let $\pi: \widetilde{V} \rightarrow V$ be the blow-up along $x$, and $E$ be the exceptional divisor. Let $L$ be an ample line bundle and $m$ a positive integer. Then
(i) $h^{0}\left(\tilde{V}, N \pi^{*} L-m E\right)=0$ if $m>N \cdot \gamma_{\mathrm{eff}, x}$, where $\gamma_{\mathrm{eff}, x}$ is defined in [McKinnon and Roth 2015], and
(ii) $h^{0}\left(\tilde{V}, N \pi^{*} L-m E\right) \geq h^{0}(V, N L)-m^{n} \operatorname{mult}_{x} V / n!+O\left(N^{n-1}\right)$ for $N \gg 0$.

Proof. Write $h^{0}\left(\tilde{V}, N \pi^{*} L-m E\right)=h^{0}\left(\tilde{V}, N \pi^{*} L-N \cdot \gamma E\right)$, where $\gamma=m / N$. The argument in [McKinnon and Roth 2015] shows that $h^{0}\left(\widetilde{V}, N \pi^{*} L-m E\right) \geq$ $h^{0}(V, N L)-m^{n} \operatorname{mult}_{x} V / n!+O\left(N^{n-1}\right)$.

The following is a restatement of Corollary 4.2 in [McKinnon and Roth 2015].
Lemma 2.2. For any ample line bundle $L, x \in V$ and positive integer $m$, we have

$$
\beta_{L, x} \geq \frac{n}{n+1}\left(\frac{L^{n}}{\operatorname{mult}_{x} V}\right)^{1 / n} \geq \frac{n}{n+1} \epsilon_{x}(L)
$$

Proof. Choose a sufficiently large $N$. By Lemma 2.1 and the Riemann-Roch theorem,

$$
\begin{equation*}
h^{0}\left(\tilde{V}, \pi^{*} N L-m E\right) \geq h^{0}(V, N L)\left(1-\frac{\operatorname{mult}_{x} V}{L^{n}}\left(\frac{m}{N}\right)^{n}\right)+O\left(N^{n-1}\right) \tag{2-1}
\end{equation*}
$$

The right-hand side of (2-1) is less than zero when $m>u=\left[N\left(L^{n} / \text { mult }_{x} V\right)^{1 / n}\right]$, hence

$$
\begin{equation*}
\sum_{m=1}^{\infty} h^{0}\left(\widetilde{V}, \pi^{*} N L-m E\right) \geq h^{0}(V, N L) \sum_{m=1}^{u}\left(1-\frac{\operatorname{mult}_{x} V}{L^{n}}\left(\frac{m}{N}\right)^{n}\right)+O\left(N^{n}\right) \tag{2-2}
\end{equation*}
$$

Consequently,

$$
\begin{align*}
\beta_{L, x} & \geq \frac{1}{N} \sum_{m=1}^{u}\left(1-\frac{\operatorname{mult}_{x} V}{L^{n}}\left(\frac{m}{N}\right)^{n}\right)+O\left(\frac{1}{N}\right) \\
& =\frac{1}{N}\left(u-\frac{\operatorname{mult}_{x} V}{L^{n}} \cdot \frac{u^{n+1}}{(n+1) N^{n}}\right)+O\left(\frac{1}{N}\right) \\
& \geq \frac{n u}{(n+1) N}+O\left(\frac{1}{N}\right) \tag{2-3}
\end{align*}
$$

Let $N$ run through all sufficiently large integers. Then we have

$$
\beta_{L, x} \geq \frac{n}{n+1}\left(\frac{L^{n}}{\operatorname{mult}_{x} V}\right)^{1 / n}
$$

Next we consider the case when $Y_{j}=D_{j}, 1 \leq j \leq q$, are effective big and nef Cartier divisors on $V$.

Definition 2.3. Suppose $X$ is a complete variety of dimension $n$. Let $D_{1}, \ldots, D_{q}$ be effective Cartier divisors on $X$ and let $D=D_{1}+D_{2}+\cdots+D_{q}$. We say that $D$ has equidegree with respect to $D_{1}, D_{2}, \ldots, D_{q}$ if $D_{i} \cdot D^{n-1}=D^{n} / q$ for all $i=1, \ldots, q$.

Lemma 2.4 [Levin 2009, Lemma 9.7]. Let $V$ be a projective variety of dimension $n$. If $D_{j}, 1 \leq j \leq q$, are big and nef Cartier divisors on $V$, then there exist positive real numbers $r_{j}$ such that $D=\sum_{j=1}^{q} r_{j} D_{j}$ has equidegree.

Since divisors $r_{j} D_{j}$ and $D_{j}$ have the same support, the above lemma tells us that we can always make the given big and nef divisors have equidegree without changing their supports. So now we assume that $D:=D_{1}+\cdots+D_{q}$ is of equidegree. To compute $\beta_{D, D_{j}}$ for $j=1, \ldots, q$, we use the following lemma.

Lemma 2.5 [Autissier 2009, Lemma 4.2]. Suppose E is a big and base-point free Cartier divisor on a projective variety $V$ and $F$ is a nef Cartier divisor on $V$ such that $F-E$ is also nef. Let $\delta>0$ be a positive real number. Then, for any positive integers $N$ and $m$ with $1 \leq m \leq \delta N$, we have
$h^{0}(N F-m E)$
$\geq \frac{F^{n}}{n!} N^{n}-\frac{F^{n-1} \cdot E}{(n-1)!} N^{n-1} m+\frac{(n-1) F^{n-2} \cdot E^{2}}{n!} N^{n-2} \min \left\{m^{2}, N^{2}\right\}+O\left(N^{n-1}\right)$, where the implicit constant depends on $\beta$.

We compute $\sum_{m \geq 1} h^{0}\left(N D-m D_{i}\right)$ for each $1 \leq i \leq q$. Let $n=\operatorname{dim} V$ and assume that $n \geq 2$. Let $b=D^{n} /\left(n D^{n-1} \cdot D_{i}\right)$ and $A=(n-1) D^{n-2} \cdot D_{i}^{2}$. Then, by Lemma 2.5,

$$
\begin{aligned}
\sum_{m=1}^{\infty} h^{0}(N & \left.D-m D_{i}\right) \\
& \geq \sum_{m=1}^{[b N]}\left(\frac{D^{n}}{n!} N^{n}-\frac{D^{n-1} \cdot D_{i}}{(n-1)!} N^{n-1} m+\frac{A}{n!} N^{n-2} \min \left\{m^{2}, N^{2}\right\}\right)+O\left(N^{n}\right) \\
& \geq\left(\frac{D^{n}}{n!} b-\frac{D^{n-1} \cdot D_{i}}{(n-1)!} \frac{b^{2}}{2}+\frac{A}{n!} g(b)\right) N^{n+1}+O\left(N^{n}\right) \\
& =\left(\frac{b}{2}+\frac{A}{D^{n}} g(b)\right) D^{n} \frac{N^{n+1}}{n!}+O\left(N^{n}\right) \\
& =\left(\frac{b}{2}+\alpha\right) N h^{0}(N D)+O\left(N^{n}\right)
\end{aligned}
$$

where $\alpha:=g(b) A / D^{n}$ and $g: \mathbb{R}^{+} \rightarrow \mathbb{R}^{+}$is the function given by $g(x)=x^{3} / 3$ if $x \leq 1$ and $g(x)=x-\frac{2}{3}$ for $x \geq 1$. Now from the assumption of equidegree $D_{i} \cdot D^{n-1}=D^{n} / q$, so $b=q / n$. Moreover, $\alpha>0$ since $\operatorname{dim} V \geq 2$ and the $D_{i}$ are
big and nef divisors. Hence

$$
\beta_{D, D_{i}}=\liminf _{N} \frac{\sum_{m \geq 1} h^{0}\left(N D-m D_{i}\right)}{N h^{0}(N D)} \geq \frac{b}{2}+\alpha .
$$

Thus we have proved the following.
Proposition 2.6. Let $V$ be a projective variety of $\operatorname{dim} V \geq 2$ and assume that $D:=\sum_{j=1}^{q} D_{j}$ has equidegree with respect to $D_{1}, \ldots, D_{q}$ which are big and nef. Then

$$
\beta_{D, D_{i}}=\liminf _{N} \frac{\sum_{m \geq 1} h^{0}\left(N D-m D_{i}\right)}{N h^{0}(N D)}>\frac{q}{2 n}+\alpha,
$$

where $\alpha$ is a computable positive number.
Proposition 2.6, together with the Main Theorem, implies the following result.
Theorem 2.7 [Hussein and Ru 2018]. Let $k$ be a number field and let $S \subseteq M_{k}$ be a finite set containing all archimedean places. Let $V$ be a projective variety of dimension $\geq 2$ over $k$ and let $D_{1}, \ldots, D_{q}$ be effective, big, and nef Cartier divisors on $V$ defined over $k$, located in $\ell$-subgeneral position. Let $r_{i}>0$ be real numbers such that $D:=\sum_{i=1}^{q} r_{i} D_{i}$ has equidegree (such numbers exist due to Lemma 2.4). Then, for $\epsilon_{0}>0$ small enough, the inequality

$$
\sum_{v \in S} \sum_{j=1}^{q} r_{j} \lambda_{D_{i}, v}(x)<\ell\left(\frac{2 \operatorname{dim} V}{q}-\epsilon_{0}\right)\left(\sum_{j=1}^{q} r_{j} h_{D_{j}}(x)\right)
$$

holds for all $x \in V(k)$ outside a proper Zariski-closed subset of $V$.

## 3. Proof of the Main Theorem

We first recall some basic properties of local Weil functions associated to closed subschemes from [Silverman 1987, Section 2]. We assume that the readers are familiar with the notion of Weil functions associated to divisors (see [Lang 1983, Chapter 10], [Hindry and Silverman 2000, B.8] or [Silverman 1987, Section 1]).

Let $Y$ be a closed subscheme on a projective variety $V$ defined over $k$. Then one can associate to each place $v \in M_{k}$ a function

$$
\lambda_{Y, v}: V \backslash \operatorname{supp}(Y) \rightarrow \mathbb{R}
$$

satisfying some functorial properties (up to an $M_{k}$-constant) described in [Silverman 1987, Theorem 2.1]. Intuitively, for each $P \in V$ and $v \in M_{k}$,

$$
\lambda_{Y, v}(P)=-\log (v \text {-adic distance from } P \text { to } Y) .
$$

The following lemma indicates the existence of local Weil functions.

Lemma 3.1. Let $Y$ be a closed subscheme of $V$. There exist effective divisors $D_{1}, \ldots, D_{r}$ such that

$$
Y=\cap D_{i}
$$

Proof. See Lemma 2.2 from [Silverman 1987].
Definition 3.2. Let $k$ be a number field, and $M_{k}$ be the set of places on $k$. Let $V$ be a projective variety over $k$ and let $Y \subset V$ be a closed subscheme of $V$. We define the (local) Weil function for $Y$ with respect to $v \in M_{k}$ as

$$
\begin{equation*}
\lambda_{Y, v}=\min _{i}\left\{\lambda_{D_{i}, v}\right\}, \tag{3-1}
\end{equation*}
$$

when $Y=\cap D_{i}$ (such $D_{i}$ exist according to the above lemma).
Lemma 3.3 [Vojta 1987, Lemma 2.5.2; Silverman 1987, Theorem 2.1(h)]. Let $Y$ be a closed subscheme of $V$ and let $\widetilde{V}$ be a blow-up of $V$ along $Y$ with exceptional divisor $E=\pi^{*} Y$. Then $\lambda_{Y, v}(\pi(P))=\lambda_{E, v}(P)+O_{v}(1)$ for $P \in \widetilde{V}$.

Note that in the original statement of Lemma 2.5.2 in [Vojta 1987], $V$ is assumed to be smooth, but from the proof it is easy to see that it works for a general projective variety from Theorem 2.1(h) in [Silverman 1987].

For our purpose, it suffices to fix a choice of local Weil functions $\lambda_{Y_{i}, v}$ for each $1 \leq i \leq q$ and $v \in S$.
Lemma 3.4. Let $Y_{1}, \ldots, Y_{q}$ be closed subschemes of a projective variety $V$ in $\ell$-subgeneral position. Then

$$
\begin{equation*}
\sum_{i=1}^{q} \lambda_{Y_{i}, v}(x) \leq \max _{I} \sum_{j \in I} \lambda_{Y_{j}, v}(x)+O_{v}(1) \tag{3-2}
\end{equation*}
$$

where I runs over all index subsets of $\{1, \ldots, q\}$ with $\ell$ elements for all $x \in V(k)$. Proof. Let $\left\{i_{1}, \ldots, i_{q}\right\}=\{1, \ldots, q\}$. Since the $Y_{i}, 1 \leq i \leq q$, are in $\ell$-subgeneral position, $\bigcap_{t=1}^{\ell+1} Y_{i_{t}}=\varnothing$. Then

$$
\begin{equation*}
\min _{1 \leq i \leq \ell+1}\left\{\lambda_{Y_{i}, v}\right\}=\left\{\lambda_{\bigcap_{t=1}^{\ell+1} Y_{i}, v}\right\}=O_{v}(1) . \tag{3-3}
\end{equation*}
$$

We note that the first equality follows from (3-1), the definition of the local Weil function; and the second equality follows from Corollary 3.3 in [Lang 1983, Chapter 10]. For $x$ with the following ordering
we have

$$
\lambda_{Y_{i_{1}}, v}(x) \geq \lambda_{Y_{i_{2}}, v}(x) \geq \cdots \geq \lambda_{Y_{i q}, v}(x)
$$

$$
\sum_{i=1}^{q} \lambda_{Y_{i}, v}(x)=\sum_{i=1}^{\ell} \lambda_{Y_{i}, v}(x)+O_{v}(1)
$$

Then assertion (3-2) follows directly as the number of subvarieties under consideration is finite.

We also need the following generalized Schmidt subspace theorem.
Theorem 3.5 [Ru and Vojta 2016, Theorem 2.7]. Let k be a number field, $S$ be a finite set of places of $k$ containing all archimedean places, $X$ be a complete variety over $k, D$ be a Cartier divisor on $X, W$ be a nonzero linear subspace of $H^{0}(X, \mathcal{O}(D)), s_{1}, \ldots, s_{q}$ be nonzero elements of $W, \epsilon>0$, and $c \in \mathbb{R}$. For each $j=1, \ldots, q$, let $D_{j}$ be the Cartier divisor $\left(s_{j}\right)$ and $\lambda_{D_{j}}$ be a Weil function for $D_{j}$. Then there is a proper Zariski-closed subset $Z$ of $X$, depending only on $k, S, X$, $D, W, s_{1}, \ldots, s_{q}, \epsilon, c$, and the choices of Weil and height functions, such that the inequality

$$
\begin{equation*}
\sum_{v \in S} \max _{J} \sum_{j \in J} \lambda_{D_{j}, v}(x) \leq(\operatorname{dim} W+\epsilon) h_{D}(x)+c \tag{3-4}
\end{equation*}
$$

holds for all $x \in(X \backslash Z)(k)$. Here the set $J$ ranges over all subsets of $\{1, \ldots, q\}$ such that the sections $\left(s_{j}\right)_{j \in J}$ are linearly independent.

We are now ready to prove the Main Theorem.
Proof of the Main Theorem. Let $\delta>0$ be a sufficiently small number. We may choose a sufficiently large integer $N$ such that, for $i=1, \ldots, q$,

$$
\begin{equation*}
\sum_{m=1}^{\infty} h^{0}\left(\tilde{V}_{i}, N \pi^{*} L-E_{i}\right) \geq\left(\beta_{L, Y_{i}}-\delta\right) N h^{0}(V, N L) \tag{3-5}
\end{equation*}
$$

where $\pi_{i}: \widetilde{V}_{i} \rightarrow V$ is the blow-up at $Y_{i}$ and $E_{i}=\pi^{-1}\left(Y_{i}\right)$ is he exceptional divisor of $\pi_{i}$.

Let $x \in V(k)$ and $v \in S$. Since the $Y_{i}, 1 \leq i \leq q$, are in $\ell$-subgeneral position, it follows from Lemma 3.4 that

$$
\begin{equation*}
\sum_{i=1}^{q} \lambda_{Y_{i}, v}(x) \leq \ell \lambda_{Y_{i_{0}}, v}(x)+O_{v}(1) \tag{3-6}
\end{equation*}
$$

for some $i_{0}$ with $1 \leq i_{0} \leq q$, where the constant $O_{v}(1)$ is independent of $x$. Note that $i_{0}$ depends on the point $x$, but $O_{v}(1)$ is independent of $x$.

Write $\widetilde{V}_{i_{0}}$ as $\widetilde{V}, \pi_{i_{0}}$ as $\pi$ and $E_{i_{0}}$ as $E$. We consider the following filtration.

$$
\begin{equation*}
H^{0}\left(\tilde{V}, \pi^{*} N L\right) \supseteq H^{0}\left(\tilde{V}, \pi^{*} N L-E\right) \supseteq H^{0}\left(\tilde{V}, \pi^{*} N L-2 E\right) \supseteq \cdots \tag{3-7}
\end{equation*}
$$

We identify $H^{0}(V, N L)$ with $H^{0}\left(\widetilde{V}, \pi^{*} N L\right)$ as vector spaces (note: according to the footnote on page 553 in [McKinnon and Roth 2015], if $X$ is not normal, then $H^{0}(V, N L)$ may only be a proper subspace of $H^{0}\left(\tilde{V}, \pi^{*} N L\right)$. However, since the volume is a birational constant, the asymptotic calculations go through without change). Choose regular sections $s_{1}, \ldots, s_{M} \in H^{0}(V, N L)$ successively so that their pull-back $\pi^{*} s_{1}, \ldots, \pi^{*} s_{M} \in H^{0}\left(\tilde{V}, \pi^{*} N L\right)$ form a basis associated to this
filtration, where $M=h^{0}\left(\tilde{V}, N \pi^{*} L\right)$. For a section $\pi^{*} s \in H^{0}\left(\tilde{V}, \pi^{*} N L-m E\right)$ (regarded as a subspace of $H^{0}\left(\tilde{V}, \pi^{*} N L\right)$ ) we have

$$
\begin{equation*}
\operatorname{div}\left(\pi^{*} s\right) \geq m E \tag{3-8}
\end{equation*}
$$

Hence, $\lambda_{\left(\pi^{*} s\right), v} \geq m \lambda_{E, v}+O_{v}(1)$. Note that although $O_{v}(1)$ here depends on $i_{0}$ (which depends on $x$ ), there are $q$ choices of such $i_{0}$ and $V$ is compact, so we can again make $O_{v}(1)$ independent of $x$. Therefore, also using Lemma 3.3 and (3-5), $\sum_{j=1}^{M} \lambda_{\left(\pi * s_{j}\right), v}$
$\geq \sum_{m=1}^{\infty} m\left(h^{0}\left(\tilde{V}, \pi^{*} N L-m E\right)-h^{0}\left(\tilde{V}, \pi^{*} N L-(m+1) E\right)\right) \lambda_{E, v}+O_{v}(1)$
$=\sum_{m=1}^{\infty} m\left(h^{0}\left(\tilde{V}, \pi^{*} N L-m E\right)-h^{0}\left(\tilde{V}, \pi^{*} N L-(m+1) E\right)\right) \lambda_{Y_{i_{0}}, v} \circ \pi+O_{v}(1)$
$=\sum_{m=1}^{\infty} h^{0}\left(\widetilde{V}, \pi^{*} N L-m E\right) \lambda_{i_{0}, v} \circ \pi+O_{v}(1)$
$\geq\left(\beta_{L, Y_{i_{0}}}-\delta\right) N h^{0}(V, N L) \lambda_{Y_{0}, v} \circ \pi+O_{v}(1)$.
The functorial property of Weil functions implies $\lambda_{\left(\pi^{*} s_{j}\right), v}=\lambda_{\left(s_{j}\right), v} \circ \pi+O_{v}(1)$. Hence, the above inequality, together with (3-6), implies that

$$
\begin{align*}
& \sum_{i=1}^{q} \lambda_{Y_{i}, v}(x) \\
& \quad \leq \frac{\ell}{N \cdot h^{0}(V, N L)\left(\min _{1 \leq i \leq q}\left\{\beta_{L, Y_{i}}\right\}-\delta\right)} \max _{J}\left\{\sum_{j \in J} \lambda_{\left(s_{j}\right), v}(x)\right\}+O_{v}(1) \tag{3-9}
\end{align*}
$$

where $J$ is a subset containing $M$ linearly independent sections taken among the collection of sections $\left\{s_{j}\left(i_{0}, v\right) \mid 1 \leq i_{0} \leq q, v \in S\right\}$ coming from the claim (3-6). It then follows from Theorem 3.5 and a suitable choice of $\delta$ that for a given $\epsilon>0$ there exists a proper algebraic subset $Z$ of $V$ defined over $k$ such that

$$
\begin{equation*}
\sum_{v \in S} \sum_{i=1}^{q} \lambda_{Y_{i}, v}(x) \leq\left(\ell \cdot \max _{1 \leq i \leq q}\left\{\beta_{L, Y_{i}}^{-1}\right\}+\epsilon\right) h_{L}(x), \tag{3-10}
\end{equation*}
$$

for all $x \in V(k) \backslash Z(k)$.
Proof of Corollary 1.5. Let $v$ be a place of $k$. The main point of the proof is to reformulate the distance function $d_{v}(\cdot, \cdot)$ defined on $V(\bar{k})$ [McKinnon and Roth 2015, Section 2] into a product of several distance functions on $V(K)$, where $K$ is a finite extension of $k$. Following the construction in [McKinnon and Roth 2015,

Section 2], we fix an extension of $v$ to $\bar{k}$. The place defines an absolute value $\|\cdot\|_{v}$ on $\bar{k}$. If $K \subset \bar{k}$ is a finite extension of $k$, then $d_{v}(\cdot, \cdot)_{K}=d_{v}(\cdot, \cdot)_{k}^{\left[K_{v}: k_{v}\right]}$. Here $d_{v}(\cdot, \cdot)_{K}$ refers to the distance function defined by using the same embedding and normalizing with respect to $K$ and $d_{v}(\cdot, \cdot)_{k}$ the distance function normalized with respect to $k$ (see [McKinnon and Roth 2015, Proposition 2.1(b)]). Assume that $V \subset \mathbb{P}^{N}$ (given by the canonical map associated to the ample line bundle $L$ ). For a given fixed point $y=\left[y_{0}: \cdots: y_{N}\right] \in V(\bar{k})$, let $K$ be the Galois closure of $k\left(y_{0}, \ldots, y_{N}\right)$ over $k$. For each $v \in M_{k}$, the inclusion map $\left.\left(i_{v}\right)\right|_{K}: K \rightarrow \bar{k}_{v}$ induces a place $w_{0}:=v$ of $K$ over $v$, and other places $w$ of $K$ over $v$ are conjugates by elements $\sigma_{w} \in \operatorname{Gal}(K / k)$ such that $\left\|\sigma_{w}(a)\right\|_{w}=\|a\|_{v}$ for all $a \in K$. Then, for $x, y \in K$,

$$
\begin{aligned}
\prod_{w \in M_{K}, w \mid v} d_{w}\left(\sigma_{w}(x), \sigma_{w}(y)\right)_{K} & =\prod_{w \in M_{K}, w \mid v} d_{v}(x, y)_{K} \\
& =\prod_{w \in M_{K}, w \mid v} d_{v}(x, y)_{k}^{\left[K_{v}: k_{v}\right]} \\
& =[K: k] d_{v}(x, y)_{k},
\end{aligned}
$$

i.e.,

$$
\begin{equation*}
d_{v}(x, y)_{k}=\prod_{w \in M_{K}, w \mid v} d_{w}\left(\sigma_{w}(x), \sigma_{w}(y)\right)_{K}^{1 /[K: k]}, \quad \text { for } x, y \in K . \tag{3-11}
\end{equation*}
$$

To compute $\alpha_{y}(L)$, we consider any sequence $\left\{x_{i}\right\} \subseteq V(k)$ of distinct points with $d_{v}\left(y, x_{i}\right) \rightarrow 0$. By (3-11), we have $d_{v}\left(y, x_{i}\right)_{k}=\prod_{w \in M_{K}, w \mid v} d_{w}\left(\sigma_{w}(y), x_{i}\right)_{K}^{1 /[K: k]}$. (Here we extend $\sigma_{w} \in \operatorname{Gal}(K / k)$ to the map from $V(K)$ to $V(K)$ by acting on the coordinates of the points.) The distance function $d_{w}(y, x)$ in [McKinnon and Roth 2015] is constructed by choosing an embedding $\phi_{L}: V \rightarrow \mathbb{P}^{N}$ into a projective space via the sections of $L$ and measuring the distance in the embedded space. For a fixed $y$ we denote $-\log d_{w}(y, \cdot)$ by $\lambda_{\phi(y), w}$, which is a local Weil function on the embedded space. We note that this fact can also be proved by a slight modification of Lemma 2.6 in [McKinnon and Roth 2015]. By the functoriality of Weil functions of closed subschemes [Silverman 1987, Theorem 2.1(h)] we have $-\log d_{w}\left(\sigma_{w}(y), x\right)=$ $\lambda_{\sigma_{w}(y), w}(x)+O(1)$. On the other hand, it is clear from the definition that $\beta_{L, y}=$ $\beta_{L, \sigma_{w}(y)}$ for very $\sigma_{w} \in \operatorname{Gal}(K / k)$. The Main Theorem then implies that for any $\epsilon>0$

$$
\begin{equation*}
\log d_{v}\left(y, x_{i}\right)=\frac{1}{[K: k]} \sum_{w \in M_{K}, w \mid v} \log d_{w}\left(y, x_{i}\right) \geq-\left(\left\{\beta_{L, y}^{-1}\right\}+\epsilon\right) h_{L}\left(x_{i}\right) \tag{3-12}
\end{equation*}
$$

holds for all $x_{i}$ outside a proper Zariski-closed subset $Z$ of $V(K)$ (note that, in this case, $\ell=1$ ). We note that $Z$ is indeed defined over $k$ since all the $x_{i}$ are in $k$. In conclusion, we have shown that for all sequences $\left\{x_{i}\right\} \subseteq V(k)$ of distinct points with $d_{v}\left(y, x_{i}\right) \rightarrow 0$, if $\alpha_{y}\left(\left\{x_{i}\right\}, L\right)<\beta_{L, y}$, then all but finitely many of the points of $\left\{x_{i}\right\}$ lie in $Z$. If (a) holds, then we are done. Therefore we assume that $\alpha_{y}(L)>\beta_{L, y}$.

Then the previous conclusion shows, in this case, that $\alpha_{y}(L)=\alpha_{y, Z}\left(\left.L\right|_{Z}\right)$. To see $Z$ is irreducible over $\bar{k}$, we first use Proposition 2.14(f) in [McKinnon and Roth 2015] to reduce $Z$ to one of the irreducible components of $Z$ over $k$, say $Y$ such that $\alpha_{y, Z}\left(\left.L\right|_{Z}\right)=\alpha_{y, Y}\left(\left.L\right|_{Y}\right)$. Without loss of generality we can assume that $Z=Y$, i.e., $Z$ itself is irreducible over $k$. We then apply Lemma 2.17 in [McKinnon and Roth 2015] to conclude that $Z$ is indeed geometrically irreducible, i.e., $Z$ is irreducible over $\bar{k}$.

## 4. The complex case

In this section, we consider the analogous result of our Main Theorem in Nevanlinna theory. Let $V$ be a complex projective variety. We use the standard notation in Nevanlinna theory (see, for example, [Ru 2016]). Note that the Weil function for divisors has been defined, so the Weil function $\lambda_{Y}$ for a subscheme $Y \subset V$ can also be defined using Lemma 3.1, similar to Definition 3.2. We define, for a holomorphic map $f: \mathbb{C} \rightarrow V$ with $f(\mathbb{C}) \not \subset Y$, the proximity function

$$
m_{f}(r, Y)=\int_{0}^{2 \pi} \lambda_{Y}\left(f\left(r e^{i \theta}\right)\right) \frac{d \theta}{2 \pi}
$$

We note that all the properties used above about the Weil functions in the arithmetic case hold for the complex case (see, for example, [Ru 2016; Ru and Vojta 2016]).

Theorem 4.1. Let $V$ be a complex projective variety and $Y_{1}, \ldots, Y_{q}$ be closed subschemes of $V$ in $\ell$-subgeneral position. Let L be a big line bundle. Let $f: \mathbb{C} \rightarrow V$ be a holomorphic map with Zariski dense image. Then for any $\epsilon>0$

$$
\begin{equation*}
\sum_{i=1}^{q} m_{f}\left(r, Y_{i}\right) \leq \ell\left(\max _{1 \leq i \leq q}\left\{\beta_{L, Y_{i}}^{-1}\right\}+\epsilon\right) T_{f, L}(r) \|, \tag{4-1}
\end{equation*}
$$

where $\|$ means that the inequality holds for all $r \in(0,+\infty)$ outside a set of finite Lebesgue measure.

To prove the theorem, we need the following result.
Theorem 4.2 [ Ru and Vojta 2016, Theorem 2.8]. Let $X$ be a complex projective $v a$ riety, $D$ be a Cartier divisor on $X, W$ be a nonzero linear subspace of $H^{0}(X, \mathcal{O}(D))$, and $s_{1}, \ldots, s_{q}$ be nonzero elements of $W$. Let $f: \mathbb{C} \rightarrow X$ be a holomorphic map with Zariski-dense image. Then

$$
\int_{0}^{2 \pi} \max _{J} \sum_{j \in J} \lambda_{\left(s_{j}\right)}\left(f\left(r e^{i \theta}\right)\right) \frac{d \theta}{2 \pi} \leq(\operatorname{dim} W) T_{f, D}(r)+O\left(\log ^{+} T_{f, D}(r)\right)+o(\log r) \|,
$$

where the set $J$ ranges over all subsets of $\{1, \ldots, q\}$ such that the sections $\left(s_{j}\right)_{j \in J}$ are linearly independent.

Proof of Theorem 4.1. Similar to the proof of the Main Theorem, let $\delta>0$ be a sufficiently small number. We choose $N$ large enough that, for $i=1, \ldots, q$,

$$
\sum_{m=1}^{\infty} h^{0}\left(\tilde{V}_{i}, N \pi_{i}^{*} L-m E_{i}\right) \geq\left(\beta_{L, Y_{i}}-\delta\right) N h^{0}(V, N L)
$$

Let $x \in V$. Since $Y_{i}, 1 \leq i \leq q$, are in $\ell$-subgeneral position, similar to Lemma 3.4, we have

$$
\begin{equation*}
\sum_{i=1}^{q} \lambda_{Y_{i}}(x) \leq \ell \lambda_{Y_{i_{0}}}(x)+O(1) \tag{4-2}
\end{equation*}
$$

for some $i_{0}$ with $1 \leq i_{0} \leq q$, where $i_{0}$ depends on the point $x$, but $O(1)$ is independent of $x$.

Let $\pi: \widetilde{V} \rightarrow V$ be the blow-up at $Y_{i_{0}}$ and $E=\pi^{-1}\left(Y_{i_{0}}\right)$ be the exceptional divisor of $\pi$. We consider the filtration of $H^{0}\left(\widetilde{V}, \pi^{*} N L\right)$ defined in (3-7). By identifying $H^{0}(V, N L)$ with $H^{0}\left(\widetilde{V}, \pi^{*} N L\right)$ as vector spaces, we can choose regular sections $s_{1}, \ldots, s_{M} \in H^{0}(V, N L)$, where $M=h^{0}(V, N L)$, successively so that their pullbacks $\pi^{*} s_{1}, \ldots, \pi^{*} s_{M} \in H^{0}\left(\tilde{V}, \pi^{*} N L\right)$ form a basis associated to this filtration. Then, in the same way as deriving (3-9), we can get

$$
\sum_{i=1}^{q} \lambda_{Y_{i}}(x) \leq \frac{\ell}{N \cdot h^{0}(V, N L)\left(\min _{1 \leq i \leq q}\left\{\beta_{L, Y_{i}}\right\}-\delta\right)} \sum_{j=1}^{q} \lambda_{\left(s_{j}\right)}(x)+O(1)
$$

Note that the basis $\left\{s_{1}, \ldots, s_{M}\right\}$ depends only on $i_{0}$, so the number of such choices is finite, since $i_{0} \in\{1, \ldots, q\}$, while $x$ varies in (4-2). We denote the set of bases as $J_{1}, \ldots, J_{T}$. Thus we get, for every $x \in V$,

$$
\sum_{i=1}^{q} \lambda_{Y_{i}}(x) \leq \frac{\ell}{N \cdot h^{0}(V, N L)\left(\min _{1 \leq i \leq q}\left\{\beta_{L, Y_{i}}\right\}-\delta\right)} \max _{1 \leq t \leq T} \sum_{j \in J_{t}} \lambda_{\left(s_{j}\right)}(x)+O(1)
$$

By taking $x=f\left(r e^{i \theta}\right)$ and then integrating, it then follows from Theorem 4.2 and a suitable choice of $\delta$ that, for the given $\epsilon>0$,

$$
\sum_{i=1}^{q} \int_{0}^{2 \pi} \lambda_{Y_{i}}\left(f\left(r e^{i \theta}\right)\right) \frac{d \theta}{2 \pi} \leq \ell\left(\max _{1 \leq i \leq q}\left\{\beta_{\mathcal{L}, Y_{i}}^{-1}\right\}+\epsilon\right) T_{f, \mathcal{L}}(r) \| .
$$

This finishes the proof.
Theorem 4.1, together with Lemma 2.2, implies the following corollary.
Corollary 4.3. Let $V$ be a complex projective variety of dimension $n$ and $a_{1}, \ldots, a_{q}$ be distinct points on $V$. Let $L$ be an ample line bundle. Let $f: \mathbb{C} \rightarrow V$ be a
holomorphic map with Zariski dense image. Then for any $\epsilon>0$,

$$
\sum_{i=1}^{q} m_{f}\left(r, a_{i}\right) \leq\left(\frac{n+1}{n} \max _{1 \leq i \leq q}\left\{\epsilon_{a_{i}}^{-1}(L)\right\}+\epsilon\right) T_{f, L}(r) \|
$$

where $\epsilon_{x}(L)$ is the Seshadri constant of $L$ at the point $x \in V$.
In particular, if $V=\mathbb{P}^{n}$, then for any $\epsilon>0$,

$$
\sum_{i=1}^{q} m_{f}\left(r, a_{i}\right) \leq\left(\frac{n+1}{n}+\epsilon\right) T_{f, L}(r) \|
$$

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# Variation of <br> anticyclotomic Iwasawa invariants in Hida families 

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Building on the construction of big Heegner points in the quaternionic setting by Longo and Vigni, and their relation to special values of Rankin-Selberg $L$ functions established by Castella and Longo, we obtain anticyclotomic analogues of the results of Emerton, Pollack and Weston on the variation of Iwasawa invariants in Hida families. In particular, combined with the known cases of the anticyclotomic Iwasawa main conjecture in weight 2 , our results yield a proof of the main conjecture for $p$-ordinary newforms of higher weights and trivial nebentypus.
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## Introduction

In the remarkable paper [Emerton et al. 2006], Emerton, Pollack and Weston obtained striking results on the behavior of the cyclotomic Iwasawa invariants attached to $p$-ordinary modular forms as they vary in Hida families. In particular, combined with Greenberg's conjecture on the vanishing of the $\mu$-invariant, their main result reduces the proof of the main conjecture to the weight two case. In this paper, we develop analogous results for newforms base-changed to imaginary quadratic fields in the definite anticyclotomic setting. In particular, combined with

[^17]Vatsal's result [2003] on the vanishing of the anticyclotomic $\mu$-invariant, and the known cases of the anticyclotomic main conjecture in weight 2 (thanks to the works of Bertolini and Darmon [2005], Pollack and Weston [2011], and Skinner and Urban [2014]), our results yield a proof of Iwasawa's main conjecture for $p$-ordinary modular forms of higher weights $k \geqslant 2$ and trivial nebentypus in the anticyclotomic setting.

Let us begin by recalling the setup of [Emerton et al. 2006], but adapted to the context at hand. Let

$$
\bar{\rho}: G_{\mathbb{Q}}:=\operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q}) \rightarrow \mathrm{GL}_{2}(\mathbb{F})
$$

be a continuous Galois representation defined over a finite field $\mathbb{F}$ of characteristic $p>3$, and assume that $\bar{\rho}$ is odd and irreducible. After the proof of Serre's conjecture [Khare and Wintenberger 2009], we know that $\bar{\rho}$ is modular, meaning that $\bar{\rho}$ is isomorphic to the $\bmod p$ Galois representation $\bar{\rho}_{f_{0}}$ associated to an elliptic newform $f_{0}$. Throughout this paper, it will be assumed that $\bar{\rho} \simeq \bar{\rho}_{f_{0}}$ for some newform $f_{0}$ of weight 2 and trivial nebentypus.

Let $N(\bar{\rho})$ be the tame conductor of $\bar{\rho}$, and let $K / \mathbb{Q}$ be an imaginary quadratic field of discriminant prime $-D_{K}<0$ to $p N(\bar{\rho})$. The field $K$ then determines a decomposition

$$
N(\bar{\rho})=N(\bar{\rho})^{+} \cdot N(\bar{\rho})^{-}
$$

with $N(\bar{\rho})^{+}\left(\right.$resp. $\left.N(\bar{\rho})^{-}\right)$only divisible by primes which are split (resp. inert) in $K$. We similarly define the decomposition $M=M^{+} \cdot M^{-}$for any positive integer $M$ prime to $D_{K}$.

As in [Pollack and Weston 2011], we consider the following conditions on a pair ( $\bar{\rho}, N^{-}$), where $N^{-}$is a fixed square-free product of an odd number of primes inert in $K$ :

Assumption (CR). (1) $\bar{\rho}$ is irreducible;
(2) $N(\bar{\rho})^{-} \mid N^{-}$;
(3) $\bar{\rho}$ is ramified at every prime $\ell \mid N^{-}$such that $\ell \equiv \pm 1(\bmod p)$.

Let $\mathcal{H}(\bar{\rho})$ be the set of all $p$-ordinary and $p$-stabilized newforms with $\bmod p$ Galois representation isomorphic to $\bar{\rho}$, and let $\Gamma:=\operatorname{Gal}\left(K_{\infty} / K\right)$ denote the Galois group of the anticyclotomic $\mathbb{Z}_{p}$-extension of $K$. Associated with each $f \in \mathcal{H}(\bar{\rho})$ of tame level $N_{f}$ with $N_{f}^{-}=N^{-}$, defined over say a finite extension $F / \mathbb{Q}_{p}$ with ring of integers $\mathcal{O}$, there is a $p$-adic $L$-function

$$
L_{p}(f / K) \in \mathcal{O} \llbracket \Gamma \rrbracket
$$

constructed by Bertolini and Darmon [1996] in weight two, and by Chida and Hsieh [2016] for higher weights. The $p$-adic $L$-function $L_{p}(f / K)$ is characterized, as
$\chi$ runs over the $p$-adic characters of $\Gamma$ corresponding to certain algebraic Hecke characters of $K$, by an interpolation property of the form

$$
\chi\left(L_{p}(f / K)\right)=C_{p}(f, \chi) \cdot E_{p}(f, \chi) \cdot \frac{L(f / K, \chi, k / 2)}{\Omega_{f, N^{-}}},
$$

where $C_{p}(f, \chi)$ is an explicit nonzero constant, $E_{p}(f, \chi)$ is a $p$-adic multiplier, and $\Omega_{f, N^{-}}$is a complex period making the above ratio algebraic. (Of course, implicit in all the above is a fixed choice of complex and $p$-adic embeddings $\mathbb{C} \stackrel{\iota_{\infty}}{\longleftrightarrow} \overline{\mathbb{Q}} \stackrel{\iota_{p}}{\longleftrightarrow} \overline{\mathbb{Q}}_{p}$.)

The anticyclotomic Iwasawa main conjecture gives an arithmetic interpretation of $L_{p}(f / K)$. More precisely, let

$$
\rho_{f}: G_{\mathbb{Q}} \rightarrow \operatorname{Aut}_{F}\left(V_{f}\right) \simeq \operatorname{GL}_{2}(F)
$$

be a self-dual twist of the $p$-adic Galois representation associated to $f$, fix an $\mathcal{O}$-stable lattice $T_{f} \subseteq V_{f}$, and set $A_{f}:=V_{f} / T_{f}$. Let $D_{p} \subseteq G_{\mathbb{Q}}$ be the decomposition group corresponding to our fixed embedding $\iota_{p}$, and let $\varepsilon_{\mathrm{cyc}}$ be the $p$-adic cyclotomic character. Since $f$ is $p$-ordinary, there is a unique one-dimensional $D_{p}$-invariant subspace $F_{p}^{+} V_{f} \subseteq V_{f}$ where the inertia group at $p$ acts via $\varepsilon_{\text {cyc }}^{k / 2} \psi$, with $\psi$ a finite order character. Let $F_{p}^{+} A_{f}$ be the image of $F_{p}^{+} V_{f}$ in $A_{f}$ and set $F_{p}^{-} A_{f}:=A_{f} / F_{p}^{+} A_{f}$. Following the terminology in [Pollack and Weston 2011], the minimal Selmer group of $f$ is defined by
$\operatorname{Sel}\left(K_{\infty}, f\right):=\operatorname{ker}\left\{H^{1}\left(K_{\infty}, A_{f}\right) \rightarrow \prod_{w \nmid p} H^{1}\left(K_{\infty, w}, A_{f}\right) \times \prod_{w \mid p} H^{1}\left(K_{\infty, w}, F_{p}^{-} A_{f}\right)\right\}$,
where $w$ runs over the places of $K_{\infty}$. By standard arguments (see [Greenberg 1989], for example), one knows that the Pontryagin dual of $\operatorname{Sel}\left(K_{\infty}, f\right)$ is finitely generated over the anticyclotomic Iwasawa algebra $\Lambda:=\mathcal{O} \llbracket \Gamma \rrbracket$. The anticyclotomic main conjecture is then the following:
Conjecture 1. The Pontryagin dual $\operatorname{Sel}\left(K_{\infty}, f\right)^{\vee}$ is $\Lambda$-torsion, and

$$
C h_{\Lambda}\left(\operatorname{Sel}\left(K_{\infty}, f\right)^{\vee}\right)=\left(L_{p}(f / K)\right)
$$

For newforms $f$ of weight 2 corresponding to elliptic curves $E / \mathbb{Q}$ with ordinary reduction at $p$, and under rather stringent assumptions on $\bar{\rho}_{f}$ which were later relaxed by Pollack and Weston [2011], one of the divisibilities predicted by Conjecture 1 was obtained by Bertolini and Darmon [2005] using Heegner points and Kolyvagin's method of Euler systems. More recently, after the work of Chida and Hsieh [2015] the divisibility

$$
C h_{\Lambda}\left(\operatorname{Sel}\left(K_{\infty}, f\right)^{\vee}\right) \supseteq\left(L_{p}(f / K)\right)
$$

is known for newforms $f$ of weight $k \leqslant p-2$ and trivial nebentypus, provided the pair $\left(\bar{\rho}_{f}, N_{f}^{-}\right)$satisfies a mild strengthening of Hypotheses (CR). This restriction
to small weights comes from the use of Ihara's lemma [Diamond and Taylor 1994], and it seems difficult to directly extend their arguments in [Chida and Hsieh 2015] to higher weights. Instead, as we shall explain in the following paragraphs, in this paper we will complete the proof of Conjecture 1 to all weights $k \equiv 2(\bmod p-1)$ by a different approach, using Howard's big Heegner points in Hida families [Howard 2007], as extended by Longo and Vigni [2011] to quaternionic Shimura curves.

Associated with every $f \in \mathcal{H}(\bar{\rho})$ there are anticyclotomic Iwasawa invariants $\mu^{\text {an }}\left(K_{\infty}, f\right), \lambda^{\text {an }}\left(K_{\infty}, f\right), \mu^{\text {alg }}\left(K_{\infty}, f\right)$, and $\lambda^{\text {alg }}\left(K_{\infty}, f\right)$. The analytic (resp. algebraic) $\lambda$-invariants are the number of zeros of $L_{p}(f / K)$ (resp. of a generator of the characteristic ideal of $\left.\operatorname{Sel}\left(K_{\infty}, f\right)^{\vee}\right)$, while the $\mu$-invariants are defined as the exponent of the highest power of $\varpi$ (with $\varpi \in \mathcal{O}$ any uniformizer) dividing the same objects. Our main results on the variation of these invariants are summarized in the following. (Recall that we assume $\bar{\rho} \simeq \bar{\rho}_{f_{0}}$ for some newform $f_{0}$ of weight 2 and trivial nebentypus.)

Theorem 2. Assume in addition that:

- $\bar{\rho}$ is irreducible;
- $\bar{\rho}$ is p-ordinary, "nonanomalous" and p-distinguished:

$$
\left.\bar{\rho}\right|_{D_{p}} \simeq\left(\begin{array}{ll}
\bar{\varepsilon} & * \\
0 & \bar{\delta}
\end{array}\right),
$$

with $\bar{\varepsilon}, \bar{\delta}: D_{p} \rightarrow \mathbb{F}^{\times}$characters such that $\bar{\delta}$ is unramified, $\bar{\delta}\left(\operatorname{Frob}_{p}\right) \neq \pm 1$ and $\bar{\delta} \neq \bar{\varepsilon} ;$

- $N(\bar{\rho})^{-}$is the square-free product of an odd number of primes.

Let $\mathcal{H}^{-}(\bar{\rho}):=\mathcal{H}^{N(\bar{\rho})^{-}}(\bar{\rho})$ consist of all newforms $f \in \mathcal{H}(\bar{\rho})$ with $N_{f}^{-}=N(\bar{\rho})^{-}$, and fix $* \in\{\mathrm{alg}, \mathrm{an}\}$. Then the following hold:
(1) For all $f \in \mathcal{H}^{-}(\bar{\rho})$, we have

$$
\mu^{*}\left(K_{\infty}, f\right)=0
$$

(2) Let $f_{1}, f_{2} \in \mathcal{H}^{-}(\bar{\rho})$ lie on the branches $\mathbb{T}\left(\mathfrak{a}_{1}\right), \mathbb{T}\left(\mathfrak{a}_{2}\right)$ (defined in $\left.\S 1 D\right)$, respectively. Then

$$
\lambda^{*}\left(K_{\infty}, f_{1}\right)-\lambda^{*}\left(K_{\infty}, f_{2}\right)=\sum_{\ell \mid N_{f_{1}}^{+} N_{f_{2}}^{+}} e_{\ell}\left(\mathfrak{a}_{2}\right)-e_{\ell}\left(\mathfrak{a}_{1}\right)
$$

where the sum is over the split primes in $K$ which divide the tame level of $f_{1}$ or $f_{2}$, and $e_{\ell}\left(\mathfrak{a}_{j}\right)$ is an explicit nonnegative invariant of the branch $\mathbb{T}\left(\mathfrak{a}_{j}\right)$ and the prime $\ell$.

Provided that $p$ splits in $K$, and under the same hypotheses on $\bar{\rho}$ as in Theorem 2, the work of Skinner and Urban [2014] establishes one of the divisibilities in their "three-variable" Iwasawa main conjecture. Combining their work with our Theorem 2, and making use of the aforementioned results of Bertolini and Darmon [2005] and Pollack and Weston [2011] in weight 2, we obtain many new cases of Conjecture 1 (cf., Corollary 5.5):

Corollary 3. Suppose that $\bar{\rho}$ is as in Theorem 2 and that $p$ splits in K. Then the anticyclotomic Iwasawa main conjecture holds for every $f \in \mathcal{H}^{-}(\bar{\rho})$ of weight $k \equiv 2(\bmod p-1)$ and trivial nebentypus.

Let us briefly explain the new ingredients in the proof of Theorem 2. As it will be clear to the reader, the results contained in Theorem 2 are anticyclotomic analogues of the results of Emerton, Pollack and Weston [Emerton et al. 2006] in the cyclotomic setting. In fact, on the algebraic side the arguments of loc.cit. carry over almost verbatim, and our main innovations in this paper are in the development of anticyclotomic analogues of their results on the analytic side. Indeed, the analytic results of [Emerton et al. 2006] are based on the study of certain two-variable $p$-adic $L$-functions à la Mazur and Kitagawa, whose construction relies on the theory of modular symbols on classical modular curves. In contrast, we need to work on a family of Shimura curves associated with definite quaternion algebras, for which cusps are not available. In the cyclotomic case, modular symbols are useful in two ways: They yield a concrete realization of the degree-one compactly supported cohomology of open modular curves, and provide a powerful tool for studying the arithmetic properties of critical values of the $L$-functions attached to modular forms. Our basic observation is that in the present anticyclotomic setting, Heegner points on definite Shimura curves provide a similarly convenient way of describing the central critical values of the Rankin $L$-series $L(f / K, \chi, s)$.

Also fundamental for the method of [Emerton et al. 2006] is the possibility to "deform" modular symbols in Hida families. In our anticyclotomic context, the construction of big Heegner points in Hida families was obtained in the work [Longo and Vigni 2011] of one of us in collaboration with Vigni, while the relation between these points and Rankin-Selberg $L$-values was established in the work [Castella and Longo 2016] by two of us. With these key results at hand, and working over appropriate quotients of the Hecke algebras considered in [Emerton et al. 2006] via the Jacquet-Langlands correspondence, we are then able to develop analogues of the arguments of loc. cit. in our setting, making use of the ramification hypotheses on $\bar{\rho}$ to ensure a multiplicity one property of certain Hecke modules, similarly as in the works of Pollack and Weston [2011] and one of us [Kim 2017].

We conclude this introduction with an overview of the contents of the paper. In the next section, we briefly recall the Hida theory that we need, following the
exposition in [Emerton et al. 2006, §1] for the most part. In Section 2, we describe a key extension of the construction of big Heegner points of [Longo and Vigni 2011] to "imprimitive" branches of the Hida family. In Section 3, we construct twovariable $p$-adic $L$-functions attached to a Hida family and to each of its irreducible components (or branches), and prove Theorem 3.10 relating the two. This theorem is the key technical result of this paper, and the analytic part of Theorem 2 follows easily from this. In Section 4, we deduce the algebraic part of Theorem 2 using the residual Selmer groups studied in [Pollack and Weston 2011, §3.2]. Finally, in Section 5 we give the applications of our results to the anticyclotomic Iwasawa main conjecture.

## 1. Hida theory

Throughout this section, we fix a positive integer $N$ admitting a factorization

$$
N=N^{+} N^{-}
$$

with $\left(N^{+}, N^{-}\right)=1$ and $N^{-}$equal to the square-free product of an odd number of primes. We also fix a prime $p \nmid 6 N$.

1A. Hecke algebras. For each integer $k \geqslant 2$, denote by $\mathfrak{h}_{N, r, k}$ the $\mathbb{Z}_{p}$-algebra generated by the Hecke operators $T_{\ell}$ for $\ell \nmid N p$, the operators $U_{\ell}$ for $\ell \mid N p$, and the diamond operators $\langle a\rangle$ for $a \in\left(\mathbb{Z} / p^{r} \mathbb{Z}\right)^{\times}$, acting on the space $S_{k}\left(\Gamma_{0,1}\left(N, p^{r}\right), \overline{\mathbb{Q}}_{p}\right)$ of cusp forms of weight $k$ on $\Gamma_{0,1}\left(N, p^{r}\right):=\Gamma_{0}(N) \cap \Gamma_{1}\left(p^{r}\right)$. For $k=2$, we abbreviate $\mathfrak{h}_{N, r}:=\mathfrak{h}_{N, r, 2}$.

Let $e^{\text {ord }}:=\lim _{n \rightarrow \infty} U_{p}^{n!}$ be Hida's ordinary projector, and define

$$
\mathfrak{h}_{N, r, k}^{\text {ord }}:=e^{\text {ord }} \mathfrak{h}_{N, r, k}, \quad \mathfrak{h}_{N, r}^{\text {ord }}:=e^{\text {ord }} \mathfrak{h}_{N, r}, \quad \mathfrak{h}_{N}^{\text {ord }}:={\underset{\longleftarrow}{\lim }}_{\mathfrak{h}_{N, r}^{\text {ord }},}
$$

where the limit is over the projections induced by the natural restriction maps.
Denote by $\mathbb{T}_{N, r, k}^{N^{-}}$the quotient of $\mathfrak{h}_{N, r, k}^{\text {ord }}$ acting faithfully on the subspace of $e^{\text {ord }} S_{k}\left(\Gamma_{0,1}\left(N, p^{r}\right), \overline{\mathbb{Q}}_{p}\right)$ consisting of forms which are new at all primes dividing $N^{-}$. Set $\mathbb{T}_{N, r}^{N^{-}}:=\mathbb{T}_{N, r, 2}^{N^{-}}$and define

$$
\mathbb{T}_{N}^{N^{-}}:={\underset{\overleftarrow{~}}{r}}^{\mathbb{T}_{N, r}} \mathbb{U}^{N^{-}}
$$

Each of these Hecke algebras is equipped with natural $\mathbb{Z}_{p} \llbracket \mathbb{Z}_{p}^{\times} \rrbracket$-algebra structures via the diamond operators, and by a well-known result of Hida, $\mathfrak{h}_{N}^{\text {ord }}$ is finite and flat over $\mathbb{Z}_{p} \llbracket 1+p \mathbb{Z}_{p} \rrbracket$.

1B. Galois representations on Hecke algebras. For each positive integer $M \mid N$ we may consider the new quotient $\mathbb{T}_{M}^{\text {new }}$ of $\mathfrak{h}_{M}^{\text {ord }}$, and the Galois representation

$$
\rho_{M}: G_{\mathbb{Q}} \rightarrow \mathrm{GL}_{2}\left(\mathbb{T}_{M}^{\mathrm{new}} \otimes \mathcal{L}\right)
$$

described in [Emerton et al. 2006, Theorem 2.2.1], where $\mathcal{L}$ denotes the fraction field of $\mathbb{Z}_{p} \llbracket 1+p \mathbb{Z}_{p} \rrbracket$.

Let $\mathbb{T}_{N}^{\prime}$ be the $\mathbb{Z}_{p} \llbracket 1+p \mathbb{Z}_{p} \rrbracket$-subalgebra of $\mathbb{T}_{N}^{N^{-}}$generated by the image under the natural projection $\mathfrak{h}_{N}^{\text {ord }} \rightarrow \mathbb{T}_{N}^{N^{-}}$of the Hecke operators of level prime to $N$. As in [Emerton et al. 2006, Proposition 2.3.2], one can show that the canonical map

$$
\mathbb{T}_{N}^{\prime} \rightarrow \prod_{M} \mathbb{T}_{M}^{\mathrm{new}}
$$

where the product is over all integers $M \geqslant 1$ with $N^{-}|M| N$, becomes an isomorphism after tensoring with $\mathcal{L}$. Taking the product of the Galois representations $\rho_{M}$ we thus obtain

$$
\rho: G_{\mathbb{Q}} \rightarrow \mathrm{GL}_{2}\left(\mathbb{T}_{N}^{\prime} \otimes \mathcal{L}\right)
$$

For any maximal ideal $\mathfrak{m}$ of $\mathbb{T}_{N}^{\prime}$, let $\left(\mathbb{T}_{N}^{\prime}\right)_{\mathfrak{m}}$ denote the localization of $\mathbb{T}_{N}^{\prime}$ at $\mathfrak{m}$ and let

$$
\rho_{\mathfrak{m}}: G_{\mathbb{Q}} \rightarrow \mathrm{GL}_{2}\left(\left(\mathbb{T}_{N}^{\prime}\right)_{\mathfrak{m}} \otimes \mathcal{L}\right)
$$

be the resulting Galois representation. If the residual representation $\bar{\rho}_{\mathfrak{m}}$ is irreducible, then $\rho_{\mathfrak{m}}$ admits an integral model (still denoted in the same manner)

$$
\rho_{\mathfrak{m}}: G_{\mathbb{Q}} \rightarrow \mathrm{GL}_{2}\left(\left(\mathbb{T}_{N}^{\prime}\right)_{\mathfrak{m}}\right)
$$

which is unique up to isomorphism.
1C. Residual representations. Let $\bar{\rho}: G_{\mathbb{Q}} \rightarrow \mathrm{GL}_{2}(\mathbb{F})$ be an odd irreducible Galois representation defined over a finite field $\mathbb{F}$ of characteristic $p>3$. As in the introduction, we assume that $\bar{\rho} \simeq \bar{\rho}_{f_{0}}$ for some newform $f_{0}$ of weight 2 , level $N$, and trivial nebentypus. Consider the following three conditions we may impose on the pair $\left(\bar{\rho}, N^{-}\right)$:

Assumption (SU). (1) $\bar{\rho}$ is $p$-ordinary: the restriction of $\bar{\rho}$ to a decomposition group $D_{p} \subseteq G_{\mathbb{Q}}$ at $p$ has a one-dimensional unramified quotient over $\mathbb{F}$;
(2) $\bar{\rho}$ is $p$-distinguished: $\left.\bar{\rho}\right|_{D_{p}} \sim\left(\begin{array}{c}\bar{\varepsilon} * \\ 0 \\ \bar{\delta}\end{array}\right)$ with $\bar{\varepsilon} \neq \bar{\delta}$;
(3) $\bar{\rho}$ is ramified at every prime $\ell \mid N^{-}$.

Fix once and for all a representation $\bar{\rho}$ satisfying Assumption (SU), together with a $p$-stabilization of $\bar{\rho}$ in the sense of [Emerton et al. 2006, Definition 2.2.10]. Let $\bar{V}$ be the two-dimensional $\mathbb{F}$-vector space which affords $\bar{\rho}$, and for any finite set of primes $\Sigma$ that does not contain $p$ or any factor of $N^{-}$, define

$$
\begin{equation*}
N(\Sigma):=N(\bar{\rho}) \prod_{\ell \in \Sigma} \ell^{m_{\ell}} \tag{1}
\end{equation*}
$$

where $N(\bar{\rho})$ is the tame conductor of $\bar{\rho}$, and $m_{\ell}:=\operatorname{dim}_{\mathbb{F}} \bar{V}_{I_{\ell}}$.

Remark 1.1. By Assumption (SU) we have the divisibility $N^{-} \mid N(\bar{\rho})$; we will further assume that $\left(N^{-}, N(\bar{\rho}) / N^{-}\right)=1$.

Combining [Emerton et al. 2006, Theorem 2.4.1] and [Emerton et al. 2006, Proposition 2.4.2] with the fact that $\bar{\rho}$ is ramified at the primes dividing $N^{-}$, one can see that there exist unique maximal ideals $\mathfrak{n}$ and $\mathfrak{m}$ of $\mathbb{T}_{N(\Sigma)}^{N^{-}}$and $\mathbb{T}_{N(\Sigma)}^{\prime}$, respectively, such that

- $\mathfrak{n} \cap \mathbb{T}_{N(\Sigma)}^{\prime}=\mathfrak{m}$;
- $\left(\mathbb{T}_{N(\Sigma)}^{\prime}\right)_{\mathfrak{m}} \simeq\left(\mathbb{T}_{N(\Sigma)}^{N^{-}}\right)_{\mathfrak{n}}$ by the natural map on localizations;
- $\bar{\rho}_{\mathfrak{m}} \simeq \bar{\rho}$.

Define the ordinary Hecke algebra $\mathbb{T}_{\Sigma}$ attached to $\bar{\rho}$ and $\Sigma$ by

$$
\mathbb{T}_{\Sigma}:=\left(\mathbb{T}_{N(\Sigma)}^{\prime}\right)_{\mathfrak{m}}
$$

Thus $\mathbb{T}_{\Sigma}$ is a local factor of $\mathbb{T}_{N(\Sigma)}^{\prime}$, and we let

$$
\rho_{\Sigma}: G_{\mathbb{Q}} \rightarrow \mathrm{GL}_{2}\left(\mathbb{T}_{\Sigma}\right)
$$

denote the Galois representation deduced from $\rho_{\mathfrak{m}}$.
Adopting the terminology of [Emerton et al. 2006, §2.4], we shall refer to $\operatorname{Spec}\left(\mathbb{T}_{\Sigma}\right)$ as "the Hida family" $\mathcal{H}^{-}(\bar{\rho})$ attached to $\bar{\rho}$ (and our chosen $p$-stabilization) that is minimally ramified outside $\Sigma$.
Remark 1.2. Note that by Assumption (SU), all the $p$-stabilized newforms in $\mathcal{H}^{-}(\bar{\rho})$ have tame level divisible by $N^{-}$.

1D. Branches of the Hida family. If $\mathfrak{a}$ is a minimal prime of $\mathbb{T}_{\Sigma}$ (for a finite set of primes $\Sigma$ as above), we put $\mathbb{T}(\mathfrak{a}):=\mathbb{T}_{\Sigma} / \mathfrak{a}$ and let

$$
\rho(\mathfrak{a}): G_{\mathbb{Q}} \rightarrow \mathrm{GL}_{2}(\mathbb{T}(\mathfrak{a}))
$$

be the Galois representation induced by $\rho_{\Sigma}$. As in [Emerton et al. 2006, Proposition 2.5.2], one can show that there is a unique divisor $N(\mathfrak{a})$ of $N(\Sigma)$ and a unique minimal prime $\mathfrak{a}^{\prime} \subseteq \mathbb{T}_{N(\mathfrak{a})}^{\text {new }}$ above $\mathfrak{a}$ such that the diagram

commutes. We call $N(\mathfrak{a})$ the tame conductor $\mathfrak{o f} \mathfrak{a}$ and set

$$
\mathbb{T}(\mathfrak{a})^{\circ}:=\mathbb{T}_{N(\mathfrak{a})}^{\text {new }} / \mathfrak{a}^{\prime} .
$$

In particular, note that $N^{-} \mid N(\mathfrak{a})$ by construction, and that the natural map $\mathbb{T}(\mathfrak{a}) \rightarrow \mathbb{T}(\mathfrak{a})^{\circ}$ is an embedding of local domains.

1E. Arithmetic specializations. For any finite $\mathbb{Z}_{p} \llbracket 1+p \mathbb{Z}_{p} \rrbracket$-algebra $\mathbb{T}$, we say that a height one prime $\wp$ of $\mathbb{T}$ is an arithmetic prime of $\mathbb{T}$ if $\wp$ is the kernel of a $\mathbb{Z}_{p}$-algebra homomorphism $\mathbb{T} \rightarrow \overline{\mathbb{Q}}_{p}$ such that the composite map

$$
1+p \mathbb{Z}_{p} \rightarrow \mathbb{Z}_{p} \llbracket 1+p \mathbb{Z}_{p} \rrbracket^{\times} \rightarrow \mathbb{T}^{\times} \rightarrow \overline{\mathbb{Q}}_{p}^{\times}
$$

is given by $\gamma \mapsto \gamma^{k-2}$ on some open subgroup of $1+p \mathbb{Z}_{p}$, for some integer $k \geqslant 2$. We then say that $\wp$ has weight $k$.

Let $\mathfrak{a} \subseteq \mathbb{T}_{\Sigma}$ be a minimal prime as above. For each $n \geqslant 1$, let $\boldsymbol{a}_{n} \in \mathbb{T}(\mathfrak{a})^{\circ}$ be the image of $T_{n}$ under the natural projection $\mathfrak{h}_{N(\Sigma)}^{\text {ord }} \rightarrow \mathbb{T}(\mathfrak{a})^{\circ}$, and form the $q$-expansion

$$
\boldsymbol{f}(\mathfrak{a})=\sum_{n \geqslant 1} \boldsymbol{a}_{n} q^{n} \in \mathbb{T}(\mathfrak{a})^{\circ} \llbracket q \rrbracket .
$$

By [Hida 1986, Theorem 1.2], if $\wp$ is an arithmetic prime of $\mathbb{T}(\mathfrak{a})$ of weight $k$, then there is a unique height one prime $\wp^{\prime}$ of $\mathbb{T}(\mathfrak{a})^{\circ}$ such that

$$
\boldsymbol{f}_{\wp}(\mathfrak{a}):=\sum_{n \geqslant 1}\left(\boldsymbol{a}_{n} \bmod \wp^{\prime}\right) q^{n} \in \mathcal{O}_{\wp}^{\circ} \llbracket q \rrbracket,
$$

where $\mathcal{O}_{\wp}^{\circ}:=\mathbb{T}(\mathfrak{a})^{\circ} / \wp^{\prime}$, is the $q$-expansion of a $p$-ordinary eigenform $f_{\wp}$ of weight $k$ and tame level $N(\mathfrak{a})$ (see [Emerton et al. 2006, Proposition 2.5.6]).

## 2. Big Heegner points

As in Section 1, we fix an integer $N \geqslant 1$ admitting a factorization $N=N^{+} N^{-}$ with $\left(N^{+}, N^{-}\right)=1$ and $N^{-}$equal to the square-free product of an odd number of primes, and fix a prime $p \nmid 6 N$. Also, we let $K / \mathbb{Q}$ be an imaginary quadratic field of discriminant $-D_{K}<0$ prime to $N p$ and such that every prime factor of $N^{+}$(resp. $N^{-}$) splits (resp. is inert) in $K$.

In this section we describe a mild extension of the construction in [Longo and Vigni 2011] (following [Howard 2007]) of big Heegner points attached to $K$. Indeed, using the results from the preceding section, we can extend the constructions of loc.cit. to branches of the Hida family which are not necessarily primitive (in the sense of [Hida 1986, §1]). The availability of such an extension is fundamental for the purposes of this paper.

2A. Definite Shimura curves. Let $B$ be the definite quaternion algebra over $\mathbb{Q}$ of discriminant $N^{-}$. We fix once and for all an embedding of $\mathbb{Q}$-algebras $K \hookrightarrow B$, and use it to identity $K$ with a subalgebra of $B$. Denote by $z \mapsto \bar{z}$ the nontrivial automorphism of $K$, and choose a basis $\{1, j\}$ of $B$ over $K$ such that

- $j^{2}=\beta \in \mathbb{Q}^{\times}$with $\beta<0$;
- $j t=\bar{t} j$ for all $t \in K$;
- $\beta \in\left(\mathbb{Z}_{q}^{\times}\right)^{2}$ for $q \mid p N^{+}$, and $\beta \in \mathbb{Z}_{q}^{\times}$for $q \mid D_{K}$.

Fix a square-root $\delta_{K}=\sqrt{-D_{K}}$, and define $\boldsymbol{\theta} \in K$ by

$$
\boldsymbol{\theta}:=\frac{1}{2} D^{\prime}+\delta_{K}, \quad \text { where } D^{\prime}:= \begin{cases}D_{K} & \text { if } 2 \nmid D_{K} \\ \frac{1}{2} D_{K} & \text { if } 2 \mid D_{K}\end{cases}
$$

Note that $\mathcal{O}_{K}=\mathbb{Z}+\mathbb{Z} \boldsymbol{\theta}$, and for every prime $q \mid p N^{+}$, define $i_{q}: B_{q}:=B \otimes_{\mathbb{Q}} \mathbb{Q}_{q} \simeq$ $\mathrm{M}_{2}\left(\mathbb{Q}_{q}\right)$ by

$$
i_{q}(\boldsymbol{\theta})=\left(\begin{array}{cc}
\operatorname{Tr}(\boldsymbol{\theta}) & -\mathrm{Nm}(\boldsymbol{\theta}) \\
1 & 0
\end{array}\right), \quad i_{q}(j)=\sqrt{\beta}\left(\begin{array}{cc}
-1 & \operatorname{Tr}(\boldsymbol{\theta}) \\
0 & 1
\end{array}\right)
$$

where Tr and Nm are the reduced trace and reduced norm maps on $B$, respectively. On the other hand, for each prime $q \nmid N p$ we fix any isomorphism $i_{q}: B_{q} \simeq \mathrm{M}_{2}\left(\mathbb{Q}_{q}\right)$ with the property that $i_{q}\left(\mathcal{O}_{K} \otimes_{\mathbb{Z}} \mathbb{Z}_{q}\right) \subset \mathrm{M}_{2}\left(\mathbb{Z}_{q}\right)$.

For each $r \geqslant 0$, let $R_{N^{+}, r}$ be the Eichler order of $B$ of level $N^{+} p^{r}$ with respect to the above isomorphisms $\left\{i_{q}: B_{q} \simeq \mathrm{M}_{2}\left(\mathbb{Q}_{q}\right)\right\}_{q \nmid N^{-}}$, and let $U_{N^{+}, r}$ be the compact open subgroup of $\widehat{R}_{N^{+}, r}^{\times}$defined by

$$
U_{N^{+}, r}:=\left\{\left(x_{q}\right)_{q} \in \widehat{R}_{N^{+}, r}^{\times} \left\lvert\, i_{p}\left(x_{p}\right) \equiv\left(\begin{array}{cc}
1 & * \\
0 & *
\end{array}\right)\left(\bmod p^{r}\right)\right.\right\} .
$$

Consider the double coset spaces

$$
\begin{equation*}
\widetilde{X}_{N^{+}, r}=B^{\times} \backslash\left(\operatorname{Hom}_{\mathbb{Q}}(K, B) \times \widehat{B}^{\times}\right) / U_{N^{+}, r}, \tag{2}
\end{equation*}
$$

where $b \in B^{\times}$acts on $(\Psi, g) \in \operatorname{Hom}_{\mathbb{Q}}(K, B) \times \widehat{B}^{\times}$by

$$
b \cdot(\Psi, g)=\left(b \Psi b^{-1}, b g\right)
$$

and $U_{N^{+}, r}$ acts on $\widehat{B}^{\times}$by right multiplication. As is well known (see, e.g., [Longo and Vigni 2011, §2.1]), $\widetilde{X}_{N^{+}, r}$ may be naturally identified with the set of $K$-rational points of certain genus zero curves defined over $\mathbb{Q}$. Nonetheless, there is a nontrivial Galois action on $\widetilde{X}_{N^{+}, r}$ defined as follows: If $\sigma \in \operatorname{Gal}\left(K^{\mathrm{ab}} / K\right)$ and $P \in \widetilde{X}_{N^{+}, r}$ is the class of a pair $(\Psi, g)$, then

$$
P^{\sigma}:=[(\Psi, \widehat{\Psi}(a) g)],
$$

where $a \in K^{\times} \backslash \widehat{K}^{\times}$is chosen so that $\operatorname{rec}_{K}(a)=\sigma$. It will be convenient to extend this action to an action of $G_{K}:=\operatorname{Gal}(\overline{\mathbb{Q}} / K)$ in the obvious manner.

Finally, we note that $\widetilde{X}_{N^{+}, r}$ is also equipped with standard actions of $U_{p}$, Hecke operators $T_{\ell}$ for $\ell \nmid N p$, and diamond operators $\langle d\rangle$ for $d \in\left(\mathbb{Z} / p^{r} \mathbb{Z}\right)^{\times}$(see [Longo and Vigni 2011, §2.4], for example).

2B. Compatible systems of Heegner points. For each integer $c \geqslant 1$, let $\mathcal{O}_{c}=$ $\mathbb{Z}+c \mathcal{O}_{K}$ be the order of $K$ of conductor $c$.

Definition 2.1. We say that a point $P \in \widetilde{X}_{N^{+}, r}$ is a Heegner point of conductor $c$ if $P$ is the class of a pair $(\Psi, g)$ with

$$
\Psi\left(\mathcal{O}_{c}\right)=\Psi(K) \cap\left(B \cap g \widehat{R}_{N^{+}, r} g^{-1}\right)
$$

and

$$
\Psi_{p}\left(\left(\mathcal{O}_{c} \otimes \mathbb{Z}_{p}\right)^{\times} \cap\left(1+p^{r} \mathcal{O}_{K} \otimes \mathbb{Z}_{p}\right)^{\times}\right)=\Psi_{p}\left(\left(\mathcal{O}_{c} \otimes \mathbb{Z}_{p}\right)^{\times}\right) \cap g_{p} U_{N^{+}, r, p} g_{p}^{-1}
$$

where $U_{N^{+}, r, p}$ denotes the $p$-component of $U_{N^{+}, r}$.
Fix a decomposition $N^{+} \mathcal{O}_{K}=\mathfrak{N}^{+} \overline{\mathfrak{N}^{+}}$, and for each prime $q \neq p$ define

- $\varsigma_{q}=1$, if $q \nmid N^{+}$;
- $\varsigma_{q}=\delta_{K}^{-1}\left(\begin{array}{ll}\boldsymbol{\theta} & \overline{\boldsymbol{\theta}} \\ 1 & 1\end{array}\right) \in \mathrm{GL}_{2}\left(K_{\mathfrak{q}}\right)=\mathrm{GL}_{2}\left(\mathbb{Q}_{q}\right)$, if $q=\mathfrak{q} \overline{\mathfrak{q}}$ splits with $\mathfrak{q} \mid \mathfrak{N}^{+}$,
and for each $s \geqslant 0$, let
- $\varsigma_{p}^{(s)}=\left(\begin{array}{rr}\boldsymbol{\theta} & -1 \\ 1 & 0\end{array}\right)\left(\begin{array}{ll}p^{s} & 0 \\ 0 & 1\end{array}\right) \in \mathrm{GL}_{2}\left(K_{\mathfrak{p}}\right)=\mathrm{GL}_{2}\left(\mathbb{Q}_{p}\right)$, if $p=\mathfrak{p} \bar{p}$ splits in $K$;
- $\varsigma_{p}^{(s)}=\left(\begin{array}{rr}0 & 1 \\ -1 & 0\end{array}\right)\left(\begin{array}{ll}p^{s} & 0 \\ 0 & 1\end{array}\right)$, if $p$ is inert in $K$.

Set $\varsigma^{(s)}:=\varsigma_{p}^{(s)} \prod_{q \neq p} \varsigma_{q}$, viewed as an element in $\widehat{B}^{\times}$via the isomorphisms $\left\{i_{q}: B_{q} \simeq M_{2}\left(\mathbb{Q}_{q}\right)\right\}_{q \nmid N^{-}}$introduced in Section 2A. Let $l_{K}: K \hookrightarrow B$ be the inclusion. Then one easily checks (see [Castella and Longo 2016, Theorem 1.2]) that for all $n, r \geqslant 0$ the points

$$
\widetilde{P}_{p^{n}, r}:=\left[\left(l_{K}, \varsigma^{(n+r)}\right)\right] \in \widetilde{X}_{N^{+}, r}
$$

are Heegner points of conductor $p^{n+r}$ with the following properties:

- Field of definition: $\widetilde{P}_{p^{n}, r} \in H^{0}\left(L_{p^{n}, r}, \widetilde{X}_{N^{+}, r}\right)$, where $L_{p^{n}, r}:=H_{p^{n+r}}\left(\boldsymbol{\mu}_{p^{r}}\right)$ and $H_{c}$ is the ring class field of $K$ of conductor $c$.
- Galois equivariance: for all $\sigma \in \operatorname{Gal}\left(L_{p^{n}, r} / H_{p^{n+r}}\right)$, we have

$$
\widetilde{P}_{p^{n}, r}^{\sigma}=\langle\vartheta(\sigma)\rangle \cdot \widetilde{P}_{p^{n}, r},
$$

where $\vartheta: \operatorname{Gal}\left(L_{p^{n}, r} / H_{p^{n+r}}\right) \rightarrow \mathbb{Z}_{p}^{\times} /\{ \pm 1\}$ is such that $\vartheta^{2}=\varepsilon_{\text {cyc }}$.

- Horizontal compatibility: if $r>1$, then

$$
\sum_{\sigma \in \operatorname{Gal}\left(L_{p^{n}, r} / L_{p^{n-1}, r}\right)} \widetilde{\alpha}_{r}\left(\widetilde{P}_{p^{n}, r}^{\sigma}\right)=U_{p} \cdot \widetilde{P}_{p^{n}, r-1}
$$

where $\widetilde{\alpha}_{r}: \widetilde{X}_{N^{+}, r} \rightarrow \widetilde{X}_{N^{+}, r-1}$ is the map induced by the inclusion $U_{N^{+}, r} \subseteq$ $U_{N^{+}, r-1}$.

- Vertical Compatibility: if $n>0$, then

$$
\sum_{\sigma \in \operatorname{Gal}\left(L_{p^{n}, r} / L_{p^{n-1}, r}\right)}{\widetilde{P_{p}^{n}, r}}_{\sigma}^{\sigma}=U_{p} \cdot \widetilde{P}_{p^{n-1}, r} .
$$

Remark 2.2. We will only consider the points $\widetilde{P}_{p^{n}, r}$ for a fixed a value of $N^{-}$ (which amounts to fixing the quaternion algebra $B / \mathbb{Q}$ ), but it will be fundamental to consider different values of $N^{+}$, and the relations between the corresponding $\widetilde{P}_{p^{n}, r}$ (which clearly depend on $N^{+}$) under various degeneracy maps.

2C. Critical character. Factor the $p$-adic cyclotomic character as

$$
\varepsilon_{\mathrm{cyc}}=\varepsilon_{\mathrm{tame}} \cdot \varepsilon_{\mathrm{wild}}: G_{\mathbb{Q}} \rightarrow \mathbb{Z}_{p}^{\times} \simeq \mu_{p-1} \times\left(1+p \mathbb{Z}_{p}\right)
$$

and define the critical character $\Theta: G_{\mathbb{Q}} \rightarrow \mathbb{Z}_{p} \llbracket 1+p \mathbb{Z}_{p} \rrbracket^{\times}$by

$$
\begin{equation*}
\Theta(\sigma)=\left[\varepsilon_{\text {wild }}^{1 / 2}(\sigma)\right] \tag{3}
\end{equation*}
$$

where $\varepsilon_{\text {wild }}^{1 / 2}$ is the unique square root of $\varepsilon_{\text {wild }}$ taking values in $1+p \mathbb{Z}_{p}$, and the map $[\cdot]: 1+p \mathbb{Z}_{p} \rightarrow \mathbb{Z}_{p} \llbracket 1+p \mathbb{Z}_{p} \rrbracket^{\times}$is given by the inclusion as group-like elements.

2D. Big Heegner points. Recall the Shimura curves $\widetilde{X}_{N^{+}, p^{r}}$ from Section 2A, and set

$$
\mathfrak{D}_{N^{+}, r}:=e^{\operatorname{ord}}\left(\operatorname{Div}\left(\tilde{X}_{N^{+}, r}\right) \otimes_{\mathbb{Z}} \mathbb{Z}_{p}\right)
$$

By the Jacquet-Langlands correspondence, $\mathfrak{D}_{N^{+}, r}$ is naturally endowed with an action of the Hecke algebra $\mathbb{T}_{N, r}^{N^{-}}$. Let $\left(\mathbb{T}_{N, r}^{N^{-}}\right)^{\dagger}$ be the free $\mathbb{T}_{N, r}^{N^{-}}$-module of rank one equipped with the Galois action via the inverse of the critical character $\Theta$, and set

$$
\mathfrak{D}_{N^{+}, r}^{\dagger}:=\mathfrak{D}_{N^{+}, r} \otimes_{\mathbb{T}_{N, r}^{N^{-}}}\left(\mathbb{T}_{N, r}^{N^{-}}\right)^{\dagger} .
$$

Let $\widetilde{P}_{p^{n}, r} \in \widetilde{X}_{N^{+}, r}$ be the system of Heegner points of Section 2B, and denote by $\mathcal{P}_{p^{n}, r}$ the image of $e^{\text {ord }} \widetilde{P}_{p^{n}, r}$ in $\mathfrak{D}_{N^{+}, r}$. By the Galois equivariance of $\widetilde{P}_{p^{n}, r}$ (see [Longo and Vigni 2011, §7.1]), we have

$$
\mathcal{P}_{p^{n}, r}^{\sigma}=\Theta(\sigma) \cdot \mathcal{P}_{p^{n}, r}
$$

for all $\sigma \in \operatorname{Gal}\left(L_{p^{n}, r} / H_{p^{n+r}}\right)$, and hence $\mathcal{P}_{p^{n}, r}$ defines an element

$$
\begin{equation*}
\mathcal{P}_{p^{n}, r} \otimes \zeta_{r} \in H^{0}\left(H_{p^{n+r}}, \mathfrak{D}_{N^{+}, r}^{\dagger}\right) \tag{4}
\end{equation*}
$$

In the next section we shall see how this system of points, for varying $n$ and $r$, can be used to construct various anticyclotomic $p$-adic $L$-functions.

## 3. Anticyclotomic $\boldsymbol{p}$-adic $\boldsymbol{L}$-functions

3A. Multiplicity one. Keep the notation introduced in Section 2. For each integer $k \geqslant 2$, denote by $L_{k}(R)$ the set of polynomials of degree less than or equal to $k-2$ with coefficients in a ring $R$, and define

$$
\mathfrak{J}_{N^{+}, r, k}:=e^{\text {ord }} H_{0}\left(\widetilde{X}_{N^{+}, r}, \mathcal{L}_{k}\left(\mathbb{Z}_{p}\right)\right)
$$

where $\mathcal{L}_{k}\left(\mathbb{Z}_{p}\right)$ is the local system on $\widetilde{X}_{N^{+}, r}$ associated with $L_{k}\left(\mathbb{Z}_{p}\right)$. The module $\mathfrak{J}_{N^{+}, r, k}$ is endowed with an action of the Hecke algebra $\mathbb{T}_{N, r, k}^{N^{-}}$and with perfect "intersection pairing":

$$
\begin{equation*}
\langle\cdot, \cdot\rangle_{k}: \mathfrak{J}_{N^{+}, r, k} \times \mathfrak{J}_{N^{+}, r, k} \rightarrow \mathbb{Q}_{p} \tag{5}
\end{equation*}
$$

(see [Chida and Hsieh 2016, Equation (3.9)]) with respect to which the Hecke operators are self-adjoint.
Theorem 3.1. Let $\mathfrak{m}$ be a maximal ideal of $\mathbb{T}_{N, r, k}^{N^{-}}$whose residual representation is irreducible and satisfies Assumption (SU). Then $\left(\mathfrak{J}_{N^{+}, r, k}\right)_{\mathfrak{m}}$ is free of rank one over $\left(\mathbb{T}_{N, r, k}^{N^{-}}\right)_{\mathfrak{m}}$. In particular, there is a $\left(\mathbb{T}_{N, r, k}^{N^{-}}\right)_{\mathfrak{m}}$-module isomorphism

$$
\left(\mathfrak{J}_{N^{+}, r, k}\right)_{\mathfrak{m}} \stackrel{\alpha_{N, r, k}}{\sim}\left(\mathbb{T}_{N, r, k}^{N^{-}}\right)_{\mathfrak{m}} .
$$

Proof. If $k=2$ and $r=1$, this follows by combining Theorem 6.2 and Proposition 6.5 of [Pollack and Weston 2011]. The general case will be deduced from this case in Section 3C using Hida theory.

Let $f \in S_{k}\left(\Gamma_{0,1}\left(N, p^{r}\right)\right)$ be an $N^{-}$-new eigenform, and suppose that $\mathfrak{m}$ is the maximal ideal of $\mathbb{T}_{N, r, k}^{N^{-}}$containing the kernel of the associated $\mathbb{Z}_{p}$-algebra homomorphism

$$
\pi_{f}:\left(\mathbb{T}_{N, r, k}^{N^{-}}\right)_{\mathfrak{m}} \rightarrow \mathcal{O}
$$

where $\mathcal{O}$ is the finite extension of $\mathbb{Z}_{p}$ generated by the Fourier coefficients of $f$. Composing $\pi_{f}$ with an isomorphism $\alpha_{N, r, k}$ as in Theorem 3.1, we obtain an $\mathcal{O}$ valued functional

$$
\psi_{f}:\left(\mathfrak{J}_{N^{+}, r, k}\right)_{\mathfrak{m}} \rightarrow \mathcal{O}
$$

By the duality (5), the map $\psi_{f}$ corresponds to a generator $g_{f}$ of the $\pi_{f}$-isotypical component of $\mathfrak{J}_{N^{+}, r, k}$, and following [Pollack and Weston 2011, §2.1] and [Chida and Hsieh 2016, §4.1] we define the Gross period $\Omega_{f, N^{-}}$attached to $f$ by

$$
\begin{equation*}
\Omega_{f, N^{-}}:=\frac{(f, f)_{\Gamma_{0}(N)}}{\left\langle g_{f}, g_{f}\right\rangle_{k}} \tag{6}
\end{equation*}
$$

Remark 3.2. By Vatsal's work [2003] (see also [Pollack and Weston 2011, Theorem 2.3] and [Chida and Hsieh 2016, §5.4]), the anticyclotomic $p$-adic $L$-functions
$L_{p}(f / K)$ in Theorem 3.14 below (normalized by the complex period $\Omega_{f, N^{-}}$) have vanishing $\mu$-invariant. The preceding uniform description of $\psi_{f}$ for all $f$ with a common maximal ideal $\mathfrak{m}$ will allow us to show that this property is preserved in Hida families.

3B. One-variable p-adic L-functions. Denote by $\Gamma$ the Galois group of the anticyclotomic $\mathbb{Z}_{p}$-extension $K_{\infty} / K$. For each $n$, let $K_{n} \subset K_{\infty}$ be defined by $\operatorname{Gal}\left(K_{n} / K\right) \simeq \mathbb{Z} / p^{n} \mathbb{Z}$ and let $\Gamma_{n}$ be the subgroup of $\Gamma$ such that $\Gamma / \Gamma_{n} \simeq \operatorname{Gal}\left(K_{n} / K\right)$.

Let $\mathcal{P}_{p^{n+1}, r} \otimes \zeta_{r} \in H^{0}\left(H_{p^{n+1+r}}, \mathfrak{D}_{N^{+}, r}^{\dagger}\right)$ be the Heegner point of conductor $p^{n+1}$, and define

$$
\begin{equation*}
\mathcal{Q}_{n, r}:=\operatorname{Cor}_{H_{p^{n+1+r}} / K_{n}}\left(\mathcal{P}_{p^{n+1}, r} \otimes \zeta_{r}\right) \in H^{0}\left(K_{n}, \mathfrak{D}_{N^{+}, r}^{\dagger}\right) \tag{7}
\end{equation*}
$$

with a slight abuse of notation, we also denote by $\mathcal{Q}_{n, r}$ its image under the natural map

$$
H^{0}\left(K_{n}, \mathfrak{D}_{N^{+}, r}^{\dagger}\right) \xrightarrow{\subseteq} \mathfrak{D}_{N^{+}, r} \longrightarrow \mathfrak{J}_{N^{+}, r}
$$

composed with localization at $\mathfrak{m}$, where $\mathfrak{J}_{N^{+}, r}:=\mathfrak{J}_{N^{+}, r, 2}$.
Definition 3.3. For any open subset $\sigma \Gamma_{n}$ of $\Gamma$, define

$$
\mu_{r}\left(\sigma \Gamma_{n}\right):=U_{p}^{-n} \cdot \mathcal{Q}_{n, r}^{\sigma} \in\left(\mathfrak{J}_{N^{+}, r}\right)_{\mathfrak{m}} .
$$

Proposition 3.4. The rule $\mu_{r}$ is a measure on $\Gamma$.
Proof. This follows immediately from the "horizontal compatibility" of Heegner points.

3C. Gross periods in Hida families. Keep the notation of Section 3A, and let

$$
\left.\left(\mathfrak{J}_{N^{+}}\right)_{\mathfrak{m}}:={\underset{r}{l i m}}_{\lim _{N^{+}, r}}\right)_{\mathfrak{m}},
$$

which is naturally equipped with an action of the big Hecke algebra $\mathbb{T}_{N}^{N^{-}}=\lim _{{ }_{r}} \mathbb{T}_{N, r}^{N^{-}}$.
Theorem 3.5. Let $\mathfrak{m}$ be a maximal ideal of $\mathbb{T}_{N}^{N^{-}}$whose residual representation is irreducible and satisfies Assumption (SU). Then $\left(\mathfrak{J}_{N^{+}}\right)_{\mathfrak{m}}$ is free of rank one over $\left(\mathbb{T}_{N}^{N^{-}}\right)_{\mathfrak{m}}$. In particular, there is a $\left(\mathbb{T}_{N}^{N^{-}}\right)_{\mathfrak{m}}$-module isomorphism

$$
\left(\mathfrak{J}_{N^{+}}\right)_{\mathfrak{m}} \stackrel{\alpha_{N}}{\simeq}\left(\mathbb{T}_{N}^{N^{-}}\right)_{\mathfrak{m}}
$$

Proof. As in [Emerton et al. 2006, Proposition 3.3.1]. Note that the version of Hida's control theorem in our context is provided by [Hida 1988, Theorem 9.4].

We can now conclude the proof of Theorem 3.1 just as in [Emerton et al. 2006, §3.3]. For the convenience of the reader, we include the argument here.

Proof of Theorem 3.1. Let $\wp_{N, r, k}$ be the product of all the arithmetic primes of $\mathbb{T}_{N}^{N^{-}}$ of weight $k$ which become trivial upon restriction to $1+p^{r} \mathbb{Z}_{p}$. By [Hida 1988, Theorem 9.4], we then have

$$
\begin{equation*}
\left(\mathfrak{J}_{N^{+}}\right)_{\mathfrak{m}} \otimes \mathbb{T}_{N}^{N^{-}} / \wp_{N, r, k} \simeq\left(\mathfrak{J}_{N^{+}, r, k}\right)_{\mathfrak{m}_{r, k}} \tag{8}
\end{equation*}
$$

where $\mathfrak{m}_{r, k}$ is the maximal ideal of $\mathbb{T}_{N, r, k}^{N^{-}}$induced by $\mathfrak{m}$. Since $\left(\mathfrak{J}_{N^{+}}\right)_{\mathfrak{m}}$ is free of rank one over $\mathbb{T}_{N}^{N^{-}}$by Theorem 3.5, it follows that $\left(\mathfrak{J}_{N^{+}, r, k}\right)_{\mathfrak{m}_{r, k}}$ is free of rank one over $\mathbb{T}_{N}^{N^{-}} / \wp_{N, r, k} \simeq \mathbb{T}_{N, r, k}^{N^{-}}$, as was to be shown.
Remark 3.6. In the above proofs we made crucial use of [Hida 1988, Theorem 9.4], which is stated in the context of totally definite quaternion algebras which are unramified at all finite places, since this is the only relevant case for the study of Hilbert modular forms over totally real number fields of even degree. However, the proofs immediately extend to the (simpler) situation of definite quaternion algebras over $\mathbb{Q}$.

3D. Two-variable p-adic L-functions. By the "vertical compatibility" satisfied by Heegner points, the points

$$
U_{p}^{-r} \cdot \mathcal{Q}_{n, r} \in\left(\mathfrak{J}_{N^{+}, r}\right)_{\mathfrak{m}}
$$

are compatible for varying $r$, thus defining an element

$$
\mathcal{Q}_{n}:={\underset{\check{l}}{r}}^{\lim _{p}^{-r}} \cdot \mathcal{Q}_{n, r} \in\left(\mathfrak{J}_{N^{+}}\right)_{\mathfrak{m}}
$$

Definition 3.7. For any open subset $\sigma \Gamma_{n}$ of $\Gamma$, define

$$
\mu\left(\sigma \Gamma_{n}\right):=U_{p}^{-n} \cdot \mathcal{Q}_{n}^{\sigma} \in\left(\mathfrak{J}_{N^{+}}\right)_{\mathfrak{m}}
$$

Proposition 3.8. The rule $\mu$ is a measure on $\Gamma$.
Proof. This follows immediately from the "horizontal compatibility" of Heegner points.

Upon the choice of an isomorphism $\alpha_{N}$ as in Theorem 3.5, we may regard $\mu$ as an element

$$
\mathcal{L}(\mathfrak{m}, N) \in\left(\mathbb{T}_{N}^{N^{-}}\right)_{\mathfrak{m}} \hat{\otimes}_{\mathbb{Z}_{p}} \mathbb{Z}_{p} \llbracket \Gamma \rrbracket .
$$

Denoting by $\mathcal{L}(\mathfrak{m}, N)^{*}$ the image of $\mathcal{L}(\mathfrak{m}, N)$ under the involution induced by $\gamma \mapsto \gamma^{-1}$ on group-like elements, we set

$$
L(\mathfrak{m}, N):=\mathcal{L}(\mathfrak{m}, N) \cdot \mathcal{L}(\mathfrak{m}, N)^{*}
$$

to which we will refer as the two-variable p-adic L-function attached to $\left(\mathbb{T}_{N}^{N^{-}}\right)_{\mathfrak{m}}$.

## 3E. Two-variable p-adic L-functions on branches of the Hida family. Let $\mathbb{T}_{\Sigma}$

 be the universal $p$-ordinary Hecke algebra$$
\begin{equation*}
\mathbb{T}_{\Sigma}:=\left(\mathbb{T}_{N(\Sigma)}^{\prime}\right)_{\mathfrak{m}} \simeq\left(\mathbb{T}_{N(\Sigma)}^{N^{-}}\right)_{\mathfrak{n}} \tag{9}
\end{equation*}
$$

associated with a $\bmod p$ representation $\bar{\rho}$ and a finite set of primes $\Sigma$ as in Section 1C.

Remark 3.9. Recall that $N^{-} \mid N(\bar{\rho})$ by Assumption (SU). Throughout the following, it will be further assumed that every prime factor of $N(\Sigma) / N^{-}$splits in $K$. In particular, every prime $\ell \in \Sigma$ splits in $K$, and any $f \in \mathcal{H}^{-}(\bar{\rho})=\operatorname{Spec}\left(\mathbb{T}_{\Sigma}\right)$ has tame level $N_{f}$ with

$$
N_{f}^{-}=N(\bar{\rho})^{-}=N^{-} .
$$

The construction of the preceding section produces a two-variable $p$-adic $L$ function

$$
L(\mathfrak{n}, N(\Sigma)) \in\left(\mathbb{T}_{N(\Sigma)}^{N^{-}}\right)_{\mathfrak{n}} \hat{\mathbb{Q}}_{\mathbb{Z}_{p}} \mathbb{Z}_{p} \llbracket \Gamma \rrbracket,
$$

which combined with the isomorphism (9) yields an element

$$
L_{\Sigma}(\bar{\rho}) \in \mathbb{T}_{\Sigma} \hat{\otimes}_{\mathbb{Z}_{p}} \mathbb{Z}_{p} \llbracket \Gamma \rrbracket .
$$

If $\mathfrak{a}$ is a minimal prime of $\mathbb{T}_{\Sigma}$, we thus obtain an element

$$
L_{\Sigma}(\bar{\rho}, \mathfrak{a}) \in \mathbb{T}(\mathfrak{a})^{\circ} \hat{\otimes}_{\mathbb{Z}_{p}} \mathbb{Z}_{p} \llbracket \Gamma \rrbracket
$$

by reducing $L_{\Sigma}(\bar{\rho}) \bmod \mathfrak{a}$ (see Section 1D). On the other hand, if we let $\mathfrak{m}$ denote the inverse image of the maximal ideal of $\mathbb{T}(\mathfrak{a})^{\circ}$ under the composite surjection

$$
\begin{equation*}
\mathbb{T}_{N(\mathfrak{a})}^{N^{-}} \rightarrow \mathbb{T}_{N(\mathfrak{a})}^{\text {new }} \rightarrow \mathbb{T}_{N(\mathfrak{a})}^{\text {new }} / \mathfrak{a}^{\prime}=\mathbb{T}(\mathfrak{a})^{\circ}, \tag{10}
\end{equation*}
$$

the construction of the preceding section yields an $L$-function

$$
L(\mathfrak{m}, N(\mathfrak{a})) \in\left(\mathbb{T}_{N(\mathfrak{a})}^{N^{-}}\right)_{\mathfrak{m}} \hat{\otimes}_{\mathbb{Z}_{p}} \mathbb{Z}_{p} \llbracket \Gamma \rrbracket
$$

giving rise, via (10), to a second element

$$
L(\bar{\rho}, \mathfrak{a}) \in \mathbb{T}(\mathfrak{a})^{\circ} \hat{\otimes}_{\mathbb{Z}_{p}} \mathbb{Z}_{p} \llbracket \Gamma \rrbracket .
$$

It is natural to compare $L_{\Sigma}(\bar{\rho}, \mathfrak{a})$ and $L(\bar{\rho}, \mathfrak{a})$, a task that is carried out in the next section, and provides the key for understanding the variation of analytic Iwasawa invariants.

3F. Comparison. Write $\Sigma=\left\{\ell_{1}, \ldots, \ell_{n}\right\}$ and for each $\ell=\ell_{i} \in \Sigma$, let $e_{\ell}$ be the valuation of $N(\Sigma) / N(\mathfrak{a})$ at $\ell$, and define the reciprocal Euler factor $E_{\ell}(\mathfrak{a}, X) \in$ $\mathbb{T}(\mathfrak{a})^{\circ}[X]$ by

$$
E_{\ell}(\mathfrak{a}, X):= \begin{cases}1 & \text { if } e_{\ell}=0 \\ 1-\left(T_{\ell} \bmod \mathfrak{a}^{\prime}\right) \Theta^{-1}(\ell) X & \text { if } e_{\ell}=1 \\ 1-\left(T_{\ell} \bmod \mathfrak{a}^{\prime}\right) \Theta^{-1}(\ell) X+\ell X^{2} & \text { if } e_{\ell}=2\end{cases}
$$

Also, writing $\ell=\mathfrak{l}$, define $E_{\ell}(\mathfrak{a}) \in \mathbb{T}(\mathfrak{a})^{\circ} \hat{\otimes}_{\mathbb{Z}_{p}} \mathbb{Z}_{p} \llbracket \Gamma \rrbracket$ by

$$
\begin{equation*}
E_{\ell}(\mathfrak{a}):=E_{\ell}\left(\mathfrak{a}, \ell^{-1} \gamma_{\mathfrak{l}}\right) \cdot E_{\ell}\left(\mathfrak{a}, \ell^{-1} \gamma_{\hat{\mathfrak{N}}}\right), \tag{11}
\end{equation*}
$$

where $\gamma_{\mathfrak{l}}, \gamma_{\mathfrak{l}}$ are arithmetic Frobenius maps at $\mathfrak{l}$, $\overline{\mathfrak{l}}$ in $\Gamma$, respectively, and put $E_{\Sigma}(\mathfrak{a}):=$ $\prod_{\ell \in \Sigma} E_{\ell}(\mathfrak{a})$.

Recall that $N^{-}|N(\mathfrak{a})| N(\Sigma)$ and set

$$
N(\mathfrak{a})^{+}:=N(\mathfrak{a}) / N^{-}, \quad N(\Sigma)^{+}:=N(\Sigma) / N^{-}
$$

both of which consist entirely of prime factors which split in $K$. The purpose of this section is to prove the following result.

Theorem 3.10. There is an isomorphism of $\mathbb{T}(\mathfrak{a})^{\circ}$-modules

$$
\mathbb{T}(\mathfrak{a})^{\circ} \otimes_{\left(\mathbb{T}_{N(\Sigma)}^{N-}\right)_{\mathfrak{n}}}\left(\mathfrak{J}_{N(\Sigma)^{+}}\right)_{\mathfrak{n}} \simeq \mathbb{T}(\mathfrak{a})^{\circ} \otimes_{\left(\mathbb{T}_{N(\mathfrak{a})}^{N-}\right)_{\mathfrak{m}}}\left(\mathfrak{J}_{N(\mathfrak{a})^{+}}\right)_{\mathfrak{m}}
$$

such that the map induced on the corresponding spaces of measures valued in these modules sends $L_{\Sigma}(\bar{\rho}, \mathfrak{a})$ to $E_{\Sigma}(\mathfrak{a}) \cdot L(\bar{\rho}, \mathfrak{a})$.

Proof. The proof follows closely the constructions and arguments in [Emerton et al. 2006, §3.8], suitably adapted to the quaternionic setting at hand. Let $r \geqslant 1$. If $M$ is any positive integer with $\left(M, p N^{-}\right)=1$, and $d^{\prime} \mid d$ are divisors of $M$, we have degeneracy maps

$$
B_{d, d^{\prime}}: \widetilde{X}_{M, r} \rightarrow \widetilde{X}_{M / d, r}
$$

induced by $(\Psi, g) \mapsto\left(\Psi, \pi_{d^{\prime}} g\right)$, where $\pi_{d^{\prime}} \in \widehat{B}^{\times}$has local component $\left(\begin{array}{ll}1 & 0 \\ 0 & \ell^{\text {val }}\left(d^{\prime}\right)\end{array}\right)$ at every prime $\ell \mid d^{\prime}$ and 1 outside $d^{\prime}$. We thus obtain a map on homology

$$
\left(B_{d, d^{\prime}}\right)_{*}: e^{\text {ord }} H_{0}\left(\tilde{X}_{M, r}, \mathbb{Z}_{p}\right) \rightarrow e^{\text {ord }} H_{0}\left(\tilde{X}_{M / d, r}, \mathbb{Z}_{p}\right)
$$

and we may define

$$
\begin{equation*}
\epsilon_{r}: e^{\text {ord }} H_{0}\left(\tilde{X}_{N(\Sigma)^{+}, r}, \mathbb{Z}_{p}\right) \rightarrow e^{\text {ord }} H_{0}\left(\tilde{X}_{N(\mathfrak{a})^{+}, r}, \mathbb{Z}_{p}\right) \tag{12}
\end{equation*}
$$

by $\epsilon_{r}:=\epsilon\left(\ell_{n}\right) \circ \cdots \circ \epsilon\left(\ell_{1}\right)$, where for every $\ell=\ell_{i} \in \Sigma$ we put

$$
\epsilon(\ell):= \begin{cases}1 & \text { if } e_{\ell}=0 \\ \left(B_{\ell, 1}\right)_{*}-\left(B_{\ell, \ell}\right)_{*} \ell^{-1} T_{\ell} & \text { if } e_{\ell}=1 \\ \left(B_{\ell^{2}, 1}\right)_{*}-\left(B_{\ell^{2}, \ell}\right)_{*} \ell^{-1} T_{\ell}+\left(B_{\ell^{2}, \ell^{2}}\right)_{*} \ell^{-1}\langle\ell\rangle_{N(\mathfrak{a}) p} & \text { if } e_{\ell}=2\end{cases}
$$

As before, let $M$ be a positive integer with $\left(M, p N^{-}\right)=1$ all of whose prime factors split in $K$, and let $\ell \nmid M p$ be a prime which also splits in $K$. We shall adopt the following simplifying notation for the system of points $\widetilde{P}_{p^{n}, r} \in \widetilde{X}_{N^{+}, r}$ constructed in Section 2B:

$$
P:=\widetilde{P}_{p^{n}, r}^{(M)} \in \widetilde{X}_{M, r}, \quad P^{(\ell)}:=\widetilde{P}_{p^{n}, r}^{(M \ell)} \in \widetilde{X}_{M \ell, r}, \quad P^{\left(\ell^{2}\right)}:=\widetilde{P}_{p^{n}, r}^{\left(M \ell^{2}\right)} \in \widetilde{X}_{M \ell^{2}, r}
$$

It is easy to check that we have the following relations in $\widetilde{X}_{M, r}$ :

$$
\begin{aligned}
& \left(B_{\ell, 1}\right)_{*}\left(P^{(\ell)}\right)=P, \quad\left(B_{\ell, \ell}\right)_{*}\left(P^{(\ell)}\right)=P^{\sigma_{\mathrm{I}}}, \quad\left(B_{\ell^{2}, 1}\right)_{*}\left(P^{\left(\ell^{2}\right)}\right)=P, \\
& \left(B_{\ell^{2}, \ell}\right)_{*}\left(P^{\left(\ell^{2}\right)}\right)=P^{\sigma_{\mathrm{I}}}, \quad\left(B_{\ell^{2}, \ell^{2}}\right)_{*}\left(P^{\left(\ell^{2}\right)}\right)=P^{\sigma_{\mathrm{I}}^{2}},
\end{aligned}
$$

where $\sigma_{\mathfrak{l}} \in \operatorname{Gal}\left(L_{p^{n}, r} / K\right)$ is a Frobenius element at a prime $\mathfrak{l} \mid \ell$. Letting $\mathcal{P}$ denote the image of $e^{\text {ord }} P$ in $\mathfrak{D}_{M, r}$, and defining $\mathcal{P}^{(\ell)} \in \mathfrak{D}_{M \ell, r}$ and $\mathcal{P}^{\left(\ell^{2}\right)} \in \mathfrak{D}_{M \ell^{2}, r}$ similarly, it follows that

$$
\begin{aligned}
\left(B_{\ell, 1}\right)_{*}\left(\mathcal{P}^{(\ell)} \otimes \zeta_{r}\right) & =\mathcal{P} \otimes \zeta_{r}, \\
\left(B_{\ell, \ell}\right)_{*}\left(\mathcal{P}^{(\ell)} \otimes \zeta_{r}\right) & =\mathcal{P}^{\sigma_{\mathrm{I}}} \otimes \zeta_{r}=\Theta^{-1}\left(\sigma_{\mathfrak{l}}\right) \cdot\left(\mathcal{P} \otimes \zeta_{r}\right)^{\sigma_{\mathrm{I}}}, \\
\left(B_{\ell^{2}, 1}\right)_{*}\left(\mathcal{P}^{\left(\ell^{2}\right)} \otimes \zeta_{r}\right) & =\mathcal{P} \otimes \zeta_{r}, \\
\left(B_{\ell^{2}, \ell}\right)_{*}\left(\mathcal{P}^{\left(\ell^{2}\right)} \otimes \zeta_{r}\right) & =\mathcal{P}^{\sigma_{\mathrm{I}}} \otimes \zeta_{r}=\Theta^{-1}\left(\sigma_{\mathfrak{l}}\right) \cdot\left(\mathcal{P} \otimes \zeta_{r}\right)^{\sigma_{\mathrm{I}}}, \\
\left(B_{\left.\ell^{2}, \ell^{2}\right)_{*}}\left(\mathcal{P}^{\left(\ell^{2}\right)} \otimes \zeta_{r}\right)\right. & =\mathcal{P}^{\sigma_{\mathrm{I}}^{2}} \otimes \zeta_{r}=\Theta^{-2}\left(\sigma_{\mathfrak{l}}\right) \cdot\left(\mathcal{P} \otimes \zeta_{r}\right)^{\sigma_{\mathrm{I}}}
\end{aligned}
$$

as elements in $\mathfrak{D}_{M, r}^{\dagger}$. Finally, setting $\mathcal{Q}:=\operatorname{Cor}_{H_{p^{n+1+r}} / K_{n}}(\mathcal{P}) \in H^{0}\left(K_{n}, \mathfrak{D}_{M, r}^{\dagger}\right)$, and defining $\mathcal{Q}^{(\ell)} \in H^{0}\left(K_{n}, \mathfrak{D}_{M \ell, r}^{\dagger}\right)$ and $\mathcal{Q}^{\left(\ell^{2}\right)} \in H^{0}\left(K_{n}, \mathfrak{D}_{M \ell^{2}, r}^{\dagger}\right)$ similarly, we see that

$$
\begin{array}{rlrl}
\left(B_{\ell, 1}\right)_{*}\left(\mathcal{Q}^{(\ell)}\right) & =\mathcal{Q}, & \left(B_{\ell, \ell}\right)_{*}\left(\mathcal{Q}^{(\ell)}\right)=\Theta^{-1}\left(\sigma_{\mathfrak{l}}\right) \cdot \mathcal{Q}^{\sigma_{\mathfrak{l}}} \\
\left(B_{\ell^{2}, 1}\right)_{*}\left(\mathcal{Q}^{\left(\ell^{2}\right)}\right) & =\mathcal{Q}, & & \\
\left(B_{\ell^{2}, \ell}\right)_{*}\left(\mathcal{Q}^{\left(\ell^{2}\right)}\right) & =\Theta^{-1}\left(\sigma_{\mathfrak{l}}\right) \cdot \mathcal{Q}^{\sigma_{\mathfrak{l}},} & \left(B_{\ell^{2}, \ell^{2}}\right)_{*}\left(\mathcal{Q}^{\left(\ell^{2}\right)}\right)=\Theta^{-2}\left(\sigma_{\mathfrak{l}}\right) \cdot \mathcal{Q}^{\sigma_{\mathfrak{l}}^{2}}
\end{array}
$$

in $H^{0}\left(K_{n}, \mathfrak{D}_{M, r}^{\dagger}\right)$. Each of these equalities is checked by an explicit calculation. For example, for the second one:

$$
\begin{aligned}
\left(B_{\ell, \ell}\right)_{*}\left(\mathcal{Q}^{(\ell)}\right) & =\left(B_{\ell, \ell}\right)_{*}\left(\operatorname{Cor}_{H_{p^{n+1+r}} / K_{n}}\left(\mathcal{P}^{(\ell)} \otimes \zeta_{r}\right)\right) \\
& =\left(B_{\ell, \ell}\right)_{*}\left(\left(\sum_{\sigma \in \operatorname{Gal}\left(H_{p^{n+1+r}} / K_{n}\right)} \Theta\left(\tilde{\sigma}^{-1}\right) \cdot\left(\mathcal{P}^{(\ell)}\right)^{\tilde{\sigma}}\right) \otimes \zeta_{r}\right) \\
& =\sum_{\sigma \in \operatorname{Gal}\left(H_{p^{n+1+r}} / K_{n}\right)} \Theta\left(\tilde{\sigma}^{-1}\right) \cdot\left(B_{\ell, \ell}\right)_{*}\left(\left(\mathcal{P}^{(\ell)}\right)^{\tilde{\sigma}} \otimes \zeta_{r}\right) \\
& =\sum_{\sigma \in \operatorname{Gal}\left(H_{p^{n+1+r}} / K_{n}\right)} \Theta\left(\tilde{\sigma}^{-1}\right) \Theta^{-1}\left(\sigma_{\mathfrak{l}}\right) \cdot\left(\mathcal{P}^{\tilde{\sigma}} \otimes \zeta_{r}\right)^{\sigma_{\mathfrak{l}}} \\
& =\Theta^{-1}\left(\sigma_{\mathfrak{l}}\right) \cdot \mathcal{Q}^{\sigma_{\mathfrak{l}}} .
\end{aligned}
$$

Now let $\mathcal{Q}_{n, r} \in \mathfrak{J}_{N(\Sigma)^{+}, r}$ be as in (7) with $N=N(\Sigma)$. Using the above formulae, we easily see that of any finite order character $\chi$ of $\Gamma$ of conductor $p^{n}$, the effect of $\epsilon_{r}$ on the element $\sum_{\sigma \in \Gamma / \Gamma_{n}} \chi(\sigma) \mathcal{Q}_{n, r}^{\sigma}$ is given by multiplication by

$$
\begin{aligned}
& \prod_{e_{\ell_{i}}=1}\left(1-(\chi \Theta)^{-1}\left(\sigma_{\mathfrak{l}_{i}}\right) \ell_{i}^{-1} T_{\ell_{i}}\right) \\
& \qquad \prod_{e_{\ell_{i}}=2}\left(1-(\chi \Theta)^{-1}\left(\sigma_{\mathfrak{I}_{i}}\right) \ell_{i}^{-1} T_{\ell_{i}}+(\chi \Theta)^{-2}\left(\sigma_{\mathfrak{\tau}_{i}}\right) \ell_{i}^{-1}\left\langle\ell_{i}\right\rangle_{N(\mathfrak{a}) p}\right)
\end{aligned}
$$

Similarly, we see $\epsilon_{r}$ has the effect of multiplying the element $\sum_{\sigma \in \Gamma / \Gamma_{n}} \chi^{-1}(\sigma) \mathcal{Q}_{n, r}^{\sigma}$ by

$$
\begin{aligned}
& \prod_{e_{i}=1}\left(1-\left(\chi^{-1} \Theta\right)^{-1}\left(\sigma_{\mathfrak{I}_{i}}\right) \ell_{i}^{-1} T_{\ell_{i}}\right) \\
& \prod_{e_{\ell_{i}}=2}\left(1-\left(\chi^{-1} \Theta\right)^{-1}\left(\sigma_{\mathfrak{\tau}_{i}}\right) \ell_{i}^{-1} T_{\ell_{i}}+\left(\chi^{-1} \Theta\right)^{-2}\left(\sigma_{\mathfrak{I}_{i}}\right) \ell_{i}^{-1}\left\langle\ell_{i}\right\rangle_{N(\mathfrak{a}) p}\right) .
\end{aligned}
$$

Hence, using the relations

$$
\chi\left(\sigma_{\overline{\mathfrak{T}}_{i}}\right)=\chi^{-1}\left(\sigma_{\mathfrak{I}_{i}}\right), \quad \Theta\left(\sigma_{\mathfrak{l}_{i}}\right)=\Theta\left(\sigma_{\overline{\mathrm{T}}_{i}}\right)=\theta\left(\ell_{i}\right), \quad \theta^{2}\left(\ell_{i}\right)=\left\langle\ell_{i}\right\rangle_{N(\mathfrak{a}) p}
$$

it follows that the effect of $\epsilon_{r}$ on the product of the above two elements is given by multiplication by

$$
\prod_{\substack{\mathfrak{i}_{i} \mid \ell_{i} \\ e_{\ell_{i}}=1}}\left(1-\chi\left(\sigma_{\mathfrak{I}_{i}}\right) \theta^{-1}\left(\ell_{i}\right) \ell_{i}^{-1} T_{\ell_{i}}\right) \prod_{\substack{\mathfrak{i}_{i} \mid \ell_{i} \\ e_{\ell_{i}}=2}}\left(1-\chi\left(\sigma_{\mathfrak{I}_{i}}\right) \theta^{-1}\left(\ell_{i}\right) \ell_{i}^{-1} T_{\ell_{i}}+\chi^{2}\left(\sigma_{\mathfrak{I}_{i}}\right) \ell_{i}^{-1}\right)
$$

Taking the limit over $r$, we thus obtain a $\mathbb{T}(\mathfrak{a})^{\circ}$-linear map

$$
\begin{equation*}
\mathbb{T}(\mathfrak{a})^{\circ} \otimes_{\left(\mathbb{T}_{N(\Sigma)}^{N-}\right)_{\mathfrak{n}}}\left(\mathfrak{J}_{\left.N(\Sigma)^{+}\right)}\right)_{\mathfrak{n}} \rightarrow \mathbb{T}(\mathfrak{a})^{\circ} \otimes_{\left(\mathbb{T}_{N(\mathfrak{a})}^{N-}\right)_{\mathfrak{m}}}\left(\mathfrak{J}_{N(\mathfrak{a})^{+}}\right)_{\mathfrak{m}} \tag{13}
\end{equation*}
$$

having an effect on the corresponding measures as stated in Theorem 3.10. Hence to conclude the proof it remains to show that (13) is an isomorphism.

By Theorem 3.5, both the source and the target of this map are free of rank one over $\mathbb{T}(\mathfrak{a})^{\circ}$, and as in [Emerton et al. 2006, p. 559] (using [Hida 1988, Theorem 9.4]), one is reduced to showing the injectivity of the dual map modulo $p$ :

$$
\begin{align*}
H^{0}\left(\widetilde{X}_{N(\mathfrak{a})^{+}, 1} ; \mathbb{F}_{p}\right)^{\text {ord }}[\mathfrak{m}] & \rightarrow\left(\mathbb{T}_{N(\mathfrak{a})}^{N^{-}} / \mathfrak{m}\right) \otimes_{\mathbb{T}_{N(\Sigma)}^{N^{-}} / \mathfrak{n}}\left(H^{0}\left(\tilde{X}_{N(\mathfrak{a})^{+}, 1} ; \mathbb{F}_{p}\right)^{\text {ord } \left.\left[\mathfrak{m}^{\prime}\right]\right)}\right. \\
& \rightarrow\left(\mathbb{T}_{N(\mathfrak{a})}^{N^{-}} / \mathfrak{m}\right) \otimes_{\mathbb{T}_{N(\Sigma)}^{N^{-}} / \mathfrak{n}}\left(H^{0}\left(\widetilde{X}_{N(\Sigma)^{+}, 1} ; \mathbb{F}_{p}\right)^{\left.\operatorname{ord}\left[\mathfrak{m}^{\prime}\right]\right)}\right. \\
& \rightarrow\left(\mathbb{T}_{N(\mathfrak{a})}^{N^{-}} / \mathfrak{m}\right) \otimes_{\mathbb{T}_{N(\Sigma)}^{N^{-}} / \mathfrak{n}}\left(H^{0}\left(\widetilde{X}_{N(\Sigma)^{+}, 1} ; \mathbb{F}_{p}\right)^{\text {ord }[\mathfrak{n}]) ;}\right. \tag{14}
\end{align*}
$$

or equivalently (by a version of [Emerton et al. 2006, Lemma 3.8.1]), to showing that the composite of the first two arrows in (14) is injective.

In turn, the latter injectivity follows from Lemma 3.11 below, where the notations are as follows:

- $M^{+}$is any positive integer with $\left(M^{+}, p N^{-}\right)=1$;
- $\ell \neq p$ is a prime;
- $n_{\ell}=1$ or 2 according to whether or not $\ell$ divides $M^{+}$;
- $N^{+}:=\ell^{n_{\ell}} M^{+}$;
and

$$
\begin{equation*}
\epsilon_{\ell}^{*}: H^{0}\left(\tilde{X}_{M^{+}, 1} ; \mathbb{F}_{p}\right)^{\text {ord }}[\mathfrak{m}] \rightarrow\left(\mathbb{T}_{M^{+} N^{-}}^{N^{-}} / \mathfrak{m}\right) \otimes_{\mathbb{U}_{N^{+} N^{-}}^{\prime}} / \mathfrak{m}^{\prime}\left(H^{0}\left(\tilde{X}_{N^{+}, 1} ; \mathbb{F}_{p}\right)^{\text {ord }}\left[\mathfrak{m}^{\prime}\right]\right) \tag{15}
\end{equation*}
$$

is the map defined by

$$
\epsilon_{\ell}^{*}:= \begin{cases}B_{\ell, 1}^{*}-B_{\ell, \ell}^{*} \ell^{-1} T_{\ell} & \text { if } n_{\ell}=1 \\ B_{\ell^{2}, 1}^{*}-B_{\ell^{2}, \ell^{*}}^{*} \ell^{-1} T_{\ell}+B_{\ell^{2}, \ell^{2}}^{*} \ell^{-1}\langle\ell\rangle_{N(\mathfrak{a}) p} & \text { if } n_{\ell}=2\end{cases}
$$

Lemma 3.11. The map (15) is injective.
Proof. As in the proof of the analogous result [Emerton et al. 2006, Lemma 3.8.2] in the modular curve case, it suffices to show the injectivity of the map

$$
\left(H^{0}\left(\tilde{X}_{M^{+}, 1} ; \mathbb{F}\right)^{\text {ord }}\left[\mathfrak{m}_{\mathbb{F}}\right]\right)^{n_{\ell}+1} \xrightarrow{\beta_{\ell}} H^{0}\left(\tilde{X}_{N^{+}, 1} ; \mathbb{F}\right)^{\text {ord }}\left[\mathfrak{m}_{\mathbb{F}}^{\prime}\right]
$$

defined by

$$
\beta_{\ell}:= \begin{cases}B_{\ell, 1}^{*} \pi_{1}+B_{\ell, \ell}^{*} \pi_{2} & \text { if } n_{\ell}=1 \\ B_{\ell^{2}, 1}^{*} \pi_{1}+B_{\ell^{2}, \ell^{*}}^{*} \pi_{2}+B_{\ell^{2}, \ell^{2}}^{*} \pi_{3} & \text { if } n_{\ell}=2\end{cases}
$$

But in our quaternionic setting the proof of this injectivity follows from [Skinner and Wiles 1999, Lemma 3.26] for $n_{\ell}=1$ and [loc.cit., Lemma 3.28] for $n_{\ell}=2$.

Applying inductively Lemma 3.11 to the primes in $\Sigma$, the proof of Theorem 3.10 follows.

3G. Analytic Iwasawa invariants. Upon the choice of an isomorphism

$$
\mathbb{Z}_{p} \llbracket \Gamma \rrbracket \simeq \mathbb{Z}_{p} \llbracket T \rrbracket
$$

we may regard the $p$-adic $L$-functions $L_{\Sigma}(\bar{\rho}, \mathfrak{a})$ and $L(\bar{\rho}, \mathfrak{a})$, as well as the Euler factor $E_{\Sigma}(\bar{\rho}, \mathfrak{a})$, as elements in $\mathbb{T}(\mathfrak{a})^{\circ} \llbracket T \rrbracket$. In this section we apply the main result of the preceding section to study the variation of the Iwasawa invariants attached to the anticyclotomic $p$-adic $L$-functions of $p$-ordinary modular forms.

For any power series $f(T) \in R \llbracket T \rrbracket$ with coefficients in a ring $R$, the content of $f(T)$ is defined to be the ideal $I(f(T)) \subseteq R$ generated by the coefficients of $f(T)$. If $\wp$ is a height one prime of $\mathbb{T}_{\Sigma}$ belonging to the branch $\mathbb{T}(\mathfrak{a})$ (in the sense that $\mathfrak{a}$ is the unique minimal prime of $\mathbb{T}_{\Sigma}$ contained in $\left.\wp\right)$, we denote by $L(\bar{\rho}, \mathfrak{a})(\wp)$ the element of $\mathcal{O}_{\wp} \llbracket \Gamma \rrbracket$ obtained from $L(\bar{\rho}, \mathfrak{a})$ via reduction modulo $\wp$. In particular, we note that $L(\bar{\rho}, \mathfrak{a})(\wp)$ has unit content if and only if its $\mu$-invariant vanishes.

Theorem 3.12. The following are equivalent:
(1) $\mu(L(\bar{\rho}, \mathfrak{a})(\wp))=0$ for some newform $f_{\wp}$ in $\mathcal{H}^{-}(\bar{\rho})$;
(2) $\mu(L(\bar{\rho}, \mathfrak{a})(\wp))=0$ for every newform $f_{\wp}$ in $\mathcal{H}^{-}(\bar{\rho})$;
(3) $L(\bar{\rho}, \mathfrak{a})$ has unit content for some irreducible component $\mathbb{T}(\mathfrak{a})$ of $\mathcal{H}^{-}(\bar{\rho})$;
(4) $L(\bar{\rho}, \mathfrak{a})$ has unit content for every irreducible component $\mathbb{T}(\mathfrak{a})$ of $\mathcal{H}^{-}(\bar{\rho})$.

Proof. The argument in [Emerton et al. 2006, Theorem 3.7.5] applies verbatim, replacing the appeal to [loc.cit., Corollary 3.6.3] by our Theorem 3.10 above.

When any of the conditions in Theorem 3.12 hold, we shall write

$$
\mu^{\mathrm{an}}(\bar{\rho})=0
$$

For a power series $f(T)$ with unit content and coefficients in a local ring $R$, the $\lambda$-invariant $\lambda(f(T))$ is defined to be the smallest degree in which $f(T)$ has a unit coefficient.

Theorem 3.13. Assume that $\mu^{\text {an }}(\bar{\rho})=0$.
(1) Let $\mathbb{T}(\mathfrak{a})$ be an irreducible component of $\mathcal{H}^{-}(\bar{\rho})$. As $\wp ~ v a r i e s ~ o v e r ~ t h e ~ a r i t h-~$ metic primes of $\mathbb{T}(\mathfrak{a})$, the $\lambda$-invariant $\lambda(L(\bar{\rho}, \mathfrak{a})(\wp))$ takes on a constant value, denoted $\lambda^{\text {an }}(\bar{\rho}, \mathfrak{a})$.
(2) For any two irreducible components $\mathbb{T}\left(\mathfrak{a}_{1}\right), \mathbb{T}\left(\mathfrak{a}_{2}\right)$ of $\mathcal{H}^{-}(\bar{\rho})$, we have that

$$
\lambda^{\mathrm{an}}\left(\bar{\rho}, \mathfrak{a}_{1}\right)-\lambda^{\mathrm{an}}\left(\bar{\rho}, \mathfrak{a}_{2}\right)=\sum_{\ell \neq p} e_{\ell}\left(\mathfrak{a}_{2}\right)-e_{\ell}\left(\mathfrak{a}_{1}\right)
$$

where $e_{\ell}(\mathfrak{a})=\lambda\left(E_{\ell}(\mathfrak{a})\right)$.
Proof. The first part follows immediately from the definitions. For the second part, the argument in [Emerton et al. 2006, Theorem 3.7.7] applies verbatim, replacing their appeal to [loc.cit., Cor. 3.6.3] by our Theorem 3.10 above.

By Theorem 3.12 and Theorem 3.13, the Iwasawa invariants of $L(\bar{\rho}, \mathfrak{a})(\wp)$ are well behaved as $\wp$ varies. However, for the applications of these results to the Iwasawa main conjecture it is of course necessary to relate $L(\bar{\rho}, \mathfrak{a})(\wp)$ to $p$-adic $L$-functions defined by the interpolation of special values of $L$-functions. This question was addressed in [Castella and Longo 2016], as we now recall.
Theorem 3.14. If $\wp$ is the arithmetic prime of $\mathbb{T}(\mathfrak{a})$ corresponding to a p-ordinary p-stabilized newform $f_{\wp}$ of weight $k \geqslant 2$ and trivial nebentypus, then

$$
L(\bar{\rho}, \mathfrak{a})(\wp)=L_{p}\left(f_{\wp} / K\right)
$$

where $L_{p}\left(f_{\wp} / K\right)$ is the p-adic L-function of Chida and Hsieh [2016]. In particular, if $\chi: \Gamma \rightarrow \mathbb{C}_{p}^{\times}$is the $p$-adic avatar of an anticyclotomic Hecke character of $K$
of infinity type $(m,-m)$ with $-k / 2<m<k / 2$, then $L(\bar{\rho}, \mathfrak{a})(\wp)$ interpolates the central critical values

$$
\frac{L\left(f_{\wp} / K, \chi, k / 2\right)}{\Omega_{f_{\wp}, N^{-}}}
$$

as $\chi$ varies, where $\Omega_{f_{\wp}, N^{-}}$is the complex Gross period (6).
Proof. This is a reformulation of the main result of [Castella and Longo 2016]. (Note that the constant $\lambda_{\wp} \in F_{\wp}^{\times}$in [Castella and Longo 2016, Theorem. 4.6] is not needed here, since the specialization map of [loc.cit., §3.1] is being replaced by the map $\left(\mathfrak{J}_{N^{+}}\right)_{\mathfrak{m}} \rightarrow\left(\mathfrak{J}_{N^{+}, r, k}\right)_{\mathfrak{m}_{r, k}}$ induced by the isomorphism (8), which preserves integrality.)
Corollary 3.15. Let $f_{1}, f_{2} \in \mathcal{H}^{-}(\bar{\rho})$ be newforms with trivial nebentypus lying in the branches $\mathbb{T}\left(\mathfrak{a}_{1}\right), \mathbb{T}\left(\mathfrak{a}_{2}\right)$, respectively. Then $\mu^{\text {an }}(\bar{\rho})=0$ and

$$
\lambda\left(L_{p}\left(f_{1} / K\right)\right)-\lambda\left(L_{p}\left(f_{2} / K\right)\right)=\sum_{\ell \neq p} e_{\ell}\left(\mathfrak{a}_{2}\right)-e_{\ell}\left(\mathfrak{a}_{1}\right)
$$

where $e_{\ell}\left(\mathfrak{a}_{j}\right)=\lambda\left(E_{\ell}\left(\mathfrak{a}_{j}\right)\right)$.
Proof. By [Chida and Hsieh 2016, Theorem. 5.7] (extending Vatsal's result [2003] to higher weights), if $f \in \mathcal{H}^{-}(\bar{\rho})$ has weight $k \leqslant p+1$ and trivial nebentypus, then $\mu\left(L_{p}(f / K)\right)=0$. By Theorems 3.12 and 3.14, this implies $\mu^{\text {an }}(\bar{\rho})=0$. The result thus follows from Theorem 3.13, using Theorem 3.14 again to replace $\lambda^{\text {an }}\left(\bar{\rho}, \mathfrak{a}_{j}\right)$ by $\lambda\left(L_{p}\left(f_{j} / K\right)\right)$.

## 4. Anticyclotomic Selmer groups

We keep the notation of the previous sections. In particular, $\bar{\rho}: G_{\mathbb{Q}} \rightarrow \mathrm{GL}_{2}(\mathbb{F})$ is an odd irreducible Galois representation satisfying Assumption (SU) and isomorphic to $\bar{\rho}_{f_{0}}$ for some newform $f_{0}$ of weight $2, \mathcal{H}^{-}(\bar{\rho})$ is the associated Hida family, and $\Sigma$ is a finite set of primes split in the imaginary quadratic field $K$.

For each $f \in \mathcal{H}^{-}(\bar{\rho})$, let $V_{f}$ denote the self-dual Tate twist of the $p$-adic Galois representation associated to $f$, fix an $\mathcal{O}$-stable lattice $T_{f} \subseteq V_{f}$, and set $A_{f}:=$ $V_{f} / T_{f}$. Since $f$ is $p$-ordinary, there is a unique one-dimensional $G_{\mathbb{Q}_{p}}$-invariant subspace $F_{p}^{+} V_{f} \subseteq V_{f}$ where the inertia group at $p$ acts via $\varepsilon_{\text {cyc }}^{k / 2} \psi$, with $\psi$ a finite order character. Let $F_{p}^{+} A_{f}$ be the image of $F_{p}^{+} V_{f}$ in $A_{f}$, and as recalled in the Introduction define the minimal Selmer group of $f$ by
$\operatorname{Sel}\left(K_{\infty}, f\right):=\operatorname{ker}\left\{H^{1}\left(K_{\infty}, A_{f}\right) \rightarrow \prod_{w \nmid p} H^{1}\left(K_{\infty, w}, A_{f}\right) \times \prod_{w \mid p} H^{1}\left(K_{\infty, w}, F_{p}^{-} A_{f}\right)\right\}$,
where $w$ runs over the places of $K_{\infty}$ and we set $F_{p}^{-} A_{f}:=A_{f} / F_{p}^{+} A_{f}$.
It is well known that $\operatorname{Sel}\left(K_{\infty}, f\right)$ is cofinitely generated over $\Lambda$. When it is also $\Lambda$-cotorsion, we define the $\mu$-invariant $\mu\left(\operatorname{Sel}\left(K_{\infty}, f\right)\right)$ (resp. $\lambda$-invariant
$\left.\lambda\left(\operatorname{Sel}\left(K_{\infty}, f\right)\right)\right)$ to the largest power of $\varpi$ dividing (resp. the number of zeros of ) the characteristic power series of the Pontryagin dual of $\operatorname{Sel}\left(K_{\infty}, f\right)$.

A distinguishing feature of the anticyclotomic setting (in comparison with cyclotomic Iwasawa theory) is the presence of primes which split infinitely in the corresponding $\mathbb{Z}_{p}$-extension. Indeed, being inert in $K$, all primes $\ell \mid N^{-}$are infinitely split in $K_{\infty} / K$. As a result, the above Selmer group differs in general from the Greenberg Selmer group of $f$, defined as

$$
\mathfrak{S e l}\left(K_{\infty}, f\right):=\operatorname{ker}\left\{H^{1}\left(K_{\infty}, A_{f}\right) \rightarrow \prod_{w \nmid p} H^{1}\left(I_{\infty, w}, A_{f}\right) \times \prod_{w \mid p} H^{1}\left(K_{\infty, w}, F_{p}^{-} A_{f}\right)\right\}
$$

where $I_{\infty, w} \subseteq G_{K_{\infty}}$ denotes the inertia group at $w$.
If $S$ is a finite set of primes in $K$, we let $\operatorname{Sel}^{S}\left(K_{\infty}, f\right)$ and $\mathfrak{S e l}^{S}\left(K_{\infty}, f\right)$ be the " $S$-primitive" Selmer groups defined as above by omitting the local conditions at the primes in $S$ (except those above $p$, when any such prime is in $S$ ). Moreover, if $S$ consists of the primes dividing a rational integer $M$, we replace the superscript $S$ by $M$ in the above notation.

Immediately from the definitions, we see that there is as exact sequence

$$
\begin{equation*}
0 \rightarrow \operatorname{Sel}\left(K_{\infty}, f\right) \rightarrow \mathfrak{S c l}\left(K_{\infty}, f\right) \rightarrow \prod_{\ell \mid N^{-}} \mathcal{H}_{\ell}^{\mathrm{un}} \tag{16}
\end{equation*}
$$

where

$$
\mathcal{H}_{\ell}^{\mathrm{un}}:=\operatorname{ker}\left\{\prod_{w \mid \ell} H^{1}\left(K_{\infty, w}, A_{f}\right) \rightarrow \prod_{w \mid \ell} H^{1}\left(I_{\infty, w}, A_{f}\right)\right\}
$$

is the set of unramified cocycles. In [Pollack and Weston 2011, §§3, 5], Pollack and Weston carried out a careful analysis of the difference between $\operatorname{Sel}\left(K_{\infty}, f\right)$ and $\mathfrak{S e l}\left(K_{\infty}, f\right)$. Even though [loc. cit.] is mostly concerned with cases in which $f$ is of weight 2 , many of their arguments apply more generally. In fact, the next result follows essentially from their work.

Theorem 4.1. Assume that $\bar{\rho}$ satisfies Hypotheses (SU). Then the following are equivalent:
(1) $\operatorname{Sel}\left(K_{\infty}, f_{0}\right)$ is $\Lambda$-cotorsion with $\mu$-invariant zero for some newform $f_{0} \in$ $\mathcal{H}^{-}(\bar{\rho})$;
(2) $\operatorname{Sel}\left(K_{\infty}, f\right)$ is $\Lambda$-cotorsion with $\mu$-invariant zero for all newforms $f \in \mathcal{H}^{-}(\bar{\rho})$;
(3) $\mathfrak{S e l}\left(K_{\infty}, f\right)$ is $\Lambda$-cotorsion with $\mu$-invariant zero for all newforms $f \in \mathcal{H}^{-}(\bar{\rho})$. Moreover, in that case $\operatorname{Sel}\left(K_{\infty}, f\right) \simeq \operatorname{Sel}\left(K_{\infty}, f\right)$.

Proof. Assume $f_{0}$ is a newform in $\mathcal{H}^{-}(\bar{\rho})$ for which $\operatorname{Sel}\left(K_{\infty}, f_{0}\right)$ is $\Lambda$-cotorsion with $\mu$-invariant zero, and set $N^{+}:=N(\Sigma) / N^{-}$. By [Pollack and Weston 2011,

Proposition 5.1], we then have the exact sequences

$$
\begin{align*}
& 0 \rightarrow \operatorname{Sel}\left(K_{\infty}, f_{0}\right) \rightarrow \operatorname{Sel}^{N^{+}}\left(K_{\infty}, f_{0}\right) \rightarrow \prod_{\ell \mid N^{+}} \mathcal{H}_{\ell} \rightarrow 0  \tag{17}\\
& 0 \rightarrow \operatorname{Sel}\left(K_{\infty}, f_{0}\right) \rightarrow \mathfrak{S e l}^{N^{+}}\left(K_{\infty}, f_{0}\right) \rightarrow \prod_{\ell \mid N^{+}} \mathcal{H}_{\ell} \rightarrow 0, \tag{18}
\end{align*}
$$

where $\mathcal{H}_{\ell}$ is the product of $H^{1}\left(K_{\infty, w}, A_{f_{0}}\right)$ over the places $w \mid \ell$ in $K_{\infty}$. Since every prime $\ell \mid N^{+}$splits in $K$ (see Remark 3.9), the $\Lambda$-cotorsionness and the vanishing of the $\mu$-invariant of $\mathcal{H}_{\ell}$ can be deduced from [Greenberg and Vatsal 2000, Proposition 2.4]. Since $\operatorname{Sel}\left(K_{\infty}, f_{0}\right)$ [ $\left.\varpi\right]$ is finite by assumption, it thus follows from (17) that $\mathrm{Sel}^{N^{+}}\left(K_{\infty}, f_{0}\right)[\varpi]$ is finite. Combined with (16) and [Pollack and Weston 2011, Corollary 5.2], the same argument using (18) shows that then $\mathfrak{S e l}^{N^{+}}\left(K_{\infty}, f_{0}\right)[\varpi]$ is also finite.

On the other hand, following the arguments in the proof [Pollack and Weston 2011, Proposition 3.6] we see that for any $f \in \mathcal{H}(\bar{\rho})$ we have the isomorphisms

$$
\operatorname{Sel}^{N^{+}}\left(K_{\infty}, \bar{\rho}\right) \simeq \operatorname{Sel}^{N^{+}}\left(K_{\infty}, f\right)[\varpi] ; \quad \mathfrak{S e l}^{N^{+}}\left(K_{\infty}, \bar{\rho}\right) \simeq \mathfrak{S e l}^{N^{+}}\left(K_{\infty}, f\right)[\varpi]
$$

As a result, the argument in the previous paragraph implies that, for any newform $f \in \mathcal{H}^{-}(\bar{\rho})$, both $\operatorname{Sel}^{N^{+}}\left(K_{\infty}, f\right)[\varpi]$ and $\mathfrak{S e l}^{N^{+}}\left(K_{\infty}, f\right)[\varpi]$ are finite, from where (using (17) and (18) with $f$ in place of $f_{0}$ ) the $\Lambda$-cotorsionness and the vanishing of both the $\mu$-invariant of $\operatorname{Sel}\left(K_{\infty}, f\right)$ and of $\mathfrak{S e l}\left(K_{\infty}, f\right)$ follows. In view of (16) and [Pollack and Weston 2011, Lemma 3.4], the result follows.

Let $w$ be a prime of $K_{\infty}$ above $\ell \neq p$ and denote by $G_{w} \subseteq G_{K_{\infty}}$ its decomposition group. Let $\mathbb{T}(\mathfrak{a})$ be the irreducible component of $\mathbb{T}_{\Sigma}$ passing through $f$, and define

$$
\delta_{w}(\mathfrak{a}):=\operatorname{dim}_{\mathbb{F}} A_{f}^{G_{w}} / \varpi .
$$

(Note that this is well defined by [Emerton et al. 2006, Lemma 4.3.1].) Assume $\ell=\overline{\mathfrak{l}}$ splits in $K$ and put

$$
\begin{equation*}
\delta_{\ell}(\mathfrak{a}):=\sum_{w \mid \ell} \delta_{w}(\mathfrak{a}), \tag{19}
\end{equation*}
$$

where the sum is over the (finitely many) primes $w$ of $K_{\infty}$ above $\ell$.
In view of Theorem 4.1, we write $\mu^{\text {alg }}(\bar{\rho})=0$ whenever any of the $\mu$-invariants appearing in that result vanish. In that case, for any newform $f$ in $\mathcal{H}^{-}(\bar{\rho})$ we may consider the $\lambda$-invariants $\lambda\left(\operatorname{Sel}\left(K_{\infty}, f\right)\right)=\lambda\left(\operatorname{Sel}^{\prime}\left(K_{\infty}, f\right)\right)$.

Theorem 4.2. Let $\bar{\rho}$ and $\Sigma$ be as above, and assume that $\mu^{\text {alg }}(\bar{\rho})=0$. If $f_{1}$ and $f_{2}$ are any two newforms in $\mathcal{H}^{-}(\bar{\rho})$ lying in the branches $\mathbb{T}\left(\mathfrak{a}_{1}\right)$ and $\mathbb{T}\left(\mathfrak{a}_{2}\right)$, respectively, then

$$
\lambda\left(\operatorname{Sel}\left(K_{\infty}, f_{1}\right)\right)-\lambda\left(\operatorname{Sel}\left(K_{\infty}, f_{2}\right)\right)=\sum_{\ell \neq p} \delta_{\ell}\left(\mathfrak{a}_{1}\right)-\delta_{\ell}\left(\mathfrak{a}_{2}\right)
$$

Proof. Since we have the divisibilities $N^{-}\left|N\left(\mathfrak{a}_{i}\right)\right| N(\Sigma)$ with the quotient $N(\Sigma) / N^{-}$ only divisible by primes that are split in $K$, the arguments of [Emerton et al. 2006, §4] apply verbatim (cf., [Pollack and Weston 2011, Theorem 7.1]).

## 5. Applications to the main conjecture

5A. Variation of anticyclotomic Iwasawa invariants. Recall the definition of the analytic invariant $e_{\ell}(\mathfrak{a})=\lambda\left(E_{\ell}(\mathfrak{a})\right)$, where $E_{\ell}(\mathfrak{a})$ is the Euler factor from Section 3 F , and of the algebraic invariant $\delta_{\ell}(\mathfrak{a})$ introduced in (19).

Lemma 5.1. Let $\mathfrak{a}_{1}, \mathfrak{a}_{2}$ be minimal primes of $\mathbb{T}_{\Sigma}$. For any prime $\ell \neq p$ split in $K$,

$$
\delta_{\ell}\left(\mathfrak{a}_{1}\right)-\delta_{\ell}\left(\mathfrak{a}_{2}\right)=e_{\ell}\left(\mathfrak{a}_{2}\right)-e_{\ell}\left(\mathfrak{a}_{1}\right) .
$$

Proof. Let $\mathfrak{a}$ be a minimal prime of $\mathbb{T}_{\Sigma}$, let $f$ be a newform in the branch $\mathbb{T}(\mathfrak{a})$, and let $\wp_{f} \subseteq \mathfrak{a}$ be the corresponding height one prime. Since $\ell=\overline{\mathfrak{l}}$ splits in $K$, we have

$$
\underset{w \mid \ell}{\bigoplus_{\ell}} H^{1}\left(K_{\infty, w}, A_{f}\right)=\left(\bigoplus_{w \mid \mathfrak{l}} H^{1}\left(K_{\infty, w}, A_{f}\right)\right) \oplus\left(\underset{w \mid \overline{\mathfrak{l}}}{\bigoplus^{1}} H^{1}\left(K_{\infty, w}, A_{f}\right)\right)
$$

and [Greenberg and Vatsal 2000, Proposition 2.4] immediately implies that

$$
C h_{\Lambda}\left(\bigoplus_{w \mid \ell} H^{1}\left(K_{\infty, w}, A_{f}\right)^{\vee}\right)=E_{\ell}\left(f, \ell^{-1} \gamma_{\imath}\right) \cdot E_{\ell}\left(f, \ell^{-1} \gamma_{\uparrow}\right)
$$

where $E_{\ell}\left(f, \ell^{-1} \gamma_{\mathfrak{l}}\right) \cdot E_{\ell}\left(f, \ell^{-1} \gamma_{\uparrow}\right)$ is the specialization of $E_{\ell}(\mathfrak{a})$ at $\wp_{f}$. The result thus follows from [Emerton et al. 2006, Lemma 5.1.5].

Theorem 5.2. Suppose that $\bar{\rho}$ satisfies Assumption (SU). If for some newform $f_{0} \in \mathcal{H}^{-}(\bar{\rho})$ we have the equalities

$$
\mu\left(\operatorname{Sel}\left(K_{\infty}, f_{0}\right)\right)=\mu\left(L_{p}\left(f_{0} / K\right)\right)=0 \quad \text { and } \quad \lambda\left(\operatorname{Sel}\left(K_{\infty}, f_{0}\right)\right)=\lambda\left(L_{p}\left(f_{0} / K\right)\right)
$$

then the equalities

$$
\mu\left(\operatorname{Sel}\left(K_{\infty}, f\right)\right)=\mu\left(L_{p}(f / K)\right)=0 \quad \text { and } \quad \lambda\left(\operatorname{Sel}\left(K_{\infty}, f\right)\right)=\lambda\left(L_{p}(f / K)\right)
$$

hold for all newforms $f \in \mathcal{H}^{-}(\bar{\rho})$.
Proof. Let $f$ be any newform in $\mathcal{H}^{-}(\bar{\rho})$. Since the algebraic and analytic $\mu$ invariants of $f_{0}$ both vanish, the vanishing of $\mu\left(\operatorname{Sel}\left(K_{\infty}, f\right)\right)$ and $\mu\left(L_{p}(f / K)\right)$ follows from Theorems 4.1 and 3.12, respectively. On the other hand, combining Theorems 3.13 and 4.2, and Lemma 5.1, we see that

$$
\lambda\left(\operatorname{Sel}\left(K_{\infty}, f\right)\right)-\lambda\left(\operatorname{Sel}\left(K_{\infty}, f_{0}\right)\right)=\lambda\left(L_{p}(f / K)\right)-\lambda\left(L_{p}\left(f_{0} / K\right)\right)
$$

and hence the equality $\lambda\left(\operatorname{Sel}\left(K_{\infty}, f_{0}\right)\right)=\lambda\left(L_{p}\left(f_{0} / K\right)\right)$ implies the same equality for $f$.

5B. Applications to the main conjecture. As an immediate consequence of Weierstrass preparation theorem, Theorem 5.2 together with one the divisibilities predicted by the anticyclotomic main conjecture implies the full anticyclotomic main conjecture.

Theorem 5.3 (Skinner-Urban). Let $f \in S_{k}\left(\Gamma_{0}(N)\right)$ be a newform of weight $k \equiv$ $2(\bmod p-1)$ and trivial nebentypus. Suppose that $\bar{\rho}_{f}$ satisfies Assumption (SU) and that $p$ splits in $K$. Then

$$
\left(L_{p}(f / K)\right) \supseteq C h_{\Lambda}\left(\operatorname{Sel}\left(K_{\infty}, f\right)^{\vee}\right)
$$

Proof. This follows from specializing the divisibility in [Skinner and Urban 2014, Theorem 3.26] to the anticyclotomic line. Indeed, let $\boldsymbol{f}=\sum_{n \geqslant 1} \boldsymbol{a}_{n}(\boldsymbol{f}) q^{n} \in \llbracket \llbracket \rrbracket$ be the $\Lambda$-adic form with coefficients in $\rrbracket:=\mathbb{T}(\mathfrak{a})^{\circ}$ associated with the branch of the Hida family containing $f$, let $\Sigma$ be a finite set of primes as in Section 3E, let $\Sigma^{\prime} \supseteq \Sigma$ be a finite set of primes of $K$ containing $\Sigma$ and all primes dividing $p N(\mathfrak{a}) D_{K}$, and assume that $\Sigma^{\prime}$ contains at least one prime $\ell \neq p$ that splits in $K$. Under these assumptions, in [Skinner and Urban 2014, Theorem 3.26] it is shown that

$$
\begin{equation*}
\left(\mathfrak{L}_{p}^{\Sigma^{\prime}}(\boldsymbol{f} / K)\right) \supseteq C h_{\Lambda_{f}\left(L_{\infty}\right)}\left(\mathfrak{S e l}^{\Sigma^{\prime}}\left(L_{\infty}, A_{\boldsymbol{f}}\right)^{\vee}\right) \tag{20}
\end{equation*}
$$

where $L_{\infty}=K_{\infty} K_{\text {cyc }}$ is the $\mathbb{Z}_{p}^{2}$-extension of $K, \Lambda_{f}\left(L_{\infty}\right)$ is the three-variable Iwasawa algebra $\square\left[\operatorname{Gal}\left(L_{\infty} / K\right) \rrbracket\right.$, and $\mathfrak{L}_{p}^{\Sigma^{\prime}}(\boldsymbol{f} / K)$ and $\mathfrak{S e l}^{\Sigma^{\prime}}\left(L_{\infty}, A_{\boldsymbol{f}}\right)$ are the " $\Sigma^{\prime}$ primitive" $p$-adic $L$-function and Selmer group defined in [Skinner and Urban 2014, $\S 3.4 .5$ ] and [Skinner and Urban 2014, §§3.1.3, 3.1.10], respectively.

Recall the character $\Theta: G_{\mathbb{Q}} \rightarrow \mathbb{Z}_{p} \llbracket 1+p \mathbb{Z}_{p} \rrbracket^{\times}$from Section 2C, regarded as a character on $\operatorname{Gal}\left(L_{\infty} / K\right)$, and let

$$
\operatorname{Tw}_{\Theta^{-1}}: \Lambda_{f}\left(L_{\infty}\right) \rightarrow \Lambda_{f}\left(L_{\infty}\right)
$$

be the $\square$-linear isomorphism induced by $\operatorname{Tw}_{\Theta^{-1}}(g)=\Theta^{-1}(g) g$ for $g \in \operatorname{Gal}\left(L_{\infty} / K\right)$. Choose a topological generator $\gamma \in \operatorname{Gal}\left(K_{\text {cyc }} / K\right)$, and expand

$$
\operatorname{Tw}_{\Theta^{-1}}\left(\mathfrak{L}_{p}^{\Sigma^{\prime}}(\boldsymbol{f} / K)\right)=\mathfrak{L}_{p, 0}^{\Sigma^{\prime}}(\boldsymbol{f} / K)+\mathfrak{L}_{p, 1}^{\Sigma^{\prime}}(\boldsymbol{f} / K)(\gamma-1)+\cdots
$$

with $\mathfrak{L}_{p, i}^{\Sigma^{\prime}}(\boldsymbol{f} / K) \in \Lambda_{f}\left(K_{\infty}\right)=\square \llbracket \Gamma \rrbracket$. In particular, note that $\mathfrak{L}_{p, 0}^{\Sigma^{\prime}}(\boldsymbol{f} / K)$ is the restriction of the twisted three-variable $p$-adic $L$-function $\operatorname{Tw}_{\Theta^{-1}}\left(\mathfrak{L}_{p}^{\Sigma^{\prime}}(\boldsymbol{f} / K)\right)$ to the "self-dual" plane.

Because of our assumptions on $f$, the $\Lambda$-adic form $\boldsymbol{f}$ has trivial tame character, and hence denoting by $\mathrm{Frob}_{\ell}$ an arithmetic Frobenius at any prime $\ell \nmid N(\mathfrak{a}) p$, the Galois representation

$$
\rho(\mathfrak{a}): G_{\mathbb{Q}} \rightarrow \mathrm{GL}\left(T_{f}\right) \simeq \mathrm{GL}_{2}\left(\mathbb{T}(\mathfrak{a})^{\circ}\right)
$$

considered in Section 1D (which is easily seen to agree with the twisted representation considered in [Skinner and Urban 2014, p. 37]) satisfies

$$
\operatorname{det}\left(X-\operatorname{Frob}_{\ell} \mid T_{f}\right)=X^{2}-\boldsymbol{a}_{\ell}(\boldsymbol{f}) X+\Theta^{2}(\ell) \ell
$$

The twist $T_{f}^{\dagger}:=T_{f} \otimes \Theta^{-1}$ is therefore self-dual. Thus combining [Rubin 2000, Lemma 6.1.2] with a straightforward variant of [Skinner and Urban 2014, Proposition 3.9] having $\operatorname{Gal}\left(K_{\infty} / K\right)$ in place of $\operatorname{Gal}\left(K_{\text {cyc }} / K\right)$, we see that divisibility (20) implies that

$$
\begin{equation*}
\left(\mathfrak{L}_{p, 0}^{\Sigma^{\prime}}(\boldsymbol{f} / K)\right) \supseteq C h_{\Lambda_{f}\left(K_{\infty}\right)}\left(\mathfrak{S e l}^{\Sigma^{\prime}}\left(K_{\infty}, A_{f}^{\dagger}\right)^{\vee}\right) \tag{21}
\end{equation*}
$$

(Here, as above, $A_{f}$ denotes the Pontryagin dual $T_{f} \otimes_{\mathbb{1}} \operatorname{Hom}_{\mathrm{cts}}\left(\mathbb{\square}, \mathbb{Q}_{p} / \mathbb{Z}_{p}\right)$, and $A_{f}^{\dagger}$ is the corresponding twist.) We next claim that, setting $\Sigma^{\prime \prime}:=\Sigma^{\prime} \backslash \Sigma$, we have

$$
\begin{equation*}
\left(\mathfrak{L}_{p, 0}^{\Sigma^{\prime}}(\boldsymbol{f} / K)\right)=\left(L_{\Sigma}(\bar{\rho}, \mathfrak{a}) \cdot \prod_{v \in \Sigma^{\prime \prime}, v \nmid p} E_{v}(\mathfrak{a})\right), \tag{22}
\end{equation*}
$$

where $L_{\Sigma}(\bar{\rho}, \mathfrak{a})$ is the two-variable $p$-adic $L$-function constructed in Section 3D, and if $v$ lies over the rational prime $\ell, E_{v}(\mathfrak{a})$ is the Euler factor given by

$$
E_{v}(\mathfrak{a})=\operatorname{det}\left(\operatorname{Id}-\operatorname{Frob}_{v} X \mid\left(V_{f}^{\dagger}\right)_{I_{v}}\right)_{X=\ell^{-1} \mathrm{Frob}_{v}}
$$

where $V_{f}:=T_{f} \otimes_{\|} \operatorname{Frac}(\mathbb{\square})$, and $\operatorname{Frob}_{v}$ is an arithmetic Frobenius at $v$. (Note that for $\ell=\mathfrak{l} \overline{\mathfrak{l}}$ split in $K, E_{\mathfrak{l}}(\mathfrak{a}) \cdot E_{\overline{\mathfrak{l}}}(\mathfrak{a})$ is simply the Euler factor (11).) Indeed, combined with Theorems 3.10 and 3.14 , equality (22) specialized to any arithmetic prime $\wp \subseteq \mathbb{T}(\mathfrak{a})$ of weight 2 is shown in [Skinner and Urban 2014, (12.3)], from where the claim follows easily from the density of these primes. (See also [Pollack and Weston 2011, Theorem 6.8] for the comparison between the different periods involved in the two constructions, which differ by a $p$-adic unit under our assumptions.)

Finally, (21) and (22) combined with Theorem 3.10 and [Greenberg and Vatsal 2000, Propositions 2.3,8] imply that

$$
\left.(L(\bar{\rho}, \mathfrak{a})) \supseteq C h_{\Lambda_{f}\left(K_{\infty}\right)}\left(\mathfrak{S e l}^{( } K_{\infty}, A_{f}^{\dagger}\right)^{\vee}\right)
$$

from where the result follows by specializing at $\wp_{f}$ using Theorem 3.14 and Theorem 4.1.

In the opposite direction, we have the following result:
Theorem 5.4 (Bertolini-Darmon). Let $f=\sum_{n=1}^{\infty} a_{n}(f) q^{n}$ be a p-ordinary newform of weight 2, level N, and trivial nebentypus. Suppose that $\bar{\rho}_{f}$ satisfies Assumption (CR) and that

$$
\begin{equation*}
a_{p}(f) \not \equiv \pm 1(\bmod p) \tag{PO}
\end{equation*}
$$

Then

$$
\left(L_{p}(f / K)\right) \subseteq C h_{\Lambda}\left(\operatorname{Sel}\left(K_{\infty}, f\right)^{\vee}\right)
$$

Proof. This is the main result of [Bertolini and Darmon 2005], as extended by Pollack and Weston [Pollack and Weston 2011] to newforms of weight 2 not necessarily defined over $\mathbb{Q}$ and under the stated hypotheses (weaker that in [Bertolini and Darmon 2005]) on $\bar{\rho}_{f}$. See also [Kim et al. 2017] for a detailed discussion on the additional "nonanomalous" hypothesis ( PO ) on $f$.

Before we combine the previous two theorems with our main results in this paper, we note that condition (PO) in Theorem 5.4 can be phrased in terms of the Galois representation $\rho_{f}$ associated to $f$. Indeed, let $f=\sum_{n=1}^{\infty} a_{n}(f) q^{n}$ be a $p$-ordinary newform as above, defined over a finite extension $F / \mathbb{Q}_{p}$ with ring of integers $\mathcal{O}$. Then

$$
\left.\rho_{f}\right|_{D_{p}} \simeq\left(\begin{array}{ll}
\varepsilon & * \\
0 & \delta
\end{array}\right)
$$

on a decomposition group $D_{p} \subseteq G_{\mathbb{Q}}$ at $p$, with $\delta: D_{p} \rightarrow \mathcal{O}^{\times}$an unramified character sending $\operatorname{Frob}_{p}$ to the unit root $\alpha_{p}$ of $X^{2}-a_{p}(f) X+p$. Since clearly $\alpha \equiv a_{p}(f)(\bmod p)$, we see that condition (PO) amounts to the requirement that

$$
\begin{equation*}
\delta\left(\operatorname{Frob}_{p}\right) \not \equiv \pm 1(\bmod p) \tag{PO}
\end{equation*}
$$

Now we are finally in a position to prove our main application to the anticyclotomic Iwasawa main conjecture for $p$-ordinary newforms.

Corollary 5.5. Suppose that $\bar{\rho}$ satisfies Assumptions (SU) and (PO) and that $p$ splits in $K$, and let $f$ be a newform in $\mathcal{H}^{-}(\bar{\rho})$ of weight $k \equiv 2(\bmod p-1)$ and trivial nebentypus. Then the anticyclotomic Iwasawa main conjecture holds for $f$.

Proof. After Theorems 5.2 and 5.3, to check the anticyclotomic main conjecture for any newform $f$ as in the statement, it suffices to check the three equalities

$$
\begin{equation*}
\mu\left(\operatorname{Sel}\left(K_{\infty}, f_{0}\right)\right)=\mu\left(L_{p}\left(f_{0} / K\right)\right)=0, \quad \lambda\left(\operatorname{Sel}\left(K_{\infty}, f_{0}\right)\right)=\lambda\left(L_{p}\left(f_{0} / K\right)\right) \tag{23}
\end{equation*}
$$

hold for some $f_{0} \in \mathcal{H}^{-}(\bar{\rho})$ of weight $k \equiv 2(\bmod p-1)$ and trivial nebentypus.
Let $\mathbb{T}(\mathfrak{a})$ be the irreducible component of $\mathcal{H}^{-}(\bar{\rho})$ containing $f$, and let $f_{0} \in$ $S_{2}\left(\Gamma_{0}(N p)\right)$ be the $p$-stabilized newform corresponding to an arithmetic prime $\wp \subseteq \mathbb{T}(\mathfrak{a})$ of weight 2 and trivial nebentypus. By Assumption (PO), the form $f_{0}$ is necessarily the $p$-stabilization of a $p$-ordinary newform $f_{0}^{\sharp} \in S_{2}\left(\Gamma_{0}(N)\right.$ ) (see, e.g., [Howard 2007, Lemma 2.1.5]). From the combination of Theorems 5.3 and 5.4, the anticyclotomic Iwasawa main conjecture holds for $f_{0}^{\sharp}$, and since we clearly have

$$
L_{p}\left(f_{0} / K\right)=L_{p}\left(f_{0}^{\sharp} / K\right) \quad \text { and } \quad \operatorname{Sel}\left(K_{\infty}, f_{0}\right) \simeq \operatorname{Sel}\left(K_{\infty}, f_{0}^{\sharp}\right)
$$

(note that the latter isomorphism relies on the absolute irreducibility of $\bar{\rho}$ ), the anticyclotomic Iwasawa main conjecture holds for $f_{0}$ as well. In particular, equalities (23) hold for this $f_{0}$, and the result follows.

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# Effective nonvanishing for Fano weighted complete intersections 

Marco Pizzato, Taro Sano and Luca Tasin


#### Abstract

We show that the Ambro-Kawamata nonvanishing conjecture holds true for a quasismooth WCI $X$ which is Fano or Calabi-Yau, i.e., we prove that, if $H$ is an ample Cartier divisor on $X$, then $|H|$ is not empty. If $X$ is smooth, we further show that the general element of $|H|$ is smooth. We then verify the AmbroKawamata conjecture for any quasismooth weighted hypersurface. We also verify Fujita's freeness conjecture for a Gorenstein quasismooth weighted hypersurface.

For the proofs, we introduce the arithmetic notion of regular pairs and highlight some interesting connections with the Frobenius coin problem.


## 1. Introduction

Complete intersections in weighted projective spaces (WCIs for short) form a natural class of varieties which are particularly interesting from the point of view of higher dimensional algebraic geometry. We refer to [Dolgachev 1982], [Mori 1975] and [Dimca 1986] for a general treatment of these varieties.

Reid [1980; 1987] and Iano-Fletcher [2000] systematically investigated notable examples of WCIs and started their classification. Several results have since been obtained concerning boundedness and classification; see for example [Johnson and Kollár 2001; Chen et al. 2011; Ballico et al. 2013; Chen 2015; Przyjalkowski and Shramov 2016].

The main motivation of this paper is to study the following conjecture in the realm of WCIs, in particular for what concerns the case of Fano and Calabi-Yau varieties.

Conjecture 1.1 (Ambro-Kawamata). Let $(X, \Delta)$ be a klt pair and $H$ be an ample Cartier divisor on $X$ such that $H-K_{X}-\Delta$ is ample. Then $|H| \neq \varnothing$.

For an introduction to this conjecture, see [Ambro 1999] and in particular [Kawamata 2000] where the 2-dimensional case is proven. In the smooth setting, Ionescu [Lanteri et al. 1993, p. 321] and Beltrametti and Sommese [1995] proposed related conjectures.

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The Ambro-Kawamata conjecture is known to be true in full generality only in dimensions 1 and 2. Several cases have been studied, especially in dimension 3; see for instance [Xie 2009; Broustet and Höring 2010; Höring 2012; Cao and Jiang 2016].

A fundamental divisor on a Fano variety $X$ is an ample Cartier divisor $H$ which is primitive and proportional to $-K_{X}$. In the classification of Fano varieties, it is important to investigate the properties of the general member of the linear system given by $H$; see for instance [Ambro 1999]. The second purpose of this note is to study this problem in the case of Fano and Calabi-Yau smooth WCIs.

The main result of this paper is the following.
Theorem 1.2. Let $X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a well-formed quasismooth weighted complete intersection which is not a linear cone and $H$ be an ample Cartier divisor on $X$. Assume that $X$ is Fano or Calabi-Yau. Then $|H| \neq \varnothing$.

Moreover, if $X$ is smooth, then the number of $a_{i}=1$ is at least $c$ and the general element of $\left|\mathcal{O}_{X}(1)\right|$ is smooth.

For a smooth Fano WCI, it was already proved in [Przyjalkowski and Shramov 2016, Lemma 3.3] that at least two weights are 1, which implies the nonvanishing for a smooth Fano WCI. In addition, it is easy to prove Conjecture 1.1 for any smooth WCI of codimension 1 and 2, see Remark 4.9.

It is particularly interesting that, in the smooth case, we can prove the smoothness of the general member of the fundamental linear system (Corollary 5.3(ii)).

Theorem 1.2 is a direct consequence of Corollaries 5.3 and 5.13. In particular, in Corollary 5.3, we show that if $X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ is a smooth wellformed Fano WCI which is not a linear cone, then the number of $i$ for which $a_{i}=1$ is at least $c+1$. By using this, we can then show that the general element of $\left|\mathcal{O}_{X}(1)\right|$ is quasismooth (from which smoothness follows easily). One can not expect a similar statement for a general member of the fundamental linear system of a singular quasismooth WCI, as Example 5.8 shows. We also give a description of the base locus of $\left|\mathcal{O}_{X}(1)\right|$ in Remark 5.5 and an example whose base locus Bs $\left|\mathcal{O}_{X}(1)\right|$ is singular and not quasismooth in Example 5.6.

In [Przyjalkowski and Shramov 2017, Corollary 4.2], the authors show that for a smooth well-formed Fano WCI $X$ the number of $a_{i}$ equal to 1 is at least $I(X)=\sum a_{i}-\sum d_{j}$ when $c \leq 2$ and write that they expect this to hold for any codimension. As a consequence of Proposition 5.2, we can confirm this expectation; see Corollary 5.11.

In the case of hypersurfaces, we can prove the following stronger result, which is the combination of Propositions 6.2 and 6.3:

Theorem 1.3. Let $X=X_{d} \subset \mathbb{P}=\mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a well-formed quasismooth hypersurface of degree $d$ which is not a linear cone.
(1) If $H$ is an ample Cartier divisor on $X$ such that $H-K_{X}$ is ample, then $|H|$ is not empty.
(2) If $X$ is Gorenstein and $H$ is an ample Cartier divisor, then $K_{X}+m H$ is globally generated for any $m \geq n$.

The second part of the statement is known as Fujita's freeness conjecture and it has been proven in the smooth setting up to dimension 5; see [Reider 1988; Ein and Lazarsfeld 1993; Kawamata 1997; Ye and Zhu 2015].

The methods. The above theorems are obtained by studying the arithmetic properties of quasismooth WCIs. More precisely, in Section 3, we prove a criterion (Proposition 3.1) for a WCI to be quasismooth, which generalizes Iano-Fletcher's criterion in codimension 1 and 2 (see [Iano-Fletcher 2000, Section 8]). We then exploit some arithmetic consequences of quasismoothness. In particular, Proposition 3.6 motivates the introduction of an $h$-regular pair (see Definition 4.1) which turns out to be a key tool in our treatment.

Given a positive integer $h$, a pair $(d ; a)=\left(d_{1}, \ldots, d_{c} ; a_{0}, \ldots, a_{n}\right) \in \mathbb{N}^{c} \times \mathbb{N}^{n+1}$ is said to be $h$-regular if for any $I=\left\{i_{1}, \ldots, i_{k}\right\} \subset\{0, \ldots, n\}$ such that $a_{I}:=$ $\operatorname{gcd}\left(a_{i_{1}}, \ldots, a_{i_{k}}\right)>1$, either $a_{I} \mid h$ or there are distinct integers $p_{1}, \ldots, p_{k}$ such that

$$
a_{I} \mid d_{p_{1}}, \ldots, d_{p_{k}}
$$

Set $\delta(d ; a):=\sum_{j=1}^{c} d_{j}-\sum_{i=0}^{n} a_{i}$. By Proposition 3.6, any quasismooth (wellformed) WCI $X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ gives rise to an $h$-regular pair $(d ; a)=$ $\left(d_{1}, \ldots, d_{c} ; a_{0}, \ldots, a_{n}\right)$, where $h$ is the smallest positive integer for which $\mathcal{O}_{X}(h)$ is Cartier. Remembering that $K_{X}=\mathcal{O}_{X}(\delta)$, the nonvanishing for a Fano or Calabi-Yau WCI follows from Proposition 5.12, which says that, if $(d ; a)$ is $h$-regular such that $a_{i} \neq d_{j}$ and $a_{i} \nmid h$ for any $i, j$, then $\delta(d ; a)>0$. A more accurate statement (Corollary 5.3) is needed to prove that, if $X$ is smooth, then the general element of $\left|\mathcal{O}_{X}(1)\right|$ is also smooth.

We now spend some words for the case $h=1$. In this case, the pair $(d ; a)$ is simply called regular. A smooth WCI $X$ gives rise to a regular pair $(d ; a)$. The nonvanishing is then equivalent to prove that

$$
\delta(d ; a) \geq G\left(a_{0}, \ldots, a_{n}\right),
$$

where $G\left(a_{0}, \ldots, a_{n}\right)$ is the Frobenius number of $a_{0}, \ldots, a_{n}$, i.e., the greatest integer which is not a nonnegative integral combination of $a_{0}, \ldots, a_{n}$. In Conjecture 4.8, we speculate that $\delta(d ; a) \geq G\left(a_{0}, \ldots, a_{n}\right)$ for a regular pair $(d ; a)$, under some natural assumptions. This would imply the Ambro-Kawamata conjecture for any smooth WCI.

We believe that this conjecture is interesting also from the arithmetic point of view, since it would give new bounds for the Frobenius number (see page 2382 for details).

## 2. Preliminaries and notation

In this section, we recall some basic facts about weighted complete intersections and fix our notation. See [Dolgachev 1982] or [Iano-Fletcher 2000] for further details.

Let $\mathbb{N}$ (resp. $\mathbb{N}_{+}$) be the set of nonnegative (resp. positive) integers. Let $a_{0}, \ldots, a_{n} \in$ $\mathbb{N}_{+}$. We define $\mathbb{P}:=\mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ to be the weighted projective space with weights $a_{0}, \ldots, a_{n}$, i.e., $\mathbb{P}=\operatorname{Proj} \mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$, where $x_{i}$ has weight $a_{i}$. We denote

$$
\mathbb{P}(\underbrace{b_{1}, \ldots, b_{1}}_{k_{1}}, \ldots, \underbrace{b_{l}, \ldots, b_{l}}_{k_{l}})
$$

by $\mathbb{P}\left(b_{1}^{\left(k_{1}\right)}, \ldots, b_{l}^{\left(k_{l}\right)}\right)$ for short.
Note that if we start with $x_{0}, \ldots, x_{n}$ to be affine coordinates on $\mathbb{A}^{n+1}$ and $\mathbb{C}^{*}$ acting on $\mathbb{A}^{n+1}$ via

$$
\lambda \cdot\left(x_{0}, \ldots, x_{n}\right)=\left(\lambda^{a_{0}} x_{0}, \ldots, \lambda^{a_{n}} x_{n}\right)
$$

for any $\lambda \in \mathbb{C}^{*}$, then $\mathbb{P}$ is just the quotient $\left(\mathbb{A}^{n+1} \backslash\{0\}\right) / \mathbb{C}^{*}$.
We always assume that $\mathbb{P}$ is well-formed, i.e., the greatest common divisor of any $n$ weights is 1 . For any $I=\left\{i_{1}, \ldots, i_{k}\right\} \subset\{0, \ldots, n\}$, the stratum $\Pi_{I}$ is defined as

$$
\Pi_{I}:=\left\{x_{i}=0: i \notin I\right\}
$$

The singular locus of $\mathbb{P}$ is the union of all strata $\Pi_{I}$ for which $a_{I}:=\operatorname{gcd}\left(a_{i}\right)_{i \in I}>1$. Any point of the interior $\Pi_{I}^{0}$ of a stratum $\Pi_{I}$ is locally isomorphic to a quotient singularity of type

$$
\frac{1}{a_{I}}\left(a_{0}, \ldots, \hat{a}_{i_{1}}, \ldots, \hat{a}_{i_{k}}, \ldots, a_{n}\right) \times \mathbb{C}^{k-1}
$$

Here, for $r \in \mathbb{N}_{+}$and $a_{1}, \ldots, a_{n} \in \mathbb{N}$ such that $\operatorname{gcd}\left(r, a_{1}, \ldots, a_{n}\right)=1$, a quotient singularity of type $1 / r\left(a_{1}, \ldots, a_{n}\right)$ means a quotient $\mathbb{C}^{n} / \mathbb{Z}_{r}$ by the action of a cyclic group $\mathbb{Z}_{r}$ of order $r$ as $g \cdot z_{i}=\zeta_{r}^{a_{i}} z_{i}$ for $i=1, \ldots, n$, where $g \in \mathbb{Z}_{r}$ is a generator and $\zeta_{r}$ is an $r$-th primitive root of unity. We also denote by $\mathbb{C}^{n} / \mathbb{Z}_{r}\left(a_{1}, \ldots, a_{r}\right)$ this quotient affine variety. Let $\pi: \mathbb{C}^{n} \rightarrow U:=\mathbb{C}^{n} / \mathbb{Z}_{r}\left(a_{1}, \ldots, a_{r}\right)$ be the quotient morphism. We have an eigendecomposition

$$
\pi_{*} \mathcal{O}_{\mathbb{C}^{n}}=\bigoplus_{i=0}^{r-1} \mathcal{F}_{i}
$$

where $\mathcal{F}_{i}:=\left\{f \in \mathcal{O}_{\mathbb{C}^{n}} \mid g \cdot f=\zeta_{r}^{i} f\right\}$ is the $\mathcal{O}_{U}$-submodule of $\pi_{*} \mathcal{O}_{\mathbb{C}^{n}}$ consisting of $\mathbb{Z}_{r}$-eigenfunctions of eigenvalue $\zeta_{r}^{i}$. Note that $\mathcal{F}_{i} \simeq \mathcal{O}_{U}\left(D_{f}\right)$, where a divisor $D_{f}:=(f=0) / \mathbb{Z}_{r}$ on $U$ is defined by a function $f \in \mathcal{F}_{i}$.

Proposition 2.1. The divisor class group of $U:=\mathbb{C}^{n} / \mathbb{Z}_{r}\left(a_{1}, \ldots, a_{r}\right)$ is

$$
\begin{equation*}
\mathrm{Cl} U \simeq \mathbb{Z}_{r} \cdot \mathcal{F}_{1} \tag{1}
\end{equation*}
$$

Proof. We have an inclusion $\iota: \mathbb{Z}_{r} \cdot \mathcal{F}_{1} \hookrightarrow \mathrm{Cl} U$. It is enough to show that this is surjective. Let $D \subset U$ be a prime divisor. Then $\pi^{-1}(D)$ is a divisor on $\mathbb{C}^{n}$ defined by some $\mathbb{Z}_{r}$-eigenfunction $f_{D} \in \mathcal{O}_{\mathbb{C}^{n}}$. Let $i(D) \in \mathbb{Z}_{r}$ be an element such that $g \cdot f_{D}=\zeta_{r}^{i(D)} f$. Then we can check that $\mathcal{O}_{U}(D) \simeq \mathcal{F}_{i(D)}$. Since $g^{i} \cdot \mathcal{F}_{1} \simeq \mathcal{F}_{i}$ for $i \in \mathbb{Z}_{r}$, we see that $\iota$ is surjective.

We can also check the isomorphism by toric computation. Since $\mathbb{C}^{n} / \mathbb{Z}_{r}\left(a_{1}, \ldots, a_{n}\right)$ is a toric variety, we can compute its class group by using the information of the cone and lattice. (cf., [Fulton 1993, p. 63, Proposition]) More precisely, it is the quotient $\mathbb{Z}^{n} / M$, where $M:=\left\{\left(m_{1}, \ldots, m_{n}\right) \in \mathbb{Z}^{n} \mid \sum_{i=1}^{n} m_{i} a_{i} \equiv 0 \bmod r\right\}$.
Definition 2.2. Let $X$ be a (closed) subvariety of codimension $c$ in $\mathbb{P}$. Then $X$ is well-formed if

$$
\operatorname{codim}_{X}(X \cap \operatorname{Sing}(\mathbb{P})) \geq 2
$$

Let $\pi: \mathbb{A}^{n+1} \backslash\{0\} \rightarrow \mathbb{P}$ be the natural projection. Then $X$ is quasismooth if $\pi^{-1}(X)$ is smooth.

The variety $X$ is said to be a weighted complete intersection (WCI for short) of multidegree $\left(d_{1}, \ldots, d_{c}\right)$ if its weighted homogenous ideal in $\mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$ is generated by a regular sequence of homogenous polynomials $\left\{f_{j}\right\}$ such that $\operatorname{deg} f_{j}=d_{j}$ for $j=1, \ldots, c$. We denote by $X_{d_{1}, \ldots, d_{c}}$ a general element of the family of WCIs of multidegree $\left(d_{1}, \ldots, d_{c}\right)$.

Finally, $X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}$ is said to be a linear cone if $d_{j}=a_{i}$ for some $i$ and $j$.
Note that by [Dimca 1986, Proposition 8], if $X$ is a well-formed quasismooth WCI, then

$$
\operatorname{Sing}(X)=X \cap \operatorname{Sing}(\mathbb{P})
$$

Proposition 2.3. If $X$ is a quasismooth $W C I$ of dimension $\geq 3$, then its divisor class group is a free $\mathbb{Z}$-module generated by $\mathcal{O}_{X}(1)$, where $\mathcal{O}_{X}(1):=\mathcal{O}_{\mathbb{P}}(1)_{\mid X} .($ We freely mix the divisorial and the sheaf notation.)

Proof. The proof is the same as [Corti et al. 2000, Lemma 3.5]. This follows from the parafactoriality of an l.c.i. local ring [Call and Lyubeznik 1994].

If $X \subset \mathbb{P}$ is a well-formed quasismooth WCI, then

$$
\omega_{X}=K_{X}=\mathcal{O}_{X}\left(\sum_{j=1}^{c} d_{j}-\sum_{i=0}^{n} a_{i}\right)
$$

see [Dolgachev 1982, Theorem 3.3.4]. We usually write $\delta:=\sum_{j=1}^{c} d_{j}-\sum_{i=0}^{n} a_{i}$.
The following result shows that the dimension of the linear system $\left|\mathcal{O}_{X}(n)\right|$ can be computed by the weights of the coordinates.

Lemma 2.4 [Iano-Fletcher 2000, Lemma 7.1]. Let $X \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a wellformed quasismooth WCI. Let $A:=k\left[x_{0}, \ldots, x_{n}\right] /\left(f_{1}, \ldots, f_{c}\right)$ be the homogeneous coordinate ring of $X$ and $A_{k}$ be the $k$-th graded part for $k \in \mathbb{Z}$. Then

$$
H^{0}\left(X, \mathcal{O}_{X}(k)\right) \simeq A_{k}
$$

Proof. See also [Dolgachev 1982, 3.4.3]. This follows since the homogeneous coordinate ring $A$ is Cohen-Macaulay and $H_{\mathfrak{m}}^{1}(A)=0$, where $\mathfrak{m}:=\left(x_{0}, \ldots, x_{n}\right)$ is the maximal ideal.

## 3. Properties of quasismooth WCIs

In the following proposition, we give a necessary and sufficient condition for quasismoothness of a WCI.

Proposition 3.1. Let $X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)=: \mathbb{P}$ be a quasismooth WCI which is not a linear cone. Let $x_{0}, \ldots, x_{n}$ be the coordinates of $\mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$. Fix $I:=\left\{i_{1}, \ldots, i_{k}\right\} \subset\{0, \ldots, n\}$ and let $\rho_{I}:=\min \{c, k\}$. For $m=\left(m_{1}, \ldots, m_{k}\right)$, let $x_{I}^{m}:=\prod_{j=1}^{k} x_{i_{j}}^{m_{j}}$. For a finite set $A$, let $|A|$ be the number of its elements. Then one of the following holds.
(Q1) There exist distinct integers $p_{1}, \ldots, p_{\rho_{I}} \in\{1, \ldots, c\}$ and $M_{1}, \ldots, M_{\rho_{I}} \in \mathbb{N}^{k}$ such that the monomial $x_{I}^{M_{j}}$ has the degree $d_{p_{j}}$ for $j=1, \ldots, \rho_{I}$.
(Q2) There exist a permutation $p_{1}, \ldots, p_{c}$ of $\{1, \ldots, c\}$, an integer $l<\rho_{I}$, and integers $e_{\mu, j} \in\{0, \ldots, n\} \backslash_{M_{j}}$ for $\mu=1, \ldots, k-l$ and $j=l+1, \ldots$, c such that there are monomials $x_{I}^{M_{j}}$ of degree $d_{p_{j}}$ for $j=1, \ldots, l$ and distinct $k-l$ monomials $\left\{x_{e_{\mu, j}} x_{I}^{M_{\mu, j}}: \mu=1, \ldots, k-l\right\}$ of degree $d_{p_{j}}$ for each $j=l+1, \ldots, c$ which satisfy the following: for any subset $J \subset\{l+1, \ldots, c\}$, we have $\left|\left\{e_{\mu, j}: j \in J, \mu=1, \ldots, k-l\right\}\right| \geq k-l+|J|-1$.

Conversely, if we have (Q1) or (Q2) for all I, then a general WCI $X_{d_{1}, \ldots, d_{c}} \subset$ $\mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ is quasismooth.

Remark 3.2. This generalizes [Iano-Fletcher 2000, Theorem 8.7] in codimension 2 case. A weaker necessary condition for the quasismoothness is written in [Chen 2015, Proposition 2.3]. Although we shall not use the new part of Proposition 3.1 in the main part of this paper, we believe it is an interesting result on its own.
Proof. The framework of the proof is similar to that of [Iano-Fletcher 2000, Theorem 8.7].

Let $F_{j}:=\left|\mathcal{O}_{\mathbb{P}}\left(d_{j}\right)\right|$ be the linear system of weighted homogeneous polynomials of degree $d_{j}$. For $j=1, \ldots, c$, let $f_{j}$ be a general homogeneous polynomial of degree $d_{j}$ such that $X=\left(f_{1}=\cdots=f_{c}=0\right) \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$. Let $C_{X}^{*} \subset \mathbb{A}^{n+1} \backslash\{0\}$ be the cone over $X$ defined by the polynomials $f_{1}, \ldots, f_{c}$ with the following diagram


Without loss of generality, we may assume $I=\{0, \ldots, k-1\}$ in the statement. Let $\Pi:=\left(x_{k}=\cdots=x_{n}=0\right) \subset \mathbb{A}^{n+1}$ be the stratum corresponding to $I$ and $\Pi^{0} \subset \Pi$ be the open toric stratum. By expanding $f_{\lambda}$ for $\lambda=1, \ldots, c$ in terms of $x_{k}, \ldots, x_{n}$, we can write

$$
f_{\lambda}=h_{\lambda}\left(x_{0}, \ldots, x_{k-1}\right)+\sum_{i=k}^{n} x_{i} g_{\lambda}^{i}\left(x_{0}, \ldots, x_{k-1}\right)+R_{\lambda}\left(x_{0}, \ldots, x_{n}\right),
$$

where $h_{\lambda}, g_{\lambda}^{i} \in \mathbb{C}\left[x_{0}, \ldots, x_{k-1}\right]$ and $R_{\lambda} \in \mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$ satisfies $\operatorname{deg}_{x_{k}, \ldots, x_{n}} R_{\lambda} \geq 2$.
Note that $X$ is quasismooth if and only if $C_{X}^{*}$ is smooth along all the coordinate strata. We shall show that $C_{X}^{*}$ is smooth along $\Pi^{0}$ when either $(\mathrm{Q} 1)$ or $(\mathrm{Q} 2)$ holds for $I$. Let $\rho:=\rho_{I}$ for short.

Suppose that (Q1) holds. Then $h_{p_{1}}, \ldots, h_{p_{\rho}}$ are nonzero on $\Pi^{0}$. If some of $h_{p_{j}}$ involves only one monomial, then we have $\Pi^{0} \cap C_{X}^{*}=\varnothing$. So we may assume that each of $h_{p_{1}}, \ldots, h_{p_{\rho}}$ involves at least 2 monomials. Thus we see that the linear systems $F_{d_{p_{1}}}, \ldots, F_{d_{p_{\rho}}}$ do not have base locus on $\Pi^{0}$. By Bertini's theorem, we see that $\left(f_{p_{1}}=\cdots=f_{p_{\rho}}=0\right) \subset \mathbb{A}^{n+1}$ is smooth along $\Pi^{0}$ when $k \geq c$. When $k<c$, we have $\left(f_{p_{1}}=\cdots=f_{p_{\rho}}=0\right) \cap \Pi^{0}=\varnothing$. Therefore $C_{X}^{*}$ is nonsingular along $\Pi^{0}$.

Next suppose that (Q2) holds. By permutation, we may assume that $p_{i}=i$. Then $h_{1}, \ldots, h_{l}$ are nonzero on $\Pi^{0}$. Hence the base locus of $F_{d_{\lambda}}$ is disjoint with $\Pi^{0}$ for $\lambda=1, \ldots, l$. By Bertini's theorem, we see that ( $f_{1}=\cdots=f_{l}=0$ ) is nonsingular along $\Pi^{0}$. We may assume that the Jacobian of $\left(f_{1}=\cdots=f_{c}=0\right) \subset \mathbb{A}^{n+1}$ at $P \in \Pi^{0}$ is of the form

$$
\left(\begin{array}{ccccc}
\frac{\partial f_{1}}{\partial x_{0}} & \cdots \frac{\partial f_{1}}{\partial x_{k-1}} & & & \\
\vdots & \vdots & & * & \\
\frac{\partial f_{l}}{\partial x_{0}} & \cdots & \frac{\partial f_{l}}{\partial x_{k-1}} & & \\
\\
& & g_{l+1}^{k} & \cdots & g_{l+1}^{n} \\
0 & & \vdots & & \vdots \\
& & g_{c}^{k} & \cdots & g_{c}^{n}
\end{array}\right)(P)
$$

since we have $h_{\lambda}=0$ for $\lambda=l+1, \ldots, c$. Note that the block matrix

$$
\left(\begin{array}{ccc}
\frac{\partial f_{1}}{\partial x_{0}} \cdots & \frac{\partial f_{1}}{\partial x_{k-1}} \\
\vdots & & \vdots \\
\frac{\partial f_{l}}{\partial x_{0}} \cdots & \frac{\partial f_{l}}{\partial x_{k-1}}
\end{array}\right)(P)
$$

has maximal rank $l$ at $P \in \Pi^{0}$ since $\left(f_{1}=\cdots=f_{l}=0\right)$ is nonsingular along $\Pi^{0}$. Hence it is enough to show that the matrix

$$
M_{P}:=\left(\begin{array}{ccc}
g_{l+1}^{k} & \cdots & g_{l+1}^{n}  \tag{2}\\
\vdots & & \vdots \\
g_{c}^{k} & \cdots & g_{c}^{n}
\end{array}\right)(P)
$$

has maximal rank $c-l$.
Note that there are at least $k-l$ elements of $K_{\lambda}:=\left\{i \in\{k, \ldots, n\}: g_{\lambda}^{i} \neq 0\right\}$ for each $\lambda=l+1, \ldots, c$. By $\left|K_{\lambda}\right| \geq k-l$, we see that each row vector of $M_{P}$ is nonzero for $P \in \Pi^{0}$. Indeed, for each $\lambda^{\prime}=l+1, \ldots, c$, the intersection

$$
\bigcap_{\lambda=1}^{l}\left(h_{\lambda}=0\right) \cap \bigcap_{i=k}^{n}\left(g_{\lambda^{\prime}}^{i}=0\right) \cap \Pi^{0}
$$

is contained in at least $k=l+(k-l)$ free linear systems on $k$-dimensional $\Pi^{0}$, and it is empty. Thus we may assume that $g_{l+1}^{k}(P) \neq 0$. We shall make elementary matrix operations on $M_{P}$ to calculate the rank of $M_{P}$.

For $\lambda=l+2, \ldots, c$, let

$$
\begin{aligned}
Z_{\lambda}(P):=\left\{Q \in \left(f_{1}=\cdots=\right.\right. & \left.f_{l}=0\right) \cap \Pi^{0}: \\
& \left.g_{l+1}^{k}(P) g_{\lambda}^{i}(Q)-g_{\lambda}^{k}(P) g_{l+1}^{i}(Q)=0(i=k+1, \ldots, n)\right\}
\end{aligned}
$$

Note that the first row $M_{P}^{1}$ and the $(\lambda-l)$-th row $M_{P}^{\lambda-l}$ of $M_{P}$ are linearly dependent if and only if $P \in Z_{\lambda}(P)$. By condition (Q2) for $J$ with $|J|=2$, there are at least $k-l$ nonzero elements of $G_{\lambda}(P):=\left\{g_{l+1}^{k}(P) g_{\lambda}^{i}-g_{\lambda}^{k}(P) g_{l+1}^{i}: i=k+1, \ldots, n\right\}$ and they define $k-l$ free linear systems on $\Pi^{0}$. Hence we obtain $Z_{\lambda}(P)=\varnothing$ and the two rows $M_{P}^{1}$ and $M_{P}^{\lambda-l}$ are linearly independent. Thus, by elementary operations on $M_{P}$, we obtain a matrix of the following form;

$$
\left(\begin{array}{cccc}
g_{l+1}^{k} & \cdots & \cdots & g_{l+1}^{n} \\
0 & h_{l+2}^{k+1} & \cdots & h_{l+2}^{n} \\
\vdots & \vdots & & \vdots \\
0 & h_{c}^{k+1} & \cdots & h_{c}^{n}
\end{array}\right)(P) .
$$

By column exchange operations, we may assume that $h_{l+2}^{k+1}(P) \neq 0$ and repeat the process to

$$
M_{P}^{\prime}:=\left(\begin{array}{ccc}
h_{l+2}^{k+1} & \cdots & h_{l+2}^{n} \\
\vdots & & \vdots \\
h_{c}^{k+1} & \cdots & h_{c}^{n}
\end{array}\right)(P)
$$

Let $G_{\lambda}^{\prime}(P):=\left\{h_{l+2}^{k+1}(P) h_{\lambda}^{i}-h_{\lambda}^{k+1}(P) h_{l+2}^{i}: i=k+2, \ldots, n\right\}$. By condition (Q2) for $J$ with $|J|=3$, there are at least $k-l$ nonzero elements of $G_{\lambda}^{\prime}(P)$ and they
define free linear systems on $\Pi^{0}$. By this, we again see that the first row and another row of $M_{P}^{\prime}$ are linearly independent.

After repeating these elementary operations, we obtain a matrix of the form

$$
\left(\begin{array}{cccc}
\alpha_{l+1} & & \\
& \ddots & * & * \\
0 & & \alpha_{c}
\end{array}\right)
$$

for some $\alpha_{l+1}, \ldots, \alpha_{c} \in \mathbb{C} \backslash\{0\}$ and see that the rank of $M_{P}$ is $c-l$. Thus $C_{X}^{*}$ is nonsingular at $P \in \Pi^{0}$.

Suppose that conditions $(\mathrm{Q} 1)$ and $(\mathrm{Q} 2)$ do not hold for some $I$. We shall show that $X$ is not quasismooth. We may again assume that $I=\{0, \ldots, k-1\}$ and $\Pi=\left(x_{k}=\cdots=x_{n}=0\right)$. Moreover, since (Q1) and (Q2) do not hold, we may assume that, for some $l<\rho_{I}$, we have $\Pi \not \subset\left(f_{\lambda}=0\right)$ for $\lambda=1, \ldots, l$ and $\Pi \subset\left(f_{\lambda}=0\right)$ for $\lambda=l+1, \ldots, c$. Then the singular locus of $C_{X}^{*}$ on $\Pi^{0}$ can be described as

$$
Z:=\left\{P \in\left(f_{1}=\cdots=f_{l}=0\right) \cap \Pi^{0}: \operatorname{rk} M_{P}<c-l\right\}
$$

where $M_{P}$ is the matrix defined in (2). By the hypothesis, we may also assume that there exists $J \subset\{l+1, \ldots, c\}$ such that there are at most $k-l+|J|-2$ nonzero elements among $\left\{g_{\lambda}^{i}: \lambda \in J, i=k, \ldots, n\right\}$. This implies that there are at most $k-l+|J|-2$ nonzero columns of the matrix $M_{P}^{J}:=\left(g_{\lambda}^{i}(P)\right)_{\lambda \in J}^{k \leq i \leq n}$. We can choose $J$ so that the number $|J|$ is minimal among such subsets of $\{l+1, \ldots, c\}$. Then, by elementary operations as in the first part of the proof, we can transfer $M_{P}^{J}$ to the form

$$
\left(\begin{array}{cccccc}
h_{l+1}^{k} & \cdots & \cdots & h_{l+1}^{k+|J|} & \cdots & h_{l+1}^{n} \\
& \ddots & & \vdots & & \vdots \\
0 & & h_{l+|J|}^{k+|J|-1} & h_{l+|J|}^{k+|J|} & \cdots & h_{l+|J|}^{n}
\end{array}\right)(P)
$$

Note that on the bottom row we have at most $k-l-1$ nonzero entries. Hence we obtain

$$
\begin{aligned}
\operatorname{dim}\left(f_{1}=\cdots=f_{l}=0\right) \cap\left(h_{l+|J|}^{k+|J|-1}=\cdots=h_{l+|J|}^{n}=0\right) & \cap \Pi^{0} \\
& \geq k-l-(k-l-1)=1
\end{aligned}
$$

Since the rank of $M_{P}^{J}$ is not maximal on the subset $\left(h_{l+|J|}^{k+|J|-1}=\cdots=h_{l+|J|}^{n}=0\right)$, we see that $C_{X}^{*}$ is singular along the above positive dimensional subset. Hence $X$ is not quasismooth in this case. This concludes the proof of Proposition 3.1.

In the following example, we use Proposition 3.1 to check quasismoothness of a given WCI.
Example 3.3. Let $X_{8,8,8} \subset \mathbb{P}\left(2^{(4)}, 3^{(5)}, 5^{(3)}\right)$ be a general WCI of codimension 3. We can check the quasismoothness of $X_{8,8,8}$ by Proposition 3.1 as follows. Consider
$I=\{4,5,6,7,8\}$, that is, $a_{4}=\cdots=a_{8}=3$. Then (Q1) does not hold for this $I$ and we have $k=5, l=0$ in (Q2). We can choose $\left\{e_{\mu, j}: j=1,2,3, \mu=1, \ldots, 5\right\} \subset$ $\{0,1,2,3,9,10,11\}$ so that (Q2) is satisfied for this $I$. We can similarly check that (Q2) holds for $I=\{9,10,11\}$. For other $I$, we have (Q1), thus we see the quasismoothness of $X_{8,8,8}$.

On the other hand, we see that $X_{8,8,8}^{\prime} \subset \mathbb{P}\left(2^{(3)}, 3^{(4)}, 5^{(3)}\right)$ is not quasismooth. Indeed, for $I=\{7,8,9\}$, that is, $a_{7}=a_{8}=a_{9}=5$, neither (Q1) nor (Q2) hold.

The following proposition treats the special situation where some weight of $\mathbb{P}$ divides none of the degrees of a WCI.
Proposition 3.4. Let $X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a (well-formed) quasismooth WCI which is not a linear cone. Assume that there exists $i_{0}$ such that $a_{i_{0}}$ does not divide $d_{j}$ for all $j$. Let $H=\mathcal{O}_{X}(h)$ be the fundamental divisor on $X$, that is, an ample Cartier divisor on $X$ which generates Pic $X$. Then
(i) $X$ has a quotient singularity of type $1 / a_{i_{0}}\left(c_{1}, \ldots, c_{n-c}\right)$ for some $c_{1}, \ldots, c_{n-c} \in$ $\mathbb{Z}_{\geq 0}$ such that $\operatorname{gcd}\left(a_{i_{0}}, c_{1}, \ldots, c_{n-c}\right)=1 ;$
(ii) $a_{i_{0}} \mid h$. As a consequence, we have $|H| \neq \varnothing$.

Proof. Let $f_{1}, \ldots, f_{c} \in \mathbb{C}\left[x_{0}, \ldots, x_{n}\right]$ be the defining equations of $X$ such that $\operatorname{deg} f_{j}=d_{j}$ for $1 \leq j \leq c$ and $X=\left(f_{1}=\cdots=f_{c}=0\right) \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$, where $\operatorname{deg} x_{i}=a_{i}$ for $0 \leq i \leq n$. By applying Proposition 3.1 to $I=\left\{i_{0}\right\}$, we see that there exist distinct integers $e_{1}, \ldots, e_{c} \in\left\{0, \ldots, \hat{i_{0}}, \ldots, n\right\}$ and positive integers $k_{1}, \ldots, k_{c}$ such that $d_{j}=k_{j} a_{i_{0}}+a_{e_{j}}$ for $1 \leq j \leq c$, i.e., we can write

$$
f_{j}=x_{i_{0}}^{k_{j}} x_{e_{j}}+g_{j}
$$

for $1 \leq j \leq c$, where $g_{j}$ is a weighted homogeneous polynomial of degree $d_{j}$.
By the inverse function theorem, we see that $X$ has a quotient singularity of type $1 / a_{i_{0}}\left(a_{0}, \ldots, \hat{a}_{i_{0}}, \ldots, \hat{a}_{e_{1}}, \ldots, \hat{a}_{e_{c}}, \ldots, a_{n}\right)$ at $P_{i_{0}}:=[0: \cdots: 1: \cdots: 0]$. We shall show that $g:=\operatorname{gcd}\left(a_{0}, \ldots, \hat{a}_{e_{1}}, \ldots, \hat{a}_{e_{c}}, \ldots, a_{n}\right)=1$. Suppose that $g>1$.

Claim 3.5. Up to a permutation on $\{1, \ldots, c\}$, we may choose $0 \leq c^{\prime} \leq c$ with the following properties:
(*) For $j=1, \ldots, c^{\prime}$, some monomial in $g_{j}$ does not contain any element of $\left\{x_{e_{j}}, \ldots, x_{e_{c}}\right\}$.
${ }^{(* *)}$ For $j=c^{\prime}+1, \ldots, c$, every monomial in $g_{j}$ contain some of $\left\{x_{e_{c^{\prime}+1}}, \ldots, x_{e_{c}}\right\}$.
Proof of the claim. If (**) holds for all $j=1, \ldots, c$ and $\left\{x_{e_{1}}, \ldots, x_{e_{c}}\right\}$, then we put $c^{\prime}:=0$. Otherwise there is some $j$ such that $1 \leq j \leq c$ and (*) holds for $\left\{x_{e_{1}}, \ldots, x_{e_{c}}\right\}$. We then exchange $\left(f_{1}, e_{1}\right)$ and $\left(f_{j}, e_{j}\right)$ and repeat the same process starting from $j=2$ till we obtain the claim, that is, check whether $\left({ }^{* *)}\right.$ holds for new $\left\{f_{2}, \ldots, f_{c}\right\}$ and $\left\{e_{2}, \ldots, e_{c}\right\}$ and so on.

Hence, for $1 \leq j \leq c^{\prime}$, there exists a monomial in $g_{j}$ of the form $h_{j}=\prod_{i \neq e_{j}, \ldots, e_{c}} x_{i}^{b_{i}}$. Then we have $a_{e_{j}} \equiv \sum_{i \neq e_{j}, \ldots, e_{c}} b_{i} a_{i} \bmod a_{i_{0}}$. Thus we can check one by one that

$$
\begin{equation*}
g \mid a_{e_{j}} \quad \text { for } 1 \leq j \leq c^{\prime} \tag{3}
\end{equation*}
$$

Now let $\Pi:=\left(x_{e_{c^{\prime}+1}}=\cdots=x_{e_{c}}=0\right) \subset \mathbb{P}$. We have $\Pi \subset \operatorname{Sing} \mathbb{P}$, in particular $\Pi \neq \mathbb{P}$.

We also have $\left.f_{j}\right|_{\Pi} \equiv 0$ for $c^{\prime}+1 \leq j \leq c$ by the property ( ${ }^{* *)}$ of $c^{\prime}$. Thus we obtain

$$
\operatorname{dim} \Pi \cap X \geq \operatorname{dim} \Pi-c^{\prime}=\operatorname{dim} \mathbb{P}-c
$$

This contradicts the fact that $X \not \subset \Pi$ since $X$ is not a linear cone. Hence we obtain $g=1$, concluding the proof of Proposition 3.4.

The following proposition is useful for calculating the fundamental divisor of a WCI and is the motivation of the definition of $h$-regular pair (see Definition 4.1).

Proposition 3.6. Let $X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a quasismooth well-formed WCI which is not a linear cone. Let $H=\mathcal{O}_{X}(h)$ be the fundamental divisor of $X$. Assume that there exists $I=\left\{i_{1}, \ldots, i_{k}\right\}$ such that $a_{I}:=\operatorname{gcd}\left(a_{i_{1}}, \ldots, a_{i_{k}}\right)>1$.

Then one of the following holds:
(i) There exist distinct integers $p_{1}, \ldots, p_{k}$ such that $a_{I} \mid d_{p_{1}}, \ldots, d_{p_{k}}$;
(ii) $a_{I} \mid h$.

Proof. We apply Proposition 3.1 to $I=\left\{i_{1}, \ldots, i_{k}\right\}$. Let

$$
P_{I}:=\left(x_{0}=\cdots \hat{x}_{i_{1}}=\cdots=\hat{x}_{i_{k}}=\cdots=x_{n}=0\right) \subset \mathbb{P}
$$

be the $(k-1)$-dimensional stratum corresponding to $I$ and $P_{I}^{0} \subset P_{I}$ be the open toric stratum.

Suppose that condition (Q1) in Proposition 3.1 holds, that is, there exist distinct integers $p_{1}, \ldots, p_{k}$ and nonnegative integers $k_{j, i}$ for $j=1, \ldots, k$ and $i \in I$ such that $d_{p_{j}}=\sum_{i \in I} k_{j, i} a_{i}$. Then we have (i) in this case.

Suppose that (Q2) holds. Then there exist a permutation $p_{1}, \ldots, p_{c}$ of $\{1, \ldots, c\}$, an integer $l<\rho:=\min \{c, k\}$, nonnegative integers $k_{j, i}$ for $j=1, \ldots, c$ and $i \in I$, and distinct integers $e_{l+1}, \ldots, e_{c}$, which satisfy the following:

- for $j=1, \ldots, l$, we have $\sum_{i \in I} k_{j, i} a_{i}=d_{p_{j}}$,
- for $j=l+1, \ldots, c$, we have $a_{e_{j}}+\sum_{i \in I} k_{j, i} a_{i}=d_{p_{j}}$.

We may assume that $\left(f_{p_{j}}=0\right) \cap P_{I}^{0} \neq \varnothing$ since $X$ is irreducible and the linear system $\left|\mathcal{O}_{\mathbb{P}}\left(d_{p_{j}}\right)\right|$ does not have a fixed component. Hence, on $p \in X \cap P_{I}^{0}$, the variety $X$ is analytic locally isomorphic to a quotient singularity of type

$$
\frac{1}{a_{I}}\left(a_{0}, \ldots, \hat{a}_{i_{1}}, \ldots, \hat{a}_{i_{k}}, \ldots, \hat{a}_{e_{l+1}}, \ldots, \hat{a}_{e_{c}}, \ldots, a_{n}\right) \times \mathbb{C}^{k-l}
$$

Now the proof is reduced to the following claim:
Claim 3.7. We have $g:=\operatorname{gcd}\left(a_{0}, \ldots, \hat{a}_{e_{l+1}}, \ldots, \hat{a}_{e_{c}}, \ldots, a_{n}\right)=1$.
Proof of Claim. Suppose that $g>1$. We shall have a similar contradiction as in the proof of Proposition 3.4. As in Claim 3.5, up to a permutation of $\{1, \ldots, c\}$, we may choose $c^{\prime}$ with $l+1 \leq c^{\prime} \leq c$ with the following properties:
(*) For $j=l+1, \ldots, c^{\prime}$, some monomial in $g_{j}$ does not contain any element of $\left\{x_{e_{j}}, \ldots, x_{e_{c}}\right\}$.
(**) For $j=c^{\prime}+1, \ldots, c$, every monomial in $g_{j}$ contain some of $\left\{x_{e_{c^{\prime}+1}}, \ldots, x_{e_{c}}\right\}$.
Let $\Pi:=\left(x_{e_{c^{\prime}+1}}=\cdots=x_{e_{c}}=0\right) \subset \mathbb{P}$. Then, as in Proposition 3.6, we see that $\left.f_{j}\right|_{\Pi} \equiv 0$ for $j=c^{\prime}+1, \ldots, c$ and $\Pi \subset \operatorname{Sing} \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ since $g \mid a_{i}$ for $i \notin\left\{e_{c^{\prime}+1}, \ldots, e_{c}\right\}$. Thus we have $\operatorname{dim} \Pi \cap X \geq \operatorname{dim} \mathbb{P}-c$ as before and it contradicts that $X \not \subset \Pi$ since $X$ is not a linear cone. Thus we obtain the claim.

The sheaf $\mathcal{O}_{X}(1)$ induces a generator of the class group of a quotient singularity of the above type. Since the class group is a cyclic group of order $a_{I}$ as in (1), we see that $a_{I} \mid h$. Thus we have finished the proof of Proposition 3.6.

The following corollary restricts Proposition 3.6 to the smooth case.
Corollary 3.8 [Przyjalkowski and Shramov 2016, Lemma 2.15]. Let $X=X_{d_{1}, \ldots, d_{c}} \subset$ $\mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a smooth WCI. Assume that there exists $I=\left\{i_{1}, \ldots, i_{k}\right\}$ such that $a_{I}:=\operatorname{gcd}\left(a_{i_{1}}, \ldots, a_{i_{k}}\right)>1$.

Then there exist distinct integers $p_{1}, \ldots, p_{k}$ such that $a_{I} \mid d_{p_{1}}, \ldots, d_{p_{k}}$.
Proof. Since $X$ is smooth, the fundamental divisor of $X$ is $\mathcal{O}_{X}(1)$, that is $h=1$ in the notation of Proposition 3.6. Thus the statement follows from Proposition 3.6.

## 4. Regular pairs and Frobenius coin problem

The following definition is motivated by Proposition 3.6 and Corollary 3.8.
Definition 4.1. Let $c \in \mathbb{N}$ and $n \in \mathbb{Z}_{\geq-1}$ be integers and $(d ; a)$ be a pair, where $d=\left(d_{1}, \ldots, d_{c}\right) \in \mathbb{N}_{+}^{c}$ and $a=\left(a_{0}, \ldots, a_{n}\right) \in \mathbb{N}_{+}^{n+1}$. Let $\bar{c}^{+}:=\{1, \ldots, c\}$ and $\bar{n}:=\{0, \ldots, n\}$.

We say that $(d ; a)$ is $h$-regular for a positive integer $h$ if, for any subset $I=$ $\left\{i_{1}, \ldots, i_{k}\right\} \subset \bar{n}$ such that $a_{I}:=\operatorname{gcd}\left(a_{i_{1}}, \ldots, a_{i_{k}}\right)>1$, one of the following holds:
(i) There exist distinct integers $p_{1}, \ldots, p_{k} \in \bar{c}^{+}$such that $a_{I} \mid d_{p_{1}}, \ldots, d_{p_{k}}$;
(ii) $a_{I} \mid h$.

If a pair is $h$-regular for $h=1$, we simply call it regular.
Remark 4.2. For technical reasons, in Definition 4.1 we admit the cases $c=0$ or $n=-1$, i.e., pairs of the form $(d ; \varnothing),(\varnothing ; a)$ and $(\varnothing, \varnothing)$.

We need to fix some notation. If $(d ; a)$ is a pair with $d=\left(d_{1}, \ldots, d_{c}\right) \in \mathbb{N}_{+}^{c}$ and $a=\left(a_{0}, \ldots, a_{n}\right) \in \mathbb{N}_{+}^{n+1}$, then we define

$$
\delta(d ; a):=\sum_{j=1}^{c} d_{j}-\sum_{i=0}^{n} a_{i} .
$$

In the case where the pair $(d ; a)$ comes from a well-formed quasismooth WCI $X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$, we have $\omega_{X} \cong \mathcal{O}_{X}(\delta(d ; a))$.

Let $q$ be a prime number. Set $I_{q}:=\left\{i \in \bar{n}: q \mid a_{i}\right\}$ and $J_{q}:=\left\{j \in \bar{c}^{+}: q \mid d_{j}\right\}$. We consider two new pairs obtained from ( $d ; a$ ). The pair ( $d^{q} ; a^{q}$ ) is given by

$$
d^{q}:=\left(\left(d_{j} / q\right)_{j \in J_{q}},\left(d_{j}\right)_{j \in \bar{c}^{+} \backslash J_{q}}\right), \quad a^{q}:=\left(\left(a_{i} / q\right)_{i \in I_{q}},\left(a_{i}\right)_{i \in \bar{n} \backslash I_{q}}\right)
$$

in which we divided by $q$ all the divisible $d_{j}$ and $a_{i}$ and the pair $(d(q), a(q))$ is given by

$$
d(q):=\left(d_{j}\right)_{j \in J_{q}}, \quad a(q):=\left(a_{i}\right)_{i \in I_{q}}
$$

in which only the divisible $d_{j}$ and $a_{i}$ appear. Note that

$$
\delta(d ; a)=\delta\left(d^{q} ; a^{q}\right)+\frac{q-1}{q} \delta(d(q) ; a(q)) .
$$

Definition 4.3. For a pair $(d ; a)$, we may choose subsets $J_{(d ; a)}=\left\{j_{1}, \ldots, j_{l}\right\} \subset \bar{c}^{+}$ and $I_{(d ; a)}=\left\{i_{1}, \ldots, i_{l}\right\} \subset \bar{n}$ uniquely for some $l \in \mathbb{N}$ so that $d_{j_{k}}=a_{i_{k}}$ for $k=1, \ldots, l$ and $d_{j} \neq a_{i}$ for all $j \in \bar{c}^{+} \backslash J_{(d ; a)}$ and $i \in \bar{n} \backslash I_{(d ; a)}$. We define a pair $(\tilde{d} ; \tilde{a})$ by

$$
\begin{equation*}
(\tilde{d} ; \tilde{a}):=\left(\left(d_{j}\right)_{j \in \bar{c}+\backslash J_{(d ; a)}} ;\left(a_{i}\right)_{i \in \bar{n} \backslash I_{(d ; a)}}\right), \tag{4}
\end{equation*}
$$

that is, we cancel the doubles $\left(d_{j}, a_{i}\right)$ with $d_{j}=a_{i}$.
Lemma 4.4. The pair $(\tilde{d} ; \tilde{a})$ is h-regular if $(d ; a)$ is h-regular.
Proof. Let $I:=\left\{i_{1}, \ldots, i_{k}\right\} \subset \bar{n} \backslash I_{(d ; a)}$ be a subset with $a_{I}>1$. Since $(d ; a)$ is $h$-regular, either (i) holds for some $\left\{p_{1}, \ldots, p_{k}\right\} \subset \bar{c}^{+}$or (ii) holds. In the latter case, there is nothing to check. Thus we consider the former case and need to find $p_{1}^{\prime}, \ldots, p_{k}^{\prime} \in \bar{c}^{+} \backslash J_{(d ; a)}$ such that $a_{I} \mid d_{p_{j}^{\prime}}$ for $j=1, \ldots, k$. Let

$$
J^{\prime}:=\left\{j \in J_{(d ; a)}: a_{I} \mid d_{j}\right\}, \quad I^{\prime}:=\left\{i \in I_{(d ; a)}: a_{I} \mid a_{i}\right\} .
$$

Then we have $\left|I^{\prime}\right|=\left|J^{\prime}\right|=: l^{\prime}$. Let $I^{\prime \prime}:=I \cup I^{\prime}$. By $a_{I}=a_{I^{\prime \prime}}$, there exist distinct integers $p_{1}, \ldots, p_{k+l^{\prime}} \in \bar{c}^{+}$such that $a_{I} \mid d_{p_{j}}$ for $j=1, \ldots, k+l^{\prime}$. Then the set $\left\{p_{1}, \ldots, p_{k+l^{\prime}}\right\} \backslash J^{\prime}$ contains $k$ elements $p_{1}^{\prime}, \ldots, p_{k}^{\prime} \in \bar{c}^{+} \backslash J_{(d ; a)}$ such that $a_{I} \mid d_{p_{j}^{\prime}}$ for $j=1, \ldots, k$. Thus (i) holds for $(\tilde{d} ; \tilde{a})$ and $I$. Hence we see that $(\tilde{d} ; \tilde{a})$ is $h$-regular.

The following straightforward lemmas show how $h$-regular pairs are very suitable for inductive arguments.

Lemma 4.5. Let $(d ; a)$ be an $h$-regular pair and $q$ be a prime not dividing $h$. Then the pairs $\left(d^{q} ; a^{q}\right)$ and $(d(q) ; a(q))$ are $h$-regular. Hence $(d(q) / q ; a(q) / q)$ is also $h$-regular.

Proof. We write the details for the pair $\left(d^{q} ; a^{q}\right)$. The proof for $(d(q) ; a(q))$ is easier.

Let $I=\left\{i_{1}, \ldots, i_{k}\right\} \subset \bar{n}$ such that $a_{I}:=\operatorname{gcd}\left(a_{i_{1}}, \ldots, a_{i_{k}}\right)>1$. By the $h$-regularity of ( $d ; a$ ), we have either condition (i) of Definition 4.1, i.e., there exist distinct integers $p_{1}, \ldots, p_{k}$ such that

$$
a_{I} \mid d_{p_{1}}, \ldots, d_{p_{k}}
$$

or (ii), i.e., $a_{I} \mid h$.
If $q \mid a_{I}$, then we have $q \mid a_{i_{\ell}}$ for all $i_{\ell} \in I$ and $a_{I} \nmid h$. Thus we have (i) and

$$
\left.\operatorname{gcd}\left(\frac{a_{i_{1}}}{q}, \ldots, \frac{a_{i_{k}}}{q}\right)=\frac{a_{I}}{q} \right\rvert\, \frac{d_{p_{1}}}{q}, \ldots, \frac{d_{p_{k}}}{q}
$$

as we wanted.
If $q \nmid a_{I}$, then

$$
\operatorname{gcd}\left(\left(a_{i} / q\right)_{i \in I_{q}^{\prime}},\left(a_{i}\right)_{i \in\left(I \backslash I_{q}^{\prime}\right)}\right)=a_{I}
$$

where $I_{q}^{\prime}:=\left\{i \in I: q \mid a_{i}\right\}$. If $a_{I} \mid h$, then there is nothing to prove, so we can assume that $a_{I} \nmid h$ and that (i) holds. Since $q \nmid a_{I}$, we get that $a_{I}$ divides $\left(d^{q}\right)_{p_{j}}$ for $j=1, \ldots, k$. This concludes the proof.

Lemma 4.6. Let $(d ; a)$ be an $h$-regular pair and $q$ be a prime dividing $h$. Then $\left(d^{q} ; a^{q}\right)$ is $h / q$-regular and $(d(q) ; a(q))$ is h-regular, hence $(d(q) / q ; a(q) / q)$ is $h / q$-regular.

Proof. We give the proof for the pair ( $d^{q} ; a^{q}$ ). Consider a set $I=\left\{i_{1}, \ldots, i_{k}\right\} \subset \bar{n}$ such that $a_{I}:=\operatorname{gcd}\left(a_{i_{1}}, \ldots, a_{i_{k}}\right)>1$ and let $a_{I}^{q}:=\operatorname{gcd}\left(a_{i}^{q}\right)_{i \in I}$ be the $\operatorname{gcd}$ of the $a_{i}$ in ( $d^{q} ; a^{q}$ ).

Assume first that $\operatorname{gcd}\left(a_{I}, q\right)=1$, so that $a_{I}^{q}=a_{I}$. If $a_{I} \mid h$, we obtain that $a_{I}^{q} \mid h / q$ and we are done. If $a_{I} \nmid h$, then there exist distinct integers $p_{1}, \ldots, p_{k}$ such that $a_{I} \mid d_{p_{1}}, \ldots, d_{p_{k}}$. We have $a_{I}^{q} \nmid h / q$ and the $d_{p_{j}}^{q}$ work.

If $a_{I}=q t$ for some positive integer $t$, then $a_{I}^{q}=t$. If $q t \mid h$, we have $t \mid h / q$. If $q t \nmid h$, then there exist distinct integers $p_{1}, \ldots, p_{k}$ such that $a_{I} \mid d_{p_{1}}, \ldots, d_{p_{k}}$. For the same integers, we have $a_{I}^{q} \mid d_{p_{1}} / q, \ldots, d_{p_{k}} / q$ so the first condition of $h / q$-regularity is satisfied and we are done.

The Frobenius coin problem. In this subsection we want to enlighten some interesting connections among the Ambro-Kawamata conjecture, regular pairs and the Frobenius coin problem.

Question 4.7 (Frobenius coin problem). Given positive integers $a_{0}, \ldots, a_{n}$ such that $\operatorname{gcd}\left(a_{0}, \ldots, a_{n}\right)=1$, find the largest integer $G=G\left(a_{0}, \ldots, a_{n}\right)$ so that there do not exist nonnegative integers $x_{0}, \ldots, x_{n}$ satisfying

$$
G=a_{0} x_{0}+\ldots+a_{n} x_{n} .
$$

Such $G$ is called the Frobenius number of $a_{0}, \ldots, a_{n}$.
For $n=1$, it is classically known that

$$
G\left(a_{0}, a_{1}\right)=a_{0} a_{1}-a_{0}-a_{1} .
$$

For $n \geq 2$, the problem is considerably harder: precise methods have been developed to compute $G\left(a_{0}, a_{1}, a_{2}\right)$ and some algorithms and (lower and upper) bounds are known for the general case (see for instance [Johnson 1960] and [Brauer and Shockley 1962]).

By Lemma 2.4, the Ambro-Kawamata conjecture for smooth WCI would follow from the following purely arithmetic statement, which we believe to be of independent interest.

Conjecture 4.8. Let $(d ; a)=\left(d_{1}, \ldots, d_{c} ; a_{0}, \ldots, a_{n}\right) \in \mathbb{N}^{c} \times \mathbb{N}^{n+1}$ be a regular pair such that $a_{i} \neq 1$ and $d_{j} \neq a_{i}$ for any $i, j$. Assume $c \leq n$ and $\operatorname{gcd}\left(a_{0}, \ldots, a_{n}\right)=1$. Then

$$
\delta(d ; a) \geq G\left(a_{0}, \ldots, a_{n}\right)
$$

One of the best known lower bounds for $G$ is given in [Brauer 1942]. Let $a_{0}, \ldots, a_{n}$ be positive coprime integers, set $g_{j}:=\operatorname{gcd}\left(a_{0}, \ldots, a_{j}\right)$ for $j=0, \ldots, n$ and consider

$$
\operatorname{Br}\left(a_{0}, \ldots, a_{n}\right):=\sum_{j=1}^{n} a_{j} \frac{g_{j-1}}{g_{j}}-\sum_{i=0}^{n} a_{i} .
$$

Brauer proved that $\operatorname{Br}\left(a_{0}, \ldots, a_{n}\right) \geq G\left(a_{0}, \ldots, a_{n}\right)$. Set $d_{j}:=a_{j} g_{j-1} / g_{j}$ for $j=1, \ldots, n$. Then it is easy to check that $(d ; a):=\left(d_{1}, \ldots, d_{n} ; a_{0}, \ldots, a_{n}\right)$ is actually a regular pair.

On the other hand, it is not difficult to see that, considering big prime numbers $p$ and $q$, the pair $(p q, 6 p, 6 q ; 2 p, 3 p, 2 q, 3 q)$ is regular, $\delta(d ; a) \geq G\left(a_{0}, \ldots, a_{n}\right)$, but $\delta(d ; a)<\operatorname{Br}\left(a_{0}, \ldots, a_{n}\right)$.

This shows that regular pairs can give better bounds for the Frobenius number with respect to the known ones. For this reason, it seems to be a challenge and interesting problem to study Conjecture 4.8.

Remark 4.9. It is not difficult to check that Conjecture 4.8 is true for $c=1,2$, which implies that the nonvanishing holds for a smooth WCI of codimension 1 or 2. For simplicity, we omit the detail in the codimension 2 case.

For $c=1$, a stronger and more general result is given in Lemma 6.1, which is the key step to prove Theorem 1.3.

## 5. Proof of Theorem 1.2

Theorem 1.2 is the combination of Corollary 5.3 and Corollary 5.13 below.
Smooth case. The pair $(d ; a)$ in the following lemma does not come from a nonempty WCI. Nevertheless this lemma is important in the proof of Proposition 5.2.

Lemma 5.1. Let $(d ; a) \in \mathbb{N}_{+}^{c} \times \mathbb{N}_{+}^{n+1}$ be a regular pair such that $a_{i} \neq d_{j}$ for any $i, j$. Let $q$ be a prime number such that $q \mid a_{i}$ and $q \mid d_{j}$ for any $i, j$. Then

$$
\delta(d ; a) \geq c q .
$$

Moreover, if the equality holds, then $c=n+1$.
Proof. Note that $c \geq n+1$ which does not occur for a nonempty WCI.
Assume first that $q$ is the only prime dividing the $a_{i}$, that is for any $i=0, \ldots, n$, we have $a_{i}=q^{\alpha_{i}}$ for some $\alpha_{i} \geq 1$. We can assume that $a_{0} \geq a_{1} \geq \ldots \geq a_{n}$. We can also order the $d_{j}$ in such a way that $v_{q}\left(d_{s}\right) \geq v_{q}\left(d_{t}\right)$ for any $s \leq t$, where $v_{q}\left(d_{j}\right)=\max \left\{e \in \mathbb{N}: q^{e} \mid d_{j}\right\}$. Then we have $a_{i} \mid d_{i+1}$ for any $i=0, \ldots, n$ and so

$$
\sum_{j=1}^{c} d_{j}-\sum_{i=0}^{n} a_{i}=\sum_{k=1}^{c-n-1} d_{n+1+k}+\sum_{i=0}^{n}\left(d_{i+1}-a_{i}\right) \geq c q
$$

and the equality is possible only if $c=n+1, d_{j}=2 q$ and $a_{i}=q$ for any $i, j$.
Assume now that $q \neq 2$ and that $q$ and 2 are the only primes dividing the $a_{i}$, that is for any $i=0, \ldots, n$ we have $a_{i}=2^{\alpha_{i}} q^{\beta_{i}}$ for some $\alpha_{i} \geq 0$ and $\beta_{i} \geq 1$ such that $\alpha_{i}>0$ for at least one $i$. We proceed by induction on $t=\max _{0 \leq i \leq n}\left\{\beta_{i}\right\}$, the greatest power of $q$ dividing at least one $a_{i}$.

Suppose $t=1$. We can assume that $v_{2}\left(a_{i}\right) \geq v_{2}\left(a_{j}\right)$ and $v_{2}\left(d_{i}\right) \geq v_{2}\left(d_{j}\right)$ for any $i \leq j$. Then again $a_{i} \mid d_{i+1}$ for any $i=0, \ldots, n$ and we conclude as before.

Suppose $t \geq 2$. Let $I_{q^{t}}:=\left\{i \in \bar{n}: q^{t} \mid a_{i}\right\}$ and $J_{q^{t}}:=\left\{j \in \bar{c}^{+}: q^{t} \mid d_{j}\right\}$. We consider the following pairs: $\left(d^{\prime} ; a^{\prime}\right)$, where $d^{\prime}=\left(\left(d_{j} / q\right)_{j \in J_{q^{t}}},\left(d_{j}\right)_{j \in \bar{c}^{+} \backslash J_{q^{t}}}\right)$ and $\left.a^{\prime}=\left(\left(a_{i} / q\right)_{i \in I_{q^{t}}},\left(a_{i}\right)_{i \in \bar{n} \backslash I_{q^{t}}}\right)\right)$ and $\left(d^{\prime \prime} ; a^{\prime \prime}\right)$, where $d^{\prime \prime}=\left(d_{j} / q\right)_{j \in J_{q^{t}}}$ and $a^{\prime \prime}=$ $\left(a_{i} / q\right)_{i \in I_{q^{t}}}$. It is straightforward to check as in Lemma 4.5 that ( $d^{\prime} ; a^{\prime}$ ) and ( $d^{\prime \prime} ; a^{\prime \prime}$ ) are regular. Consider the regular pair ( $\tilde{d}^{\prime} ; \tilde{a}^{\prime}$ ) constructed in (4) which satisfies $\tilde{d}_{j}^{\prime} \neq \tilde{a}_{i}^{\prime}$ for any $i \in \bar{n} \backslash I_{\left(d^{\prime} ; a^{\prime}\right)}, j \in \bar{c}^{+} \backslash J_{\left(d^{\prime} ; a^{\prime}\right)}$, where $I_{\left(d^{\prime} ; a^{\prime}\right)} \subset \bar{n}$ and $J_{\left(d^{\prime} ; a^{\prime}\right)} \subset \bar{c}^{+}$ are the subsets defined in Definition 4.3.

Let

$$
\begin{aligned}
& m:=\mid\left\{j \in J_{q^{t}}: d_{j} / q=a_{i} \mid \text { for some } i \in \bar{n} \backslash I_{q^{t}}\right\} \mid, \\
& \bar{m}:=\mid\left\{i \in I_{q^{t}}: d_{j}=a_{i} / q \mid \text { for some } j \in \bar{c}^{+} \backslash J_{q^{t}}\right\} \mid .
\end{aligned}
$$

Note that $\left|I_{\left(d^{\prime} ; a^{\prime}\right)}\right|=\left|J_{\left(d^{\prime} ; a^{\prime}\right)}\right| \leq m+\bar{m}$. Let $k:=\left|J_{q^{t}}\right|$. By induction on $t$, we may assume that we have $\delta\left(d^{\prime} ; a^{\prime}\right)=\delta\left(\tilde{d}^{\prime} ; \tilde{a}^{\prime}\right) \geq(c-m-\bar{m}) q$ and $\delta\left(d^{\prime \prime} ; a^{\prime \prime}\right) \geq k q$. We have that $\left|I_{q^{t}}\right| \leq\left|J_{q^{t}}\right|$ since $(d ; a)$ is regular. Since $m \leq k$ and $\bar{m} \leq\left|I_{q^{t}}\right| \leq k$, we obtain

$$
\begin{align*}
\delta(d ; a)=\delta\left(d^{\prime} ; a^{\prime}\right)+(q-1) \delta\left(d^{\prime \prime} ; a^{\prime \prime}\right) & \geq(c-2 k) q+(q-1) k q \\
& =c q-2 k q+k q^{2}-k q \\
& =c q+k q(q-3) \geq c q \tag{5}
\end{align*}
$$

because $q \geq 3$. The equality is possible only if we have it for both ( $\tilde{d}^{\prime} ; \tilde{a}^{\prime}$ ) and ( $d^{\prime \prime} ; a^{\prime \prime}$ ). This implies by induction on $t$ that $c=n+1$ in this case.

We now pass to the general case. For any prime $p$, different from $q$ and 2, let $e_{p}:=\max \left\{e \in \mathbb{N}: p^{e} \mid a_{i}\right.$ for some $\left.i\right\}$. The proof is by induction on $D=\sum_{p} e_{p}$, where the index varies over all prime numbers different from $q$ and 2 . The case $D=0$ has already been treated in the first part of the proof. So assume $D \geq 1$ and that the inequality holds up to $D-1$. Consider $\left(d^{p} ; a^{p}\right)$ and let

$$
\begin{aligned}
& m_{p}:=\mid\left\{j \in \bar{c}^{+}: d_{j} / p=a_{i} \text { for some } i \in \bar{n} \backslash I_{p}\right\} \mid \\
& \bar{m}_{p}:=\mid\left\{i \in \bar{n}: d_{j}=a_{i} / p \text { for some } j \in \bar{c}^{+} \backslash J_{p}\right\} \mid
\end{aligned}
$$

Let us again consider the pair ( $\tilde{d}^{p} ; \tilde{a}^{p}$ ) as in Definition 4.3 by removing subsets $J_{\left(d^{p} ; a^{p}\right)} \subset \bar{c}^{+}$and $I_{\left(d^{p} ; a^{p}\right)} \subset \bar{n}$. Then this satisfies the hypothesis $\left(\tilde{d}^{p}\right)_{j} \neq\left(\tilde{a}^{p}\right)_{i}$ for any $i$ and $j$. We again have that $\left|J_{\left(d^{p} ; a^{p}\right)}\right| \leq m_{p}+\bar{m}_{p}$. By induction on $D$, we obtain $\delta\left(d^{p} ; a^{p}\right)=\delta\left(\tilde{d}^{p} ; \tilde{a}^{p}\right) \geq\left(c-m_{p}-\bar{m}_{p}\right) q$. Now consider the pair $(d(p) / p ; a(p) / p)$. Again by induction on $D$, we obtain $\delta(d(p) / p ; a(p) / p) \geq s q$, where $s:=\mid\{j \in$ $\left.\bar{c}^{+}: p \mid d_{j}\right\} \mid$. We see that $m_{p} \leq s$ by the definition of $m_{p}$. Let $s^{\prime}:=\left|\left\{i \in \bar{n}: p \mid a_{i}\right\}\right|$. We see that $s^{\prime} \leq s$ by the regularity of $(d ; a)$ and that $\bar{m}_{p} \leq s^{\prime}$ by the definition of $\bar{m}_{p}$. Thus we have $\bar{m}_{p} \leq s$. By these inequalities and $p \geq 3$, we conclude that

$$
\begin{aligned}
\delta(d ; a) & =\delta\left(d^{p} ; a^{p}\right)+(p-1) \delta(d(p) / p ; a(p) / p) \\
& \geq\left(c-m_{p}-\bar{m}_{p}\right) q+(p-1) s q \\
& =c q+p s q-m_{p} q-\bar{m}_{p} q-s q \geq c q+p s q-3 s q \geq c q
\end{aligned}
$$

as we wanted. Again, the equality is possible only if $c=n+1$.
By using Lemma 5.1, we prove the following key proposition.
Proposition 5.2. Let $(d ; a) \in \mathbb{N}_{+}^{c} \times \mathbb{N}_{+}^{n+1}$ be a regular pair such that $a_{i}>1$ and $a_{i} \neq d_{j}$ for any $i, j$. Then the following holds.
(i) We have

$$
\begin{equation*}
\delta(d ; a) \geq c \tag{6}
\end{equation*}
$$

(ii) If $\operatorname{gcd}\left(a_{0}, \ldots, a_{n}\right)=1$, then the equality holds only if $(d ; a)$ is of the form $\left(6^{(s)}, 1^{(c-s)} ; 2^{(s)}, 3^{(s)}\right)$ for some integer $s$.

Proof. (i) The proof is by induction on $n$, the case $n=0$ being obvious. We can assume that no prime divides every $a_{i}$, otherwise we are in the case of Lemma 5.1. In particular, we may assume that there is a prime $q \neq 2$ which divides some $a_{i}$. Let

$$
\begin{aligned}
m & :=\mid\left\{j \in \bar{c}^{+}: d_{j} / q=a_{i} \text { for some } i \in \bar{n} \backslash I_{q}\right\} \mid, \\
\bar{m} & :=\mid\left\{i \in \bar{n}: d_{j}=a_{i} / q \text { for some } j \in \bar{c}^{+} \backslash J_{q}, d_{j} \neq 1\right\} \mid, \\
\ell & :=\left|\left\{i \in \bar{n}: a_{i}=q\right\}\right|, \quad s:=\left|\left\{j \in \bar{c}^{+}: q \mid d_{j}\right\}\right|=\left|J_{q}\right|
\end{aligned}
$$

We note that $m \leq s$ by definition and $\ell+\bar{m} \leq s$ by the regularity.
Case 1: Suppose that $\ell+m+\bar{m} \geq 1$. Then the pair ( $d^{q} ; a^{q}$ ) has some redundant $a_{i}$, in the sense that $a_{i} / q=1, d_{j} / q=a_{i}$ or $d_{j}=a_{i} / q$ for some $i, j$. That is, we consider a regular pair ( $\tilde{d}^{q}, \tilde{a}^{q}$ ) and, by removing all $\tilde{a}_{i}^{q}=1$, we obtain a new regular pair $\left(\hat{d}^{q} ; \hat{a}^{q}\right) \in \mathbb{N}_{+}^{\hat{c}} \times \mathbb{N}_{+}^{\hat{n}+1}$ for some $\hat{c} \leq c$ and $\hat{n} \leq n$. Note that $\hat{n}<n$ by the hypothesis $\ell+m+\bar{m} \geq 1$. Let $\ell_{1}:=\left|\left\{j \in \bar{c}^{+}: d_{j}=1\right\}\right|$ and $\ell^{\prime}:=\min \left\{\ell, \ell_{1}\right\}$. Then we see that $\hat{c} \geq c-m-\bar{m}-\ell^{\prime}$ by the construction of $\left(\tilde{d}^{q} ; \tilde{a}^{q}\right)$. Since we have $\left|\left\{i \in \bar{n} \backslash I_{\left(d^{q} ; a^{q}\right)}: \tilde{a}_{i}^{q}=1\right\}\right|=\ell-\ell^{\prime}$, we obtain, by induction on $n$, that

$$
\delta\left(d^{q} ; a^{q}\right)=\delta\left(\hat{d}^{q} ; \hat{a}^{q}\right)-\left(\ell-\ell^{\prime}\right) \geq \hat{c}-\left(\ell-\ell^{\prime}\right) \geq c-\ell-m-\bar{m}
$$

By applying Lemma 5.1 to $(d(q), a(q))$, we obtain

$$
\delta(d(q) ; a(q)) \geq s q
$$

By these and $\ell+m+\bar{m} \leq 2 s$, we obtain

$$
\begin{align*}
\delta(d ; a)=\delta\left(d^{q} ; a^{q}\right)+\frac{q-1}{q} \delta(d(q) ; a(q)) & \geq c-\ell-m-\bar{m}+(q-1) s \\
& \geq c+q s-3 s \geq c \tag{7}
\end{align*}
$$

since $q \geq 3$.
Case 2: Suppose now that $\ell+m+\bar{m}=0$. Then the pair ( $d^{q} ; a^{q}$ ) satisfies the assumptions of the proposition. We note that

$$
\delta(d ; a)=\delta\left(d^{q} ; a^{q}\right)+\frac{q-1}{q} \delta(d(q) ; a(q))>\delta\left(d^{q} ; a^{q}\right)
$$

since we have $\delta(d(q) ; a(q))>0$ by Lemma 5.1. So we can replace the pair ( $d ; a$ ) with ( $d^{q} ; a^{q}$ ) without changing the number $c$ of the $j$ and we can repeat the argument from the beginning of the proof (possibly changing the prime $q$ ) till either we end up in Case 1 or we reach the situation of Lemma 5.1. In both cases, we are done and obtain (6).
(ii) We now study when the identity holds in the case $\operatorname{gcd}\left(a_{0}, \ldots, a_{n}\right)=1$.

Note that the case $n=1$ is clear, being equivalent to asking $a_{0} a_{1}-a_{0}-a_{1}=1$.

Assume $n \geq 2$ and let $q \neq 2$ be a prime number such that $q \mid a_{i}$ for some $i$. We shall follow the proof of the inequality. In particular, we look at

$$
\delta(d ; a)=\delta\left(d^{q} ; a^{q}\right)+\frac{q-1}{q} \delta(d(q) ; a(q))
$$

With the same notation as above, we note that the equality can hold only if we are in Case 1 and, by Lemma 5.1, the number $\left|\left\{i \in \bar{n}: q \mid a_{i}\right\}\right|$ must be equal to $s=\left|J_{q}\right|$. Moreover we obtain $q=3$ by (7). This implies that the only possible prime numbers that divide at least one $a_{i}$ are 2 and 3 .

We must also have $m=s$ and $\ell+\bar{m}=s$. By $m=s$, we see that any $d_{j} / 3 \in \mathbb{N}$ must be equal to some $a_{i}$ which is not divisible by 3 . Hence we can write

$$
(d ; a)=\left(3 \cdot 2^{\beta_{1}}, \ldots, 3 \cdot 2^{\beta_{s}}, 2^{\beta_{s+1}}, \ldots, 2^{\beta_{c}} ; 2^{\beta_{1}}, \ldots, 2^{\beta_{s}}, 3 \cdot 2^{\alpha_{s+1}}, \ldots, 3 \cdot 2^{\alpha_{n+1}}\right)
$$

for some nonnegative integers $\alpha_{i}$ and $\beta_{i}$. Then
$\delta(d ; a)=\delta\left(2^{\beta_{s+1}}, \ldots, 2^{\beta_{c}} ; 2^{\alpha_{s+1}}, \ldots, 2^{\alpha_{n+1}}\right)+\frac{2}{3} \delta\left(3 \cdot 2^{\beta_{1}}, \ldots, 3 \cdot 2^{\beta_{s}} ; 3 \cdot 2^{\alpha_{s+1}}, \ldots, 3 \cdot 2^{\alpha_{n+1}}\right)$.
Also note that $\ell+\bar{m}=s$ implies that $n+1-s=\left|\left\{i \in \bar{n}: 3 \mid a_{i}\right\}\right|=\ell+\bar{m}=s$, thus $n+1=2 s$. By the regularity of ( $d(3) ; a(3))$ and the assumption $d_{j} \neq a_{i}$ for any $i, j$, to have the equality $\delta(d(3) ; a(3))=3 s$ we need $\beta_{j}=1$ for $j=1, \ldots, s$ and $\alpha_{i}=0$ for $i=s+1, \ldots, n+1$, which implies

$$
\delta(d ; a)=\delta\left(2^{\beta_{s+1}}, \ldots, 2^{\beta_{c}} ; 1, \ldots, 1\right)+2 s=\sum_{i=1}^{c-s} 2^{\beta_{s+i}}-(n+1-s)+2 s
$$

i.e., $c=\delta(d ; a)=\sum_{i=1}^{c-s} 2^{\beta_{s+i}}+s$.

Hence, we must have $\beta_{j}=0$ for $j=s+1, \ldots, c$, which finishes the proof.
As a corollary of Proposition 5.2, we obtain the nonemptiness of $|\mathcal{O}(1)|$ and the smoothness of its general member on a smooth Fano or Calabi-Yau WCI.
Corollary 5.3. Let $X:=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a well-formed smooth Fano or Calabi-Yau WCI which is not a linear cone. Let $c_{1}:=\left|\left\{i \in \bar{n}: a_{i}=1\right\}\right|$. Then the following hold.
(i) We have $c_{1} \geq c$. Moreover the equality is possible only if $X$ is Calabi-Yau of type $X_{6, \ldots, 6} \subset \mathbb{P}\left(1^{(c)}, 2^{(c)}, 3^{(c)}\right)$.
(ii) The linear system $\left|\mathcal{O}_{X}(1)\right|$ is nonempty and its general member $H$ is smooth.

Proof. (i) We may assume that $a_{0} \leq \cdots \leq a_{n}$. Thus we have $a_{0}=\cdots=a_{c_{1}-1}=1$. Since $X$ is smooth, we see that $\left(d_{1}, \ldots, d_{c} ; a_{c_{1}}, \ldots, a_{n}\right)$ is regular. By this and Proposition 5.2(i), we obtain

$$
\delta\left(d_{1}, \ldots, d_{c} ; a_{c_{1}}, \ldots, a_{n}\right) \geq c .
$$

By the assumptions, $0 \geq \delta(d ; a) \geq c-c_{1}$, and this implies the former statement.

Let $(d ; a)$ be a regular pair which satisfies $c_{1}=c$. Let $(\hat{d} ; \hat{a})$ be the regular pair obtained by removing all $a_{i}=1$. Then $(\hat{d} ; \hat{a})$ satisfies the hypothesis of Proposition 5.2(ii) since $(d ; a)$ defines a smooth WCI. Hence, by Proposition 5.2(ii), we see that $(\hat{d} ; \hat{a})=\left(6^{(c)} ; 2^{(c)}, 3^{(c)}\right)$ and $(d ; a)=\left(6^{(c)} ; 2^{(c)}, 3^{(c)}, 1^{(c)}\right)$.
(ii) By the latter part of Proposition 5.2, we can assume that $X$ is not of the form $X_{6}, \ldots, 6 \subset \mathbb{P}\left(2^{(c)}, 3^{(c)}, 1^{(c)}\right)$; otherwise the conclusion is immediate. In particular, we may assume $c_{1} \geq c+1$.

By (i), we see that $\left|\mathcal{O}_{X}(1)\right| \neq \varnothing$. Since $X$ is smooth and well-formed, we have Sing $\mathbb{P}\left(a_{0}, \ldots, a_{n}\right) \cap X=\varnothing$. Thus $H \cap \operatorname{Sing} \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)=\varnothing$. Hence it is enough to check $H$ is quasismooth at $P:=\Pi(p)$, where $p \in \Pi^{-1}\left(\left(x_{0}=\ldots=x_{c_{1}-1}=0\right) \cap X\right)$ and $\Pi: \mathbb{A}^{n+1} \backslash\{0\} \rightarrow \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ is the quotient map.

Set $H_{i}:=X \cap\left(x_{i}=0\right)$ for $i=0, \ldots, c_{1}-1$. We shall look at the Jacobi matrices of $X$ and $H_{i} \subset \mathbb{P}\left(a_{0}, \ldots, a_{i-1}, a_{i+1}, \ldots, a_{n}\right)$. Let $f_{1}, \ldots, f_{c}$ be the defining equations of $X$ such that $\operatorname{deg} f_{j}=d_{j}$. For $i=0, \ldots, n$, set

$$
\boldsymbol{v}_{i}(p):=\left(\begin{array}{c}
\partial f_{1} / \partial x_{i} \\
\vdots \\
\partial f_{c} / \partial x_{i}
\end{array}\right)(p)
$$

The Jacobi matrix $J_{X}(p)$ and $J_{H_{i}}(p)$ of $X$ and $H_{i}$ can be written as

$$
J_{X}(p)=\left(\boldsymbol{v}_{0}(p), \ldots, \boldsymbol{v}_{n}(p)\right), \quad J_{H_{i}}(p)=\left(\boldsymbol{v}_{0}(p), \ldots, \boldsymbol{v}_{i-1}(p), \boldsymbol{v}_{i+1}(p), \ldots, \boldsymbol{v}_{n}(p)\right)
$$

Since $X$ is quasismooth, there exist linearly independent vectors

$$
\boldsymbol{v}_{i_{1}}(p), \ldots, \boldsymbol{v}_{i_{c}}(p)
$$

Since $c_{1} \geq c+1$, we can choose $i$ so that $i \notin\left\{i_{1}, \ldots, i_{c}\right\}$. Then we see that $H_{i}$ is quasismooth at $P:=\Pi(p)$. Thus a general member $H$ is also quasismooth at $P$.

Remark 5.4. Let $X \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a smooth WCI as in Corollary 5.3. For $I \subset \bar{n}$ such that $a_{I}=1$, it may a priori happen that (Q1) does not hold, but (Q2) holds. That is why we make an argument as in Corollary 5.3 (ii).

Remark 5.5. Let $X_{d_{1}, \ldots, d_{c}}$ be a smooth WCI as in Corollary 5.3. Motivated by a question by Andreas Höring, we consider the description of the base locus Bs $\left|\mathcal{O}_{X}(1)\right|$.

Up to reordering $d_{1}, \ldots, d_{c}$, we can assume that there is an integer $c^{\prime} \leq c$ with the following properties: for $1 \leq j \leq c^{\prime}$, there are weighted homogeneous polynomials $f_{j}\left(x_{c_{1}}, \ldots, x_{n}\right)$ of degree $d_{j}$ and, for $c^{\prime}+1 \leq j \leq c$, all monomials of degree $d_{j}$ contain one of the variables $x_{0}, \ldots, x_{c_{1}-1}$ of weights 1 . Since the base locus Bs $\left|\mathcal{O}_{X}(1)\right|$ is $\left(x_{0}=\cdots=x_{c_{1}-1}=0\right) \cap X_{d_{1}, \ldots, d_{c}}$, it is isomorphic to a general WCI $Y_{d_{1}, \ldots, d_{c^{\prime}}} \subset \mathbb{P}\left(a_{c_{1}}, \ldots, a_{n}\right)$.

Thus the base locus is again a WCI. However this is not necessarily (quasi)smooth in general. We shall see this in Example 5.6.

Example 5.6. Let $X:=X_{231,231,26} \subset \mathbb{P}:=\mathbb{P}\left(3,3,7,7,11,11,1^{(447)}\right)$ be a general WCI. We can check that this is a smooth Fano WCI as follows: for $I=\{0,1\},\{2,3\}$ or $\{4,5\}$, (that is, two variables of weights 3,7 or 11 ), we have (Q1) for $d_{1}=231$, $d_{2}=231$. Also, for $I=\{0,1,2,3\}$ or $\{0,1,4,5\}$, we have $(\mathrm{Q} 1)$ for $d_{1}=231$, $d_{2}=231, d_{3}=26$ since $26=7 \cdot 2+3 \cdot 4=11+3 \cdot 5$. For $I=\{2,3,4,5\}$, we have (Q2) for $d_{1}=231, d_{2}=231, d_{3}=26=7 \cdot 2+11+1$. By Proposition 3.1, we see that $X$ is quasismooth, and smooth since $X \cap \operatorname{Sing} \mathbb{P}=\varnothing$.

The base locus Bs $\left|\mathcal{O}_{X}(1)\right|$ is a WCI $Y:=Y_{231,231,26} \subset \mathbb{P}^{\prime}:=\mathbb{P}(3,3,7,7,11,11)$. This is not quasismooth. Indeed, for $I=\{2,3,4,5\}$, neither (Q1) nor (Q2) holds because of the lack of suitable degree 26 polynomials. In fact, $Y$ is a nonnormal surface singular along a curve $\left(x_{0}=x_{1}=f_{1}=f_{2}=0\right) \subset \mathbb{P}^{\prime}$, where $f_{1}, f_{2}$ are part of defining polynomials of degrees 231 and $x_{0}, x_{1}$ are the variables of weights 3 .

Hence we can not expect smoothness of the base locus of the fundamental linear system even if it contains a smooth member.

Remark 5.7. Let $W=W_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a smooth WCI which is not a linear cone, where $a_{i}>1$ for any $i=0, \ldots, n$. By Corollary 5.3 we know that $W$ is not Fano. Then we can consider a WCI

$$
X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}, 1^{(\ell)}\right)
$$

where $\ell=\delta(W)+1$. In this way $X$ is a smooth Fano with $-K_{X}=\mathcal{O}_{X}(1)$ and Bs $\left|\mathcal{O}_{X}(1)\right|$ is exactly $W$.

In Corollary 5.3 we showed that for a smooth Fano WCI, the general member of the fundamental divisor is quasismooth. This is not true in general for a quasismooth Fano WCI as the following example shows.

Example 5.8. Let $X=X_{35} \subset \mathbb{P}\left(5,7,2^{(k)}, 3^{(k)}\right)$ where $k \geq 5$. Then $X$ is a quasismooth Fano WCI with fundamental divisor $\mathcal{O}_{X}(6)$, but $X_{35,6} \subset \mathbb{P}\left(5,7,2^{(k)}, 3^{(k)}\right)$ is not quasismooth. However, we see that a general member of $\left|\mathcal{O}_{X}(6)\right|$ has only terminal singularities. Indeed it has an isolated singularity at $[*: *: 0: \cdots: 0]$ which is locally isomorphic to $0 \in\left(x_{1}^{3}+\cdots+x_{k}^{3}+x_{k+1}^{2}+\cdots+x_{2 k}^{2}=0\right) \subset \mathbb{C}^{2 k}$.

It is also natural to look at the general element of $\left|-K_{X}\right|$ in the case of a Fano variety X. For instance, Shokurov [1979] and Reid [1983] proved that a Fano 3-fold with only canonical Gorenstein singularities admits an anticanonical member with only Du Val singularities. Here we give an example of a smooth Fano WCI whose anticanonical members are singular (and not quasismooth). See also [Höring and Voisin 2011, 2.12] for an example of a Fano 4-fold with singular fundamental divisor.

Example 5.9. (cf., [Sano 2014, Example 2.9]) For $m \in \mathbb{Z}_{>0}$, let $X$ be a weighted hypersurface $X=X_{(2 m+1)(2 m+2)} \subset \mathbb{P}\left(1^{(1+2 m(2 m+1))}, 2 m+1,2 m+2\right)$ of degree $(2 m+1)(2 m+2)$. Then we see that $-K_{X}=\mathcal{O}_{X}(2)$ and the linear system $\left|-i K_{X}\right|$ does not contain a smooth member for $i=1, \ldots, m$. Thus we can not expect a smooth element of the plurianticanonical system on a Fano manifold. However, in the above example, we can find a member with only terminal singularities. Moreover, the base locus of $|H|$ consists of a point.

Remark 5.10. It is well known that following the arguments in [Ambro 1999, Section 5] or [Kawamata 2000, Section 5] and assuming Conjecture 1.1, it is possible to show that the general element of $\left|-m K_{X}\right|$ has always only klt singularities for $m>0$ such that $-m K_{X}$ is Cartier (we thank Chen Jiang for pointing this fact out to us).

Finally, we also get the following corollary, which generalizes [Przyjalkowski and Shramov 2017, Corollary 4.2] to any codimension.

Corollary 5.11. Let $X:=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a well-formed smooth Fano or Calabi-Yau WCI which is not a linear cone. Let $c_{1}:=\left|\left\{i \in \bar{n}: a_{i}=1\right\}\right|$. Then $c_{1}>I(X):=-\delta(d ; a)=\sum_{i=0}^{n} a_{i}-\sum_{j=1}^{c} d_{j}$.
Proof. Consider the regular pair $(d ; a)$ associated with $X$. We may assume that $a_{0} \leq \cdots \leq a_{n}$, so that $a_{0}=\cdots=a_{c_{1}-1}=1$. Let $\left(d^{\prime} ; a^{\prime}\right)$ be the pair $\left(d ; a_{c_{1}}, \ldots, a_{n}\right)$, where we took away every 1 from ( $d ; a$ ). This pair is regular with no $a_{i}=1$ and so by Proposition 5.2 we get

$$
\delta\left(d^{\prime} ; a^{\prime}\right) \geq c>0
$$

which implies

$$
\delta(d ; a)=\delta\left(d^{\prime} ; a^{\prime}\right)-c_{1}>-c_{1}
$$

i.e., $I(X)<c_{1}$, as we wanted.

General case. The following is a key proposition to deduce the nonvanishing in the quasismooth Fano case.
Proposition 5.12. Let $h \in \mathbb{N}_{+}$and $(d ; a) \in \mathbb{N}_{+}^{c} \times \mathbb{N}_{+}^{n+1}$ be an $h$-regular pair with $c \geq 1$. If $a_{i} \nmid h$ for any $i=0, \ldots, n$ and $a_{i} \neq d_{j}$ for any $i, j$, then

$$
\delta(d ; a)>0
$$

Proof. Let us write $h=p_{1}^{\alpha_{1}} \cdots p_{k}^{\alpha_{k}}$, where the $p_{i}$ are distinct prime numbers. The proof is by induction on $\alpha=\sum_{i=1}^{k} \alpha_{i} \geq 0$. If $\alpha=0$, then the pair $(d ; a)$ is regular and the statement follows from Proposition 5.2, so we assume $\alpha \geq 1$.

Let $p$ be a prime number dividing $h$ and consider $\left(d^{p} ; a^{p}\right)$. As usual,

$$
\delta(d ; a)=\delta\left(d^{p} ; a^{p}\right)+\frac{p-1}{p} \delta(d(p) ; a(p)) .
$$

By Lemma 4.6, $\left(d^{p} ; a^{p}\right)$ and $(d(p) / p ; a(p) / p)$ are $h / p$-regular. Note that there does not exist $i$ such that $a_{i} / p=1$ by the hypothesis $a_{i} \nmid h$. Thus, after cancellation on $\left(d^{p} ; a^{p}\right)$ (see Definition 4.3), we see that $\left(\tilde{d^{p}} ; \tilde{a}^{p}\right)$ and $(d(p) / p ; a(p) / p)$ are $h / p$-regular and satisfy the hypothesis of the proposition.

If $p \mid a_{i}$ or $p \mid d_{j}$ for some $i$ or $j$, then we obtain $\delta(d(p) / p ; a(p) / p)>0$ by the induction hypothesis and conclude $\delta(d ; a)>0$ by induction since we have either $\delta\left(d^{p} ; a^{p}\right)=\delta\left(\tilde{d^{p}} ; \tilde{a^{p}}\right)>0$ or $\left(\tilde{d^{p}} ; \tilde{a^{p}}\right)$ is empty.

If $p \nmid a_{i}$ and $p \nmid d_{j}$ for any $i, j$, then $\delta\left(d^{p} ; a^{p}\right)>0$ by the induction hypothesis since $\left(d^{p} ; a^{p}\right)=(d ; a)$ is $h / p$-regular. Moreover $(d(p) ; a(p))$ is empty. Hence we can again conclude that $\delta(d ; a)>0$.

Corollary 5.13. Let $X=X_{d_{1}, \ldots, d_{c}} \subset \mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a well-formed quasismooth WCI which is Fano or Calabi-Yau and which is not a linear cone. Then $|H| \neq \varnothing$ for any ample Cartier divisor $H$ on $X$.
Proof. Write $H=\mathcal{O}_{X}(h)$. If there exists $i \in \bar{n}$ such that $a_{i} \mid h$, then we are done. Otherwise, we are in the situation of Proposition 5.12, and so the variety can not be Fano or Calabi-Yau since $(d ; a)$ is $h$-regular by Proposition 3.6.

## 6. Weighted hypersurfaces

The following lemma gives a proof of a generalized version of Conjecture 4.8 in the case $c=1$.

Lemma 6.1. Let $a_{0}, \ldots, a_{n}$ be positive integers, $n \geq 1$ and set

$$
h:=\operatorname{lcm}_{i \neq j}\left(\operatorname{gcd}\left(a_{i}, a_{j}\right)\right)
$$

Assume that $a_{i} \nmid h$ for any $i$ and set $f:=\operatorname{lcm}\left(a_{0}, \ldots, a_{n}\right)$. Then

$$
f-\sum_{i=0}^{n} a_{i} \geq \operatorname{lcm}\left(a_{s}, a_{t}\right)-a_{s}-a_{t}
$$

for any s and $t$.
Proof. We first note that, for any proper subset $I$ of $\bar{n}:=\{0, \ldots, n\}$, we have $f \neq \operatorname{lcm}_{i \in I}\left(a_{i}\right)$. In fact, suppose that the equality holds and let $k \in \bar{n} \backslash I$. For any prime power $p^{e}$ such that $e \geq 1$ and $p^{e} \mid a_{k}$, we have $p^{e} \mid f$. In particular, we have $p^{e} \mid a_{\ell}$ for some $\ell \in I$. This implies that $p^{e} \mid \operatorname{gcd}\left(a_{k}, a_{\ell}\right)$ and so $a_{k} \mid h$, which is a contradiction. In particular, $f \geq 2 \mathrm{lcm}_{i \in I}\left(a_{i}\right)$.

The proof of the lemma is by induction on $n$ and the case $n=1$ is trivial, so assume $n \geq 2$. Let $s, t \in \bar{n}$ be such that $s \neq t$. Then

$$
f-\sum_{i=0}^{n} a_{i}=\frac{f}{2}-a_{s}-a_{t}+\frac{f}{2}-\sum_{i \neq s, t} a_{i} \geq \operatorname{lcm}\left(a_{s}, a_{t}\right)-a_{s}-a_{t}+\operatorname{lcm}_{i \neq s, t}\left(a_{i}\right)-\sum_{i \neq s, t} a_{i}
$$

If $n=2$, then we have

$$
\operatorname{lcm}_{i \neq s, t}\left(a_{i}\right)-\sum_{i \neq s, t} a_{i}=0
$$

and we are done.
If $n \geq 3$, then we have

$$
\operatorname{lcm}_{i \neq s, t}\left(a_{i}\right)-\sum_{i \neq s, t} a_{i} \geq \operatorname{lcm}\left(a_{s^{\prime}}, a_{t^{\prime}}\right)-a_{s^{\prime}}-a_{t^{\prime}}
$$

for $s^{\prime}, t^{\prime} \in \bar{n} \backslash\{s, t\}$ such that $s^{\prime} \neq t^{\prime}$ by induction on $n$ and $\operatorname{lcm}\left(a_{s^{\prime}}, a_{t^{\prime}}\right)-a_{s^{\prime}}-a_{t^{\prime}} \geq 0$ because we are assuming that $a_{i} \nmid h$ for any $i$.
Proposition 6.2. Let $X=X_{d} \subset \mathbb{P}=\mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a well-formed, quasismooth hypersurface of degree $d$ which is not a linear cone. Let $H$ be an ample Cartier divisor on $X$ such that $H-K_{X}$ is ample.

Then $|H|$ is not empty.
Proof. Write $\mathcal{O}_{X}(H)=\mathcal{O}_{X}(h)$ for a positive integer $h$.
By Proposition 3.4, we can assume that $a_{i} \mid d$ for any $i$. Then $X$ is a Cartier divisor which intersects any stratum $P_{\{i, j\}}$ in some interior point. The condition of $H$ to be Cartier is then equivalent to

$$
\operatorname{lcm}_{i \neq j}\left(\operatorname{gcd}\left(a_{i}, a_{j}\right)\right) \mid h
$$

If there exists $a_{i}$ such that $a_{i} \mid h$, then we are done. So assume that $a_{i} \nmid h$ for any $i$ and let $f:=\operatorname{lcm}\left(a_{0}, \ldots, a_{n}\right)$. By Lemma 6.1, we get

$$
f-\sum_{i=0}^{n} a_{i} \geq \operatorname{lcm}\left(a_{s}, a_{t}\right)-a_{s}-a_{t}
$$

for any $s$ and $t$. Since $h>f-\sum_{i=0}^{n} a_{i}$ (because $H-K_{X}$ is ample and $f \mid d$ ) and $g:=\operatorname{gcd}\left(a_{s}, a_{t}\right) \mid h$ for any $s \neq t$, we can use the Frobenius number $G\left(a_{s} / g, a_{t} / g\right)=$ $(1 / g)\left(\operatorname{lcm}\left(a_{s}, a_{t}\right)-a_{s}-a_{t}\right)$ as on page 2382 to conclude that there are nonnegative integers $\lambda_{s}, \lambda_{t}$ such that

$$
\lambda_{s} a_{s}+\lambda_{t} a_{t}=h
$$

which implies that $|H|$ is not empty by Lemma 2.4.
In the following, we prove the basepoint freeness on a Gorenstein weighted hypersurface.
Proposition 6.3. Let $X=X_{d} \subset \mathbb{P}=\mathbb{P}\left(a_{0}, \ldots, a_{n}\right)$ be a well-formed, quasismooth hypersurface of degree $d$ which is not a linear cone such that $K_{X}$ is Cartier. Let $H$ be the fundamental divisor of $X$ and $h$ be the positive integer such that $H=\mathcal{O}_{X}(h)$.

Then $L=K_{X}+m H$ is globally generated for any $m \geq n$.

Proof. Suppose by contradiction that there is a point $p=\left[p_{0}: \cdots: p_{n}\right] \in \mathrm{Bs}|L|$ and take $\ell$ such that $L=\mathcal{O}_{X}(\ell)$.

Note that if $p_{s} \neq 0$ for some $s$, then $a_{s} \nmid h$, otherwise $x_{s}^{e} \in|L|$ for some positive integer $e$ and so $p \notin \mathrm{Bs}|L|$. Also note that, for all $i \in \bar{n}$ such that $a_{i} \nmid h$, we have $a_{i} \mid d$ by Proposition 3.4.

Assume first that there exists a unique $s \in \bar{n}$ such that $p_{s} \neq 0$. Since $p \in X$ and $a_{s} \nmid h$, we get that $a_{s} \mid d$. Let $f_{d}$ be the defining polynomial of $X_{d}$. If $f_{d}$ contains a monomial $x_{s}^{d / a_{s}}$, then we obtain $p \notin X_{d}$ and this is a contradiction. If $f_{d}$ does not contain such a monomial, then it should contain a monomial of the form $x_{s}^{k} x_{i}$ for some $k>0$ and $i \neq s$ by the quasismoothness of $X_{d}$. Then we see that $a_{s} \mid a_{i}$ by $a_{s} \mid d$, and $X_{d}$ has a quotient singularity of index $a_{s}$. Thus we obtain $a_{s} \mid h$ and this is a contradiction.

Hence we can assume that there exist $s$ and $t$ such that $s \neq t, p_{s} \neq 0$ and $p_{t} \neq 0$, thus $a_{s}, a_{t} \nmid h$. We have

$$
\ell=d-\sum_{i=0}^{n} a_{i}+m h=d-\sum_{a_{i} \nmid h} a_{i}-\sum_{a_{i} \mid h} a_{i}+m h
$$

Assume that $-\sum_{a_{i} \mid h} a_{i}+m h \geq 1$. Since $a_{i} \mid d$ for all $i$ such that $a_{i} \nmid h$, we can apply Lemma 6.1 to conclude that

$$
\ell>\operatorname{lcm}\left(a_{s}, a_{t}\right)-a_{s}-a_{t},
$$

which implies that $x_{s}^{e_{s}} x_{t}^{e_{t}} \in|L|$ for some nonnegative integers $e_{s}$ and $e_{t}$. So we again have $p \notin \mathrm{Bs}|L|$.

Assume now that $-\sum_{a_{i} \mid h} a_{i}+m h \leq 0$. Then we can check that $\left|\left\{i: a_{i}=h\right\}\right| \geq$ $n-1$, because $m \geq n$. Moreover, since $\mathbb{P}$ is well-formed, the greatest common factor of any $n$ weights is 1 . By these, when $\left|\left\{i: a_{i}=h\right\}\right|=n$, we have $h=1$ and $\mathbb{P}=\mathbb{P}\left(a_{0}, 1, \ldots, 1\right)$ for some $a_{0}>1$. When $\left|\left\{i: a_{i}=h\right\}\right|=n-1$, we have $h=2$ and $\mathbb{P}=\mathbb{P}(1,1,2, \ldots, 2)$. In both cases, we can check that $L$ is basepoint free, and we have derived a contradiction.

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# Generalized Kuga-Satake theory and good reduction properties of Galois representations 

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In previous work, we described conditions under which a single geometric representation $\Gamma_{F} \rightarrow H\left(\overline{\mathbb{Q}}_{\ell}\right)$ of the Galois group of a number field $F$ lifts through a central torus quotient $\tilde{H} \rightarrow H$ to a geometric representation. In this paper, we prove a much sharper result for systems of $\ell$-adic representations, such as the $\ell$-adic realizations of a motive over $F$, having common "good reduction" properties. Namely, such systems admit geometric lifts with good reduction outside a common finite set of primes. The method yields new proofs of theorems of Tate (the original result on lifting projective representations over number fields) and Wintenberger (an analogue of our main result in the case of a central isogeny $\widetilde{H} \rightarrow H)$.

## 1. Introduction

Let $F$ be a number field, and let $\Gamma_{F}=\operatorname{Gal}(\bar{F} / F)$ be its absolute Galois group with respect to a fixed algebraic closure $\bar{F}$. A fundamental theorem of Tate (see [Serre 1977, §6]) asserts that $H^{2}\left(\Gamma_{F}, \mathbb{Q} / \mathbb{Z}\right)$ vanishes; as a result, all (continuous, $\ell$-adic) projective representations of $\Gamma_{F}$ lift to genuine representations, and more generally, whenever $\widetilde{H} \rightarrow H$ is a surjection of linear algebraic groups over $\overline{\mathbb{Q}}_{\ell}$ with kernel equal to a central torus in $\widetilde{H}$, all representations $\rho_{\ell}: \Gamma_{F} \rightarrow H\left(\overline{\mathbb{Q}}_{\ell}\right)$ lift to $\widetilde{H}\left(\overline{\mathbb{Q}}_{\ell}\right)$.

The $\ell$-adic representations of greatest interest in number theory are those with conjectural connections to the theories of motives and automorphic forms; if the monodromy group of $\rho_{\ell}$ is semisimple, then it is expected - by conjectures of Fontaine-Mazur, Tate, Grothendieck-Serre, and Langlands - that the $\rho_{\ell}$ arising from pure motives or automorphic forms are precisely those that are geometric in the sense of Fontaine-Mazur, i.e., unramified outside a finite set of places of $F$, and de Rham at all places dividing $\ell$. The paper [Patrikis 2016c] established a variant of Tate's lifting theorem for such geometric Galois representations. There are

[^18]obstructions when $F$ has real embeddings, but at least for totally imaginary $F$, any geometric $\rho_{\ell}: \Gamma_{F} \rightarrow H\left(\overline{\mathbb{Q}}_{\ell}\right)$ satisfying a natural "Hodge symmetry" requirement admits a geometric lift $\tilde{\rho}_{\ell}: \Gamma_{F} \rightarrow \widetilde{H}\left(\overline{\mathbb{Q}}_{\ell}\right)$ [Patrikis 2016c, Theorem 3.2.10]. This geometric lifting theorem leads to a precise expectation for the corresponding lifting problem for motivic Galois representations. Namely, if $\mathcal{G}_{F, E}$ denotes the motivic Galois group for pure motives over $F$ with coefficients in a number field $E$ we will make this setup precise in Section 2, but for now the reader may take homological motives under the standard conjectures - and if $\widetilde{H} \rightarrow H$ is now a surjection of groups over $E$ with central torus kernel, then we conjecture [Patrikis 2016c, Conjecture 4.3.1] that any motivic Galois representation $\rho: \mathcal{G}_{F, E} \rightarrow H$ lifts to $\widetilde{H}$, at least after some finite extension of coefficients:


There is essentially one classical example (with several variants) of this conjecture, a well-known construction of Kuga and Satake [1967], which associates to a complex, for our purposes projective, K3 surface $X$ a complex abelian variety $\mathrm{KS}(X)$, related by an inclusion of Hodge-structures $H^{2}(X, \mathbb{Q}) \subset H^{1}(\mathrm{KS}(X), \mathbb{Q})^{\otimes 2}$. In the motivic Galois language, finding $\operatorname{KS}(X)$ amounts (when $F=\mathbb{C}$ ) to finding a lift $\tilde{\rho}$ of the representation $\rho_{X}: \mathcal{G}_{\mathbb{C}, \mathbb{Q}} \rightarrow H=\mathrm{SO}\left(H^{2}(X)(1)\right)$, through the surjection $\widetilde{H}=\operatorname{GSpin}\left(H^{2}(X)(1)\right) \rightarrow H$. Progress on the general conjecture, when the motives in question do not lie in the Tannakian subcategory of motives generated by abelian varieties, seems to require entirely new ideas. ${ }^{1}$

The aim of this paper is to establish a Galois-theoretic result which is necessary for this conjecture to hold, but considerably more delicate than the basic geometric lifting theorem of [Patrikis 2016c, Theorem 3.2.10]. Namely, any motive $M$ over $F$ has good reduction outside a finite set of primes: for any choice of variety $X$ in whose cohomology $M$ appears, $X$ spreads out as a smooth projective scheme over $\mathcal{O}_{F}[1 / N]$ for some integer $N$. In particular, by the base-change theorems of étale cohomology [Deligne 1977] and the crystalline $p$-adic comparison isomorphism [Faltings 1989], for any motivic Galois representation $\rho: \mathcal{G}_{F, E} \rightarrow H$, the $\lambda$-adic realizations $\rho_{\lambda}: \Gamma_{F} \rightarrow H\left(E_{\lambda}\right)$ have good reduction outside a finite set of primes $S$, in the sense (also see Definition 1.1) that each $\rho_{\lambda}$ factors through $\Gamma_{F, S \cup S_{\lambda}}$ and is crystalline at all places of $S_{\lambda} \backslash\left(S_{\lambda} \cap S\right)$; here $S_{\lambda}$ denotes the primes of $F$ with the same residue characteristic as $\lambda$, and $\Gamma_{F, S \cup S_{\lambda}}$ is the Galois group of the maximal extension of $F$ inside $\bar{F}$ that is unramified outside of $S \cup S_{\lambda}$. Certainly a necessary

[^19]condition for the generalized Kuga-Satake conjecture to hold is that the realizations $\left\{\rho_{\lambda}\right\}_{\lambda}$ of $\rho$ should lift to geometric representations $\tilde{\rho}_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ that likewise have good reduction outside a common finite set of places $P$. This is what we will show, as a consequence of a more general result. To state it, we first make a couple of definitions.
Definition 1.1. A collection $\left\{\rho_{\lambda}: \Gamma_{F} \rightarrow H\left(\bar{E}_{\lambda}\right)\right\}_{\lambda}$, as $\lambda$ varies over finite places of $E$, of geometric Galois representations is ramification-compatible if there exist
(1) a finite set $S$ of places of $F$ such that each $\rho_{\lambda}$ is unramified outside of $S \cup S_{\lambda}$, i.e., factors through
$$
\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)
$$
and for $v$ in $S_{\lambda}$ but not in $S,\left.\rho_{\lambda}\right|_{\Gamma_{F_{v}}}$ is crystalline; and
(2) a central cocharacter $\omega: \boldsymbol{G}_{m, E} \rightarrow H$ and a collection of conjugacy classes
$$
\left\{\left[\mu_{\tau}: \boldsymbol{G}_{m, \bar{E}} \rightarrow H_{\bar{E}}\right]\right\}_{\tau: F \hookrightarrow \bar{E}}
$$
satisfying $\left[\mu_{\tau}\right]=\omega \cdot\left[\mu_{c \tau}^{-1}\right]$ for any choice of complex conjugation $c \in \operatorname{Gal}(\bar{E} / \mathbb{Q})$, such that for all $E$-embeddings $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$, inducing via $\tau$ some $\tau_{\iota_{\lambda}}: F_{v} \hookrightarrow \bar{E}_{\lambda}$, the conjugacy class $\left[\mu_{\tau} \otimes_{\bar{E}_{l_{\lambda}}} \bar{E}_{\lambda}\right.$ ] is equal to the conjugacy class of $\tau_{\iota_{\lambda}}$-labeled Hodge-Tate cocharacters associated to $\left.\rho_{\lambda}\right|_{\Gamma_{F_{v}}}$.
If a single representation $\rho_{\lambda}$ satisfies the condition in item (1), we say $\rho_{\lambda}$ has good reduction outside $S$. If it satisfies the condition in item (2) (for some collection of cocharacters $\omega,\left\{\mu_{\tau}\right\}$ ), then we say it satisfies Hodge symmetry.

Remark 1.2. - Note that the $\rho_{\lambda}$ need not be "compatible" in the usual sense (frobenii acting compatibly): if the "coefficients" of the $\rho_{\lambda}$ are bounded in a rather strong sense - there exists a common number field over which their frobenius characteristic polynomials are defined - one expects that our collection of $\rho_{\lambda}$ should partition (dividing up the $\lambda$ 's) into finitely many compatible systems. ${ }^{2}$

- The Hodge symmetry requirement of part (2) of Definition 1.1 is not the most general constraint that pertains to a compatible system of $\ell$-adic representations. It will always hold for the $\lambda$-adic realizations of motives, as we will see in Lemma 2.3, when we take the ambient group to be the motivic Galois group $G_{\rho}$ of the underlying motivic Galois representation $\rho$, and not some larger group. But there may be compatible systems (of motivic origin) where the criterion in part (2) of Definition 1.1 fails; for example, consider $\rho_{\ell}: \Gamma_{\mathbb{Q}} \rightarrow \mathrm{PGL}_{3}\left(\overline{\mathbb{Q}}_{\ell}\right)$ given

[^20]by the projectivization of $\kappa_{\ell} \oplus 1 \oplus 1$, with $\kappa_{\ell}$ denoting the $\ell$-adic cyclotomic character. The relation $\left[\mu_{\tau}\right]=\omega \cdot\left[\mu_{c \tau}^{-1}\right]$ implies in this case $\left[\mu_{\tau}\right]=\left[\mu_{\tau}^{-1}\right]$, which is false. A way around this is given in [Patrikis 2016c, §3.2], where the Hodge symmetry hypothesis is formulated in a way that conjecturally holds for any geometric representation $\rho_{\ell}: \Gamma_{F} \rightarrow H\left(\overline{\mathbb{Q}}_{\ell}\right)$, regardless of its (reductive) algebraic monodromy group. The proof of [Patrikis 2016c, Theorem 3.2.10] thus requires a slightly trickier group-theoretic argument than the one we require here. We have opted in this paper to keep the simpler condition (2) above, so as to focus on what is new in the arguments, and because of its obvious centrality from a motivic point of view (in particular, its sufficiency for Corollary 1.6).
Here is the main theorem:
Theorem 1.3. Let $E$ be a number field, and let $\widetilde{H} \rightarrow H$ be a surjection of linear algebraic groups over $E$ with kernel equal to a central torus in $\tilde{H}$. Let $F$ be a totally imaginary number field, and let $S$ be a finite set of places of $F$ containing the archimedean places. Fix a set of cocharacters $\left\{\mu_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$ satisfying the "Hodge symmetry" condition of part (2) of Definition 1.1. Then there exists a finite set of places $P \supset S$ such that for any place $\lambda$ of $E$, any embedding $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$, and any geometric representation $\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)$ such that

- $\rho_{\lambda}$ has good reduction outside $S$, and
- the conjugacy classes of labeled Hodge-Tate cocharacters of $\rho_{\lambda}$ are induced via $\iota_{\lambda}$ from $\left\{\mu_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$ (again, see Definition 1.1 for details),
the representation $\rho_{\lambda}$ admits a geometric lift $\tilde{\rho}_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \tilde{H}\left(\bar{E}_{\lambda}\right)$ having good reduction outside $P$.

In particular, if $\left\{\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)\right\}_{\lambda}$ is a ramification-compatible system, then there exist a finite set of places $P \supset S$ and lifts $\tilde{\rho}_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ such that $\left\{\tilde{\rho}_{\lambda}\right\}_{\lambda}$ is a ramification-compatible system.
Remark 1.4. All results of this paper, once we take into account the caveat of [Patrikis 2016c, §2.8] (see too [Patrikis 2015, Proposition 5.5]), admit straightforward variants when $F$ has real places. Thus, for real $F$, the analogue of Theorem 1.3 either holds exactly as written, or after replacing $F$ by any totally imaginary (e.g., composite with a quadratic imaginary) extension. We do not want to discuss this at any length here, but we simply remind the reader that the prototypical example in which $F$ is totally real, and Theorem 1.3 fails as stated, is that of the projective motivic Galois representation associated to a mixed-parity Hilbert modular form.

The proof of this theorem is completed in Theorem 1.3. The typical application is to the collection of Galois representations $\left\{\rho_{\lambda}\right\}_{\lambda}$ associated to a motivic Galois representation; we make this precise in Corollary 1.6.

Next we describe applications of Theorem 1.3 to the more general problem of lifting through surjections $H^{\prime} \rightarrow H$ with central kernel of multiplicative type. Here in general we cannot expect as strong a result as Theorem 1.3. First, the Hodge-Tate cocharacters of $\rho_{\lambda}$ may not lift to $H^{\prime}$, in which case there can be no geometric lifts to $H^{\prime}$. Second, even if the Hodge-Tate cocharacters lift, the Galois representations may only lift after a finite base change of $F$. For example, if $\rho_{\lambda}$ is the projectivization of the Galois representation associated to a weight 3 modular form, $\operatorname{det}\left(\rho_{\lambda}\right): \Gamma_{\mathbb{Q}} \rightarrow \bar{E}_{\lambda}^{\times}$does not admit a square root. A beautiful result of Wintenberger [1995, Théorème 2.1.4, Théorème 2.1.7] shows that when $H^{\prime} \rightarrow H$ is an isogeny, a result similar to Theorem 1.3 holds, as long as in the conclusion $F$ is replaced by a suitable finite extension. Here we treat the general case of multiplicative-type quotients:

Theorem 1.5 (see Corollary 3.18). Let $H^{\prime} \rightarrow H$ be a surjection of linear algebraic groups over $E$ whose kernel is central and of multiplicative type. Let $F$ be a number field, and let $S$ be a finite set of places of $F$ containing the archimedean places. Fix a set of cocharacters $\left\{\mu_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$ as in part (2) of Definition 1.1, and moreover, assume that each $\mu_{\tau}$ lifts to a cocharacter of $H^{\prime}$.

Then there exist a finite set of places $P \supset S$, and a finite extension $F^{\prime} / F$, such that any geometric representation $\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)$ having good reduction outside $S$, and whose Hodge-Tate cocharacters arise from the set $\left\{\mu_{\tau}\right\}_{\bar{\tau}: F \hookrightarrow \bar{E}}$ via some embedding $\bar{E} \hookrightarrow \bar{E}_{\lambda}$, admits a geometric lift $\tilde{\rho}_{\lambda}: \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow H^{\prime}\left(\bar{E}_{\lambda}\right)$ having good reduction outside $P$.

In particular, if $\left\{\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)\right\}_{\lambda}$ is a ramification-compatible system with Hodge cocharacters $\left\{\mu_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$, then there exist a finite set of places $P \supset S$, a finite extension $F^{\prime} / F$, and lifts $\tilde{\rho}_{\lambda}: \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow H^{\prime}\left(\bar{E}_{\lambda}\right)$ such that $\left\{\tilde{\rho}_{\lambda}\right\}_{\lambda}$ is a ramification-compatible system.

We deduce Wintenberger's original result in Corollary 3.16. Our proof differs in an essential way, as it passes through Theorem 1.3, which cannot be deduced from the methods of [Wintenberger 1995]. Our problem resembles Wintenberger's in that both lead to a basic difficulty of annihilating cohomological obstruction classes in infinitely many Galois cohomology groups, one for each $\lambda$, but needing to do so in an "independent-of- $\lambda$ " fashion. The arguments themselves, however, are in fact orthogonal to one another: Wintenberger always kills cohomology by making a finite base change on $F$, whereas that is precisely what we are forbidden from doing if we want the more precise results of Theorem 1.3. Moreover, our methods also yield a novel proof of Tate's original vanishing theorem (see Corollary 3.9). In fact, Corollary 3.9 establishes a more precise form of Tate's theorem: the latter of course shows that the image under the canonical map $H^{2}\left(\Gamma_{F, S}, \mathbb{Z} / N\right) \rightarrow H^{2}\left(\Gamma_{F}, \mathbb{Q} / \mathbb{Z}\right)$ is zero, and our refinement quantifies how much additional ramification must be
added, and how much the coefficients must be enlarged, in order to annihilate $H^{2}\left(\Gamma_{F, S}, \mathbb{Z} / N\right)$. Our arguments thus achieve, from scratch, a satisfying common generalization of the theorems of Wintenberger and Tate.

In Corollaries 3.12 and 3.14, we give a couple of applications to lifting $\lambda$-adic realizations such that the associated "similitude characters" (e.g., determinant or Clifford norm) of the lifts form strongly compatible systems. Note that even in the case of the classical Kuga-Satake construction, this compatibility is only achieved as a consequence of having an arithmetic descent of the (Hodge-theoretically defined) Kuga-Satake abelian variety; such a descent depends on the deformation theory of K3 surfaces and monodromy arguments (due to Deligne [1972] and André [1996a]).

Our final result is the promised motivic application:
Corollary 1.6. Let $F$ be a totally imaginary number field, let $E$ be a number field, and let $\mathcal{G}_{F, E}$ denote the motivic Galois group, defined by André's motivated cycles, of pure motives over $F$ with coefficients in $E$ (see Section 2). Let $\widetilde{H} \rightarrow H$ be a surjection of linear algebraic groups over $E$ whose kernel is a central torus in $\tilde{H}$, and let $\rho: \mathcal{G}_{F, E} \rightarrow H$ be any motivic Galois representation, with associated $\lambda$-adic realizations $\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)$ for some finite set $S$ of places of $F$. Then there exist a finite, independent of $\lambda$, set $P \supset S$ of places of $F$ and, for all $\lambda$, lifts

such that each $\tilde{\rho}_{\lambda}$ is de Rham at all places in $S_{\lambda}$, and is moreover crystalline at all places in $S_{\lambda} \backslash\left(S_{\lambda} \cap P\right)$.

Now suppose $H^{\prime} \rightarrow H$ is a surjection of linear algebraic groups whose kernel is central but of arbitrary multiplicative type, and let $\rho: \mathcal{G}_{F, E} \rightarrow H$ again be a motivic Galois representation. Assume that the labeled Hodge cocharacters of $\rho$ (see Definition 2.1) lift to $H^{\prime}$. Then there exist a finite set $P \supset S$ of places of $F, a$ finite extension $F^{\prime} / F$, and, for all $\lambda$, lifts

such that $\tilde{\rho}_{\lambda}$ is de Rham at all places above $S_{\lambda}$, and is moreover crystalline at all places above $S_{\lambda} \backslash\left(S_{\lambda} \cap P\right)$.

Remark 1.7. Even admitting a strong finiteness conjecture, that there are finitely many isomorphism classes of $\rho$ (as in Corollary 1.6), having coefficients in $E$, prescribed Hodge-Tate cocharacters, and good reduction outside a fixed finite set $S$, Theorem 1.3 still says rather more than Corollary 1.6, since even for fixed $\lambda$ it applies to infinitely many distinct $\rho_{\lambda}$ simultaneously (because we have not bounded the coefficients: recall the example of modular forms of weight two whose nebentypus characters have unbounded conductor, even though supported on the fixed finite set $S$ of primes).

We close this introduction by emphasizing what we do not prove. The realizations $\left\{\rho_{\lambda}\right\}_{\lambda}$ of $\rho$ should, moreover, form a weakly compatible system of Galois representations in the sense that the conjugacy class of $\rho_{\lambda}\left(f r_{v}\right)$ is defined over $\bar{E}$ and is suitably independent of $\lambda$ (for $v$ outside $S \cup S_{\lambda}$ ), and in turn one would hope to construct lifts $\tilde{\rho}_{\lambda}$ with the same frobenius compatibility. This problem seems to be out of reach: I know of no way to establish such results using only Galois-theoretic techniques, although indeed they would follow (assuming the standard conjectures) from the generalized Kuga-Satake conjecture. Alternatively, it is possible to establish results of this nature in settings where $\rho_{\lambda}$ and $\tilde{\rho}_{\lambda}$ are constructed as automorphic Galois representations: the most significant example of this is the recent work of Kret and Shin [2016] associating GSpin-value Galois representations to certain (discrete series at infinity, Steinberg at some finite place) cuspidal automorphic representations of $\mathrm{GSp}_{2 n}\left(\boldsymbol{A}_{F}\right)$, for $F$ totally real. Crucially, their construction uses the already-known construction of ( $\mathrm{SO}_{2 n+1}$-valued) Galois representations for the restrictions of such automorphic representations to $\mathrm{Sp}_{2 n}\left(\boldsymbol{A}_{F}\right)$, so it is very much in the spirit of the generalized Kuga-Satake lifting problem.

## 2. Hodge symmetry

In this section we establish a motivic setting in which our general Galois-theoretic results apply; this setting will both serve as motivation for subsequent sections and allow us to deduce Corollary 1.6 from Theorem 1.3 (and Corollary 3.18). The reader who does not find the motivic language illuminating can safely skip this section.

Rather than working with (pure) homological motives and assuming the standard conjectures, we work with a category of motives that is unconditionally semisimple and Tannakian - and in which we can prove unconditional results - but that would, under the standard conjectures, turn out to be equivalent to the category of homological motives. Namely, let $\mathcal{M}_{F, E}$ denote André's category of motivated motives over $F$ with coefficients in $E$; see [André 1996b]. We begin by elaborating on the consequences of Hodge symmetry in $\mathcal{M}_{F, E}$. Throughout this discussion, it will be convenient to fix embeddings $\tau_{0}: F \hookrightarrow \bar{E}$ and $\iota_{\infty}: \bar{E} \hookrightarrow \mathbb{C}$. The composite
$\iota_{\infty} \tau_{0}: F \hookrightarrow \mathbb{C}$ yields a Betti fiber functor $H_{\iota_{\infty} \tau_{0}}: \mathcal{M}_{F, E} \rightarrow \operatorname{Vect}_{E}$, making $\mathcal{M}_{F, E}$ into a neutral Tannakian category over $E$. We denote by $\mathcal{G}=\mathcal{G}_{B}\left(\iota_{\infty} \tau_{0}\right)$ the associated Tannakian group (tensor automorphisms of the fiber functor), so that $H_{l_{\infty} \tau_{0}}$ induces an equivalence of tensor categories $\mathcal{M}_{F, E} \xrightarrow{\sim} \operatorname{Rep}(\mathcal{G})$.

We will consider other cohomological realizations on $\mathcal{M}_{F, E}$, and their comparisons with the Betti fiber functor. Let $H_{d R}: \mathcal{M}_{F, E} \rightarrow \mathrm{Fil}_{F \otimes_{\mathbb{Q}} E}$ denote the de Rham realization, taking values in filtered $F \otimes_{\mathbb{Q}} E$-modules, and for each place $\lambda$ of $E$, let $H_{\lambda}$ denote the $\lambda$-adic realization, which takes values in finite $E_{\lambda}$-modules with a continuous action of $\Gamma_{F}$. For all embeddings $\tau: F \hookrightarrow \bar{E}$, we obtain an $\bar{E}$-valued fiber functor

$$
\omega_{d R, \tau}: M \mapsto \operatorname{gr}\left(e_{\tau} H_{d R}(M)\right)
$$

where $e_{\tau}$ is the idempotent induced by $\tau \otimes 1: F \otimes_{\mathbb{Q}} E \rightarrow \bar{E}$. Let $\mathcal{G}_{d R}(\tau)=\operatorname{Aut}^{\otimes}\left(\omega_{d R, \tau}\right)$ be the associated Tannakian group over $\bar{E}$. Of course this fiber functor factors through the category $\mathrm{Gr}_{\bar{E}}$ of graded $\bar{E}$-vector spaces

so we obtain a corresponding homomorphism $\mu_{\tau}: \boldsymbol{G}_{m, \bar{E}} \rightarrow \mathcal{G}_{d R}(\tau)$. Without specifying $\tau$, we obtain a fiber functor (see [Deligne and Milne 1982, §3]) $\omega_{d R}=\mathrm{gr}{ }^{\bullet} H_{d R}$ on $\mathcal{M}_{F, E}$ valued in projective $F \otimes_{\mathbb{Q}} E$-modules. By [loc. cit., Theorem 3.2], the functor $\underline{\operatorname{Hom}}^{\otimes}\left(H_{\iota_{\infty} \tau_{0}}, \omega_{d R}\right)$ is a $\mathcal{G}$-torsor over $F \otimes_{\mathbb{Q}} E$. In particular, for all $\tau: F \hookrightarrow \bar{E}$, we can choose a point of $\underline{\operatorname{Hom}}^{\otimes}\left(H_{l_{\infty} \tau_{0}} \otimes_{E} \bar{E}, \omega_{d R, \tau}\right)$ to induce a cocharacter $\mu_{\tau}$ of $\mathcal{G}_{\bar{E}}$, and the conjugacy class $\left[\mu_{\tau}\right.$ ] of $\mu_{\tau}$ is independent of this choice.
Definition 2.1. For each $\tau: F \hookrightarrow \bar{E}$, we call any $\mu_{\tau}: \boldsymbol{G}_{m, \bar{E}} \rightarrow \mathcal{G}_{\bar{E}}$ as above a $\tau$-labeled Hodge cocharacter; it is a representative of the conjugacy class of cocharacters $\left[\mu_{\tau}\right]$, the latter being canonically independent of any of the above choices of isomorphisms of fiber functors.
Lemma 2.2. For all $\sigma \in \operatorname{Gal}(\bar{E} / E),\left[\mu_{\tau}\right]=\left[\mu_{\sigma \tau}\right]$. In particular, $\left[\mu_{\tau}\right]$ only depends on the restriction of $\tau$ to the maximal $C M$ (or totally real) subfield $F_{c m}$ of $F$.
Proof. We decompose $F \otimes_{\mathbb{Q}} E=\prod_{i} E_{i}$ into a product of fields, writing $p_{i}$ for the projection onto $E_{i}$. Any $E$-algebra homomorphism $\tau: F \otimes_{\mathbb{Q}} E \rightarrow \bar{E}$ factors through $p_{i(\tau)}$ for a unique $i(\tau)$, and then the $\operatorname{Gal}(\bar{E} / E)$-orbit of $\tau$ is precisely those $E$-algebra homomorphisms (i.e., embeddings $F \hookrightarrow \bar{E}$ ) $\tau^{\prime}: F \otimes_{\mathbb{Q}} E \rightarrow \bar{E}$ such that $i(\tau)=i\left(\tau^{\prime}\right)$. The first claim follows, since both $\omega_{d R, \tau}$ and $\omega_{d R, \sigma \tau}$ can be factored through $p_{i(\tau)} \circ \omega_{d R}$. The second claim follows from the first, and the fact that all motives arise by scalar extension from motives with coefficients in CM (or totally real) fields [Patrikis 2016c, Lemma 4.1.22].

Next note that the canonical weight-grading on $\mathcal{M}_{F, E}$ induces a central weight homomorphism

$$
\omega: \boldsymbol{G}_{m, E} \rightarrow \mathcal{G}
$$

and likewise for any other choice of fiber functor and Tannakian group (because $\omega$ is central, it is in fact canonically independent of any choice of isomorphism between fiber functors). Hodge symmetry then results from the complex conjugation action on Betti cohomology, interpreted via the Betti-de Rham comparison isomorphism, which is a distinguished $\mathbb{C}$-point of $\underline{\operatorname{Hom}}^{\otimes}\left(\omega_{d R, \tau} \otimes_{\bar{E}, \iota_{\infty}} \mathbb{C}, H_{l_{\infty} \tau} \otimes_{E, \iota_{\infty}} \mathbb{C}\right)$. Namely, complex conjugation on complex-analytic spaces induces (see [Patrikis 2016c, Lemma 4.1.24]) natural isomorphisms (without restricting to particular graded pieces for the weight and Hodge filtrations, these are isomorphisms of fiber functors over $\mathbb{C}$ )

$$
\begin{equation*}
\operatorname{gr}^{p}\left(e_{\tau} H_{d R}^{w}(M)\right) \otimes_{\bar{E}, \iota_{\infty}} \mathbb{C} \xrightarrow{\sim} \operatorname{gr}^{w-p}\left(e_{c \tau} H_{d R}^{w}(M) \otimes_{\bar{E}, \iota_{\infty}} \mathbb{C}\right), \tag{1}
\end{equation*}
$$

where $c \in \operatorname{Aut}(\bar{E})$ is the choice of complex conjugation for which $\overline{\iota_{\infty}} \tau=\iota_{\infty} c \tau$. We deduce the following relation:
Lemma 2.3. For any embedding $\tau: F \hookrightarrow \bar{E}$, and any choice of complex conjugation $c \in \operatorname{Aut}(\bar{E})$, the conjugacy classes of cocharacters $\left[\mu_{\tau}\right]$ and $\left[\mu_{c \tau}\right]$ satisfy

$$
\left[\mu_{\tau}\right]=\omega \cdot\left[\mu_{c \tau}^{-1}\right]
$$

where $\omega$ is the weight cocharacter.
Proof. For the choice of complex conjugation specified by $\overline{\iota_{\infty}} \tau=\iota_{\infty} c \tau$, the relation $\left[\mu_{\tau}\right]=\omega \cdot\left[\mu_{c \tau}^{-1}\right]$ follows, after base extension $\iota_{\infty}: \bar{E} \rightarrow \mathbb{C}$, from (1) above; but this relation necessarily descends to $\bar{E}$, since the conjugacy classes of cocharacters are defined over any algebraically closed subfield of $\mathbb{C}$. It only remains to observe that [ $\mu_{c \tau}$ ] is independent of the choice of complex conjugation on $\bar{E}$. This follows from the second assertion of Lemma 2.2.

The comparison isomorphisms of $p$-adic Hodge theory then imply that the analogue of Lemma 2.3 also holds for the associated Hodge-Tate cocharacters. For any place $\lambda$ of $E$, fix an algebraic closure $\bar{E}_{\lambda}$. Embeddings $\tau: F \hookrightarrow \bar{E}$ and $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$ then induce $\tau_{\iota_{\lambda}}: F_{v} \hookrightarrow \bar{E}_{\lambda}$ for a suitable place $v$ of $F$ of the same residue characteristic $p$ as $\lambda$. Meanwhile, the restriction to $\Gamma_{F_{v}}$ of the $\lambda$-adic realization induces

$$
\begin{equation*}
\mathcal{M}_{F, E} \xrightarrow{\left.H_{\lambda}\right|_{\Gamma_{F_{v}}}} \operatorname{Rep}_{E_{\lambda}}^{d R}\left(\Gamma_{F_{v}}\right) \xrightarrow{D_{d R}} \operatorname{Fil}_{F_{v} \otimes_{\mathbb{Q}_{p}} E_{\lambda}} \xrightarrow{e_{\tau_{\lambda}}} \operatorname{Fil}_{\bar{E}_{\lambda}} \xrightarrow{\mathrm{gr}} \operatorname{Gr}_{\bar{E}_{\lambda}} \rightarrow \operatorname{Vect}_{\bar{E}_{\lambda}}, \tag{2}
\end{equation*}
$$

where $D_{d R}: \operatorname{Rep}_{E_{\lambda}}^{d R}\left(\Gamma_{F_{v}}\right) \rightarrow \operatorname{Fil}_{F_{v} \otimes_{\mathbb{Q}_{p}} E_{\lambda}}$ denotes Fontaine's functor restricted to the category of de Rham representations. (Here we have invoked Faltings' $p$-adic de Rham comparison isomorphism from [Faltings 1989], and the fact - already noted by André [1996b] - that it extends to a comparison isomorphism on all of $\mathcal{M}_{F, E}$;
for details of the latter point, see [Patrikis 2016c, Lemma 4.1.25].) Of course, $\operatorname{Rep}_{E_{\lambda}}^{d R}\left(\Gamma_{F_{v}}\right)$ also has its standard forgetful fiber functor (let us say $\bar{E}_{\lambda}$-valued), yielding a Tannakian group $\Gamma_{v, \lambda}^{d R}$ for de Rham $\Gamma_{F_{v}}$-representations over $\bar{E}_{\lambda}$; by choosing an isomorphism between the two $\bar{E}_{\lambda}$-valued fiber functors on $\operatorname{Rep}_{E_{\lambda}}^{d R}\left(\Gamma_{F_{v}}\right)$, we obtain a canonical conjugacy class (recalling [Deligne and Milne 1982, Theorem 3.2]) of " $\tau_{\iota_{\lambda}}$-labeled Hodge-Tate cocharacters" [ $\mu_{\tau_{i_{\lambda}}}$ ] of $\Gamma_{v, \lambda}^{d R}$. Specializing, this construction defines the labeled Hodge-Tate cocharacters of any de Rham Galois representation $\rho: \Gamma_{F_{v}} \rightarrow H\left(E_{\lambda}\right)$, for any affine algebraic group $H$ over $E_{\lambda}$.

To relate the $\tau_{\iota_{\lambda}}$-labeled Hodge-Tate cocharacters in the motivic setting to the Hodge cocharacters previously discussed, note that the de Rham comparison isomorphism [Faltings 1989] yields a natural isomorphism of tensor functors $\mathcal{M}_{F, E} \rightarrow \mathrm{Gr}_{\bar{E}_{\lambda}}:$

$$
\operatorname{gr}\left(e_{\tau}\left(H_{d R}(M) \otimes_{E} \bar{E}\right) \otimes_{\bar{E}, \iota_{\lambda}} \bar{E}_{\lambda}\right) \cong \operatorname{gr}\left(e_{\tau_{\iota_{\lambda}}}\left(D_{d R}\left(\left.H_{\lambda}(M)\right|_{\Gamma_{F_{v}}}\right) \otimes_{E_{\lambda}} \bar{E}_{\lambda}\right)\right)
$$

We deduce the following corollary.
Corollary 2.4. For any embedding $\tau: F \hookrightarrow \bar{E}$, and any embedding $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$, there is an equality of conjugacy classes

$$
\left[\mu_{\tau} \otimes_{\bar{E}, \iota_{\lambda}} \bar{E}_{\lambda}\right]=\left[\mu_{\tau_{\iota_{\lambda}}}\right]
$$

In particular, for all $\lambda$, and for all E-embeddings $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$, the conjugacy classes $\left[\mu_{\tau_{\iota_{\lambda}}}\right]$ are independent of $\left(\lambda, \iota_{\lambda}\right)$ when regarded as valued in the common group $\mathcal{G}_{\bar{E}}$.

## 3. Lifting

In this section we prove our main results. First we recall the setting and some notational conventions that will be in effect for the rest of the paper.

Let $F$ be a totally imaginary field (see Remark 1.4), let $S$ be a finite set of places of $F$ containing the infinite places, and let $E$ be any number field. Fix an algebraic closure $\bar{F}$, and set $\Gamma_{F}=\operatorname{Gal}(\bar{F} / F)$. We write $F(S)$ for the maximal extension of $F$ inside $\bar{F}$ that is unramified outside of $S$, and we set $\Gamma_{F, S}=\operatorname{Gal}(F(S) / F)$. We denote the ring of $S$-integers in $F$ by $\mathcal{O}_{F}[1 / S]$. We also set $F_{S}=\prod_{v \in S} F_{v}$. If $L$ is a finite extension of $F$ (inside $\bar{F}$ ), we then abusively continue to write $S$ for the set of all places of $L$ above those in $S$, with corresponding notation $L(S), \Gamma_{L, S}$, etc. For a place $\lambda$ of $E$, let $S_{\lambda}$ denote the set of places of $F$ with the same residue characteristic as $\lambda$. We freely use the terminology established in Definition 1.1.

Torus quotients. In this subsection we prove Theorem 1.3. We begin in the next few paragraphs by gathering together all of the "independent of $\lambda$ and $\rho_{\lambda}$ " data, and the auxiliary constructions we make on top of this data. Continuing with
$F, S$, and $E$ as above, also fix a surjection $\widetilde{H} \rightarrow H$ of linear algebraic groups over $E$ whose kernel is a central torus, which we denote by $C$. (We will without comment also write $C$ for the base change to various algebraically closed fields containing $E$.) Next fix an isogeny complement $H_{1}$ of $C$ in $\widetilde{H}$ (for existence of such $H_{1}$, see [Conrad 2011, Proposition 5.3, Step 1]); thus, $H_{1} \cdot C=\widetilde{H}$, and $H_{1} \cap C$ is finite. For technical reasons, we will later want to include all primes dividing $\#\left(H_{1} \cap C\right)(\bar{E})$ in the set of bad primes $S$; this will be indicated at the necessary point (see the discussion following Lemma 3.5), but it does no harm simply to add these primes to $S$ from now. Consider the quotient map $\xi: \widetilde{H} \rightarrow \widetilde{H} / H_{1} ; \widetilde{Z} \widetilde{Z}^{\vee}=\widetilde{H} / H_{1}$ is a torus, and there is an isogeny $C \rightarrow \widetilde{Z}^{\vee}$, with kernel $C \cap H_{1}$. Fix a split torus $\widetilde{Z}$ over $F$ whose dual group (constructed over $\bar{E}$ ) is isomorphic to $\widetilde{Z}^{\vee} \otimes_{E} \bar{E}$, and fix such an identification (implicit from now on).

Fix a set of cocharacters $\left\{\mu_{\tau}: \boldsymbol{G}_{m, \bar{E}} \rightarrow H_{\bar{E}}\right\}_{\tau: F \hookrightarrow \bar{E}}$, and a central cocharacter $\omega: \boldsymbol{G}_{m, \bar{E}} \rightarrow H_{\bar{E}}$, satisfying the Hodge symmetry requirement of item (2) of Definition 1.1. Denote by $F_{c m}$ the maximal CM subfield of $F$. The condition in Definition 1.1 implies that the cocharacter $\mu_{\tau}$ depends only on the restriction of $\tau$ to $F_{c m}$; we denote this restriction by $\tau_{c m}: F_{c m} \hookrightarrow \bar{E}$. We fix a set of representatives $I$ of $\operatorname{Hom}\left(F_{c m}, \bar{E}\right)$ modulo complex conjugation, and for each $\sigma \in I$, we fix a lift $\widetilde{\mu_{\sigma}}$ to $\widetilde{H}$ of $\mu_{\sigma}$, as well as a (central) lift $\widetilde{\omega}$ of $\omega$. Note that this is possible, because $C$ is a torus. If $\tau: F \hookrightarrow \bar{E}$ restricts to a $\sigma \in I$, we then set $\widetilde{\mu_{\tau}}=\widetilde{\mu_{\sigma}}$; if not, then $c \tau: F \hookrightarrow \bar{E}$ restricts to a $\sigma \in I$, and we then set $\widetilde{\mu_{\tau}}=\widetilde{\omega_{\mu_{\sigma}}}{ }^{-1}$.
Lemma 3.1. Fix once and for all an embedding $\iota_{\infty}: \bar{E} \hookrightarrow \mathbb{C}$. There exists an algebraic automorphic representation $\psi$ of $\widetilde{Z}\left(\boldsymbol{A}_{F}\right)$ such that for all $\tau: F \hookrightarrow \bar{E}$, inducing $\tau_{\iota_{\infty}}: F_{v} \hookrightarrow \mathbb{C}$ by composition with $\iota_{\infty}$, the local component $\psi_{v}: F_{v}^{\times} \rightarrow \mathbb{C}^{\times}$ is given by

$$
\psi_{v}(z)=\tau_{l_{\infty}}(z)^{\xi\left(\widetilde{\mu_{\tau}}\right)} \overline{\tau_{l_{\infty}}}(z)^{\xi\left(\widetilde{\mu_{c \tau}}\right)}
$$

(Recall that $\xi$ is the quotient $\widetilde{H} \rightarrow \widetilde{Z}^{\vee}$.)
Proof. We readily reduce to the case $\widetilde{Z}=\boldsymbol{G}_{m}$, where it follows from the description, due to Weil [1956], of the possible archimedean components of algebraic Hecke characters. (This is where Hodge-symmetry is required.)

From now on we fix such a $\psi$, and we let $T$ denote the finite set of places of $F$ such that $\psi$ is unramified outside $T$. For any embedding $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$, we can then consider the $\lambda$-adic realization ${ }^{3}$

$$
\psi_{\iota_{\lambda}}: \Gamma_{F, T \cup S_{\lambda}} \rightarrow \widetilde{Z}^{\vee}\left(\bar{E}_{\lambda}\right)
$$

Each $\psi_{\iota_{\lambda}}$ is a geometric Galois representation, with good reduction outside $T$, and for any $\tau: F \hookrightarrow \bar{E}$, inducing $\tau_{\iota_{\lambda}}: F_{v} \hookrightarrow \bar{E}_{\lambda}$, the Hodge-Tate cocharacter of $\psi_{\iota_{\lambda}}$ associated to $\tau_{\iota_{\lambda}}$ is $\xi\left(\widetilde{\mu_{\tau}}\right) \otimes_{\bar{E}, \iota_{\lambda}} \bar{E}_{\lambda}$.

[^21]Now we consider any geometric representation

$$
\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)
$$

having good reduction outside $S$, along with an embedding $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$ such that the Hodge-Tate cocharacters of $\rho_{\lambda}$ arise from the collection $\left\{\mu_{\tau}: \boldsymbol{G}_{m, \bar{E}} \rightarrow H_{\bar{E}}\right\}_{\tau: F \hookrightarrow \bar{E}}$ via $\iota_{\lambda}$. Because the kernel of $\widetilde{H} \rightarrow H$ is a central torus, a fundamental theorem of Tate (see [Serre 1977, §6]) ensures in this case that $\rho_{\lambda}$, as a representation of $\Gamma_{F}$, lifts to $\widetilde{H}$. As we will see, our arguments in fact imply Tate's theorem (Corollary 3.9), so we do not need to assume it in what follows.

We can define an obstruction class $\mathcal{O}\left(\rho_{\lambda}\right)$ to lifting $\rho_{\lambda}$ to a continuous representation $\Gamma_{F, S \cup S_{\lambda}} \rightarrow H_{1}\left(\bar{E}_{\lambda}\right)$ in the usual way: choose a topological (but not grouptheoretic) lift $\rho_{\lambda}^{\prime}$, and then form the 2-cocycle $(g, h) \mapsto \rho_{\lambda}^{\prime}(g h) \rho_{\lambda}^{\prime}(h)^{-1} \rho_{\lambda}^{\prime}(g)^{-1}$, defining

$$
\mathcal{O}\left(\rho_{\lambda}\right) \in H^{2}\left(\Gamma_{F, S \cup S_{\lambda}}, H_{1} \cap C\right)
$$

Here and in what follows, we simply write $H_{1} \cap C$ for the $\bar{E}_{\lambda}$-points of this finite group scheme.
Remark 3.2. Here lies the essential difficulty to be overcome: while Tate's theorem allows us to annihilate the cohomology classes $\mathcal{O}\left(\rho_{\lambda}\right)$ - after allowing some additional ramification and enlarging the subgroup $H_{1} \cap C$ of $C$ — we have to carry out this annihilation in a way that is independent of $\lambda$, and moreover, for fixed $\lambda$ independent of $\rho_{\lambda}$. Simultaneous annihilation of the $\mathcal{O}\left(\rho_{\lambda}\right)$ using only a uniform, finite enlargement of the allowable ramification set and of the subgroup of $C$ in fact does not seem to be possible; we will as a first step have to define modified versions of these obstruction classes that take into account the Hodge numbers of $\psi$.

Before proceeding, we reinterpret the obstruction $\mathcal{O}\left(\rho_{\lambda}\right)$ (we will only use the local version of what follows; in particular, the arguments of the present section depend only on the local version of Tate's theorem, which is an almost immediate consequence of local duality).
Lemma 3.3. Let $v$ be a finite place of $F$, and suppose that $\tilde{\rho}_{\lambda}: \Gamma_{F_{v}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ is any continuous homomorphism lifting $\left.\rho_{\lambda}\right|_{\Gamma_{F_{v}}}$. Then $\left.\mathcal{O}\left(\rho_{\lambda}\right)\right|_{\Gamma_{F_{v}}}$ is equal to the inverse of $\mathcal{O}\left(\xi\left(\tilde{\rho}_{\lambda}\right)\right)$, the obstruction associated to lifting $\xi\left(\tilde{\rho}_{\lambda}\right): \Gamma_{F_{v}} \rightarrow \widetilde{Z}^{\vee}\left(\bar{E}_{\lambda}\right)$ to $C$. (The same holds if we replace $\Gamma_{F_{v}}$ by $\Gamma_{F}$, but we do not require this.)
Proof. Before beginning the proof proper, we make precise our convention for coboundary maps: the inverse appearing in the conclusion of the lemma is crucial, and it is easy to get confused if one is not careful with the definitions. Let $\Gamma$ be a group and $M$ a (for simplicity) trivial $\Gamma$-module. For a function $\alpha: \Gamma^{n} \rightarrow M$, set

$$
\begin{aligned}
& \delta(\alpha)\left(g_{1}, \ldots, g_{n+1}\right)=\alpha\left(g_{2}, \ldots, g_{n}\right)+\sum_{i=1}^{n}(-1)^{i} \alpha\left(g_{1}, \ldots, g_{i} g_{i+1}, \ldots, g_{n+1}\right) \\
&+(-1)^{n+1} \alpha\left(g_{1}, \ldots, g_{n}\right)
\end{aligned}
$$

For $n=1$, this says $\delta(\alpha)(g, h)=\alpha(h) \alpha(g h)^{-1} \alpha(g)$, and in the situation considered above $(g, h) \mapsto \rho_{\lambda}^{\prime}(g h) \rho_{\lambda}^{\prime}(h)^{-1} \rho_{\lambda}^{\prime}(g)^{-1}$ is in fact a 2-cocycle.

Tate's theorem implies that, for sufficiently large $m$, the image of $\mathcal{O}\left(\rho_{\lambda}\right)$ in $H^{2}\left(\Gamma_{F_{v}}, C[m]\right)$ vanishes, i.e., $\mathcal{O}\left(\rho_{\lambda}\right)=\delta(\phi)$ for some $\phi: \Gamma_{F_{v}} \rightarrow C[m]$. The product $\rho_{\lambda}^{\prime} \cdot \phi$ is then a homomorphism $\Gamma_{F_{v}} \rightarrow\left(H_{1} \cdot C[m]\right)\left(\bar{E}_{\lambda}\right)$ lifting $\rho_{\lambda}$; we set $\tilde{\rho}_{\lambda}=\rho_{\lambda}^{\prime} \cdot \phi$. Clearly $\xi\left(\tilde{\rho}_{\lambda}\right)=\xi(\phi)$, and then $\mathcal{O}\left(\xi\left(\tilde{\rho}_{\lambda}\right)\right)$ is (tautologically) represented by the cocycle $(g, h) \mapsto \phi(g h) \phi(h)^{-1} \phi(g)^{-1}$, i.e., by $\delta(\phi)^{-1}=\mathcal{O}\left(\rho_{\lambda}\right)^{-1} \in Z^{2}\left(\Gamma_{F}, H_{1} \cap C\right) .{ }^{4}$ This proves the claim for our particular lift $\tilde{\rho}_{\lambda}$, but any other lift $\tilde{\rho}_{\lambda}^{1}$ gives rise to the same obstruction $\mathcal{O}\left(\xi\left(\tilde{\rho}_{\lambda}^{1}\right)\right.$ ). (The global claim holds for the same reasons, if we admit the global version of Tate's theorem.)

To address the difficulty indicated in Remark 3.2, we begin by using the abelian representations coming from $\psi$ to construct a second obstruction class. Namely, consider the realization $\psi_{\iota_{\lambda}}$, which, for notational simplicity, from now on we simply denote by $\psi_{\lambda}$. The automorphic representation $\psi$ is unramified outside the finite set of places $T$ of $F$, so $\psi_{\lambda}$ is a geometric representation $\Gamma_{F, T \cup S_{\lambda}} \rightarrow \widetilde{Z}^{\vee}\left(\bar{E}_{\lambda}\right)$, which has good reduction outside $T$ (i.e., is crystalline at primes of $S_{\lambda}$ not in $T$ ). Via the isogeny $C \rightarrow \widetilde{Z}^{\vee}$, we can then form a cohomology class measuring the obstruction to lifting $\psi_{\lambda}$ to $C$ : let $\psi_{\lambda}^{\prime}$ denote a topological lift $\Gamma_{F, T \cup S_{\lambda}} \rightarrow C\left(\bar{E}_{\lambda}\right)$, defining as before a cohomology class

$$
\mathcal{O}\left(\psi_{\lambda}\right) \in H^{2}\left(\Gamma_{F, T \cup S_{\lambda}}, H_{1} \cap C\right)
$$

We can in turn define (via inflation) a cohomology class

$$
\mathcal{O}\left(\rho_{\lambda}, \psi_{\lambda}\right)=\mathcal{O}\left(\rho_{\lambda}\right) \cdot \mathcal{O}\left(\psi_{\lambda}\right) \in H^{2}\left(\Gamma_{F, S \cup T \cup S_{\lambda}}, H_{1} \cap C\right)
$$

which is represented by the 2-cocycle (recall that $C$ is central in $\widetilde{H}$ )

$$
(g, h) \mapsto\left(\rho_{\lambda}^{\prime} \cdot \psi_{\lambda}^{\prime}\right)(g h)\left(\rho_{\lambda}^{\prime} \cdot \psi_{\lambda}^{\prime}\right)(h)^{-1}\left(\rho_{\lambda}^{\prime} \cdot \psi_{\lambda}^{\prime}\right)(g)^{-1}
$$

(Note, however, that the function $g \mapsto\left(\rho_{\lambda}^{\prime} \cdot \psi_{\lambda}^{\prime}\right)(g)$ is valued in $\widetilde{H}$, not in $H_{1}$.)
We need one more lemma before getting to the crucial local result (Lemma 3.5).
Lemma 3.4. For all places $v \in S_{\lambda}$, and for any choice of embedding $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$, there exists a de Rham lift

of $\left.\rho_{\lambda}\right|_{\Gamma_{v}}$ such that for all embeddings $\tau_{\lambda}: F_{v} \hookrightarrow \bar{E}_{\lambda}$, the $\tau_{\lambda}$-labeled Hodge-Tate

[^22]cocharacter of $\tilde{\rho}_{\lambda}$ is (conjugate to) $\widetilde{\mu_{\tau}} \otimes_{\bar{E}, \iota_{\lambda}} \bar{E}_{\lambda}$, where $\tau: F \hookrightarrow \bar{E}$ is defined by the diagram


Moreover, if $\rho_{\lambda}$ is crystalline, then $\tilde{\rho}_{\lambda}$ may be taken to be crystalline.
Proof. For each $\tau_{\lambda}: F_{v} \hookrightarrow \bar{E}_{\lambda}$, set for notational simplicity $\widetilde{\mu_{\tau_{\lambda}}}=\widetilde{\mu_{\tau}} \otimes_{\bar{E}, l_{\lambda}} \bar{E}_{\lambda}$, where $\tau$ is determined as in the diagram, and where $\widetilde{\mu_{\tau}}$ is the lift of $\mu_{\tau}$ we have fixed above. The proof of [Patrikis 2016c, Corollary 3.2.12] shows that for any collection of cocharacters lifting the Hodge cocharacters of $\rho_{\lambda}$, and in particular for our $\widetilde{\mu_{\tau_{\lambda}}}$, there exists a Hodge-Tate lift $\tilde{\rho}_{\lambda}: \Gamma_{F_{v}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ whose $\tau_{\lambda}$-labeled Hodge-Tate cocharacter is $\widetilde{\mu_{\tau_{\lambda}}}$. Now consider the isogeny lifting problem


Since $\left(\rho_{\lambda}, \xi\left(\tilde{\rho}_{\lambda}\right)\right)$ admits a Hodge-Tate lift (namely, $\tilde{\rho}_{\lambda}$ ), and is itself de Rham ( $\rho_{\lambda}$ is de Rham by assumption, and any abelian Hodge-Tate representation is de Rham), we can apply [Conrad 2011, Corollary 6.7] to deduce the existence of a de Rham lift $\tilde{\rho}_{\lambda}^{\prime}$, which clearly has the same Hodge-Tate cocharacters as $\tilde{\rho}_{\lambda}$, since they differ by a finite-order twist. If we further assume $\rho_{\lambda}$ is crystalline, then we need only a minor modification to this argument: some power $\xi\left(\tilde{\rho}_{\lambda}\right)^{d}$ is crystalline, so if we instead consider the problem of lifting the crystalline representation ( $\rho_{\lambda},[d] \xi\left(\tilde{\rho}_{\lambda}\right)$ ) through the composite isogeny

$$
\widetilde{H} \rightarrow H \times \widetilde{Z}^{\vee} \xrightarrow{\mathrm{id} \times[d]} H \times \widetilde{Z}^{\vee},
$$

then again [Conrad 2011, Corollary 6.7] applies to produce a crystalline lift of $\rho_{\lambda}$ with the desired Hodge-Tate cocharacters.

Here is the key lemma:
Lemma 3.5. For any place $v \in S_{\lambda}$ not belonging to the finite set $S \cup T$, the restriction $\left.\mathcal{O}\left(\rho_{\lambda}, \psi_{\lambda}\right)\right|_{\Gamma_{F_{v}}}$ is trivial.

Proof. Under the assumption on $v$, both $\rho_{\lambda}$ and $\psi_{\lambda}$ are crystalline at $v$. Lemma 3.4 above shows that $\rho_{\lambda} \mid \Gamma_{F_{v}}$ admits a crystalline lift $\tilde{\rho}_{\lambda}: \Gamma_{F_{v}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ such that $\xi\left(\tilde{\rho}_{\lambda}\right)$ has the same (labeled) Hodge-Tate cocharacters as $\left.\psi_{\lambda}\right|_{\Gamma_{F_{v}}}$. Since they are both crystalline, it follows (see [Chai et al. 2014, 3.9.7 Corollary]) that $\left.\xi\left(\tilde{\rho}_{\lambda}\right) \cdot \psi_{\lambda}^{-1}\right|_{\Gamma_{v}}$
is unramified; this is an elaboration of the familiar fact that a crystalline character whose Hodge-Tate weights are zero must be unramified. In particular, replacing the initial lift $\left.\tilde{\rho}_{\lambda}\right|_{\Gamma_{F_{v}}}$ by an unramified twist, we may assume $\xi\left(\tilde{\rho}_{\lambda}\right)=\psi_{\lambda}$ as homomorphisms $\Gamma_{F_{v}} \rightarrow \widetilde{Z}^{\vee}\left(\bar{E}_{\lambda}\right)$. But recall that Lemma 3.3 implies that $\mathcal{O}\left(\rho_{\lambda}\right)=\mathcal{O}\left(\xi\left(\tilde{\rho}_{\lambda}\right)\right)^{-1}$, so we deduce that $\left.\mathcal{O}\left(\rho_{\lambda}\right) \cdot \mathcal{O}\left(\psi_{\lambda}\right)\right|_{\Gamma_{V}}$ is trivial.

Since the set of places $S \cup T$ is finite, by the local version of Tate's theorem, the vanishing of $H^{2}\left(\Gamma_{F_{v}}, \mathbb{Q} / \mathbb{Z}\right)$ for all places $v$ of $F$, we may enlarge $H_{1} \cup C$ to some $C[m]$ inside the torus $C$ so as to kill the image of $H^{2}\left(\Gamma_{F_{v}}, H_{1} \cap C\right) \rightarrow H^{2}\left(\Gamma_{F_{v}}, C[m]\right)$ for all $v \in S \cup T$. (We emphasize that $m$ only depends on the set of places $S \cup T$ of $F$ and the finite group $H_{1} \cap C$.) It follows then from Lemma 3.5 that if $\lambda$ does not belong to $S \cup T$, then $\mathcal{O}\left(\rho_{\lambda}, \psi_{\lambda}\right)$ in fact belongs to

$$
\amalg_{S \cup T \cup S_{\lambda}}^{2}(F, C[m])=\operatorname{ker}\left(H^{2}\left(\Gamma_{F, S \cup T \cup S_{\lambda}}, C[m]\right) \rightarrow \bigoplus_{v \in S \cup T \cup S_{\lambda}} H^{2}\left(\Gamma_{F_{v}}, C[m]\right)\right) .
$$

We can moreover guarantee that this holds regardless of $\lambda$ by an additional finite enlargement of $m$ (since the number of exceptional $\lambda$ is finite). Furthermore, by including the primes dividing $\#\left(H_{1} \cap C\right)$ in $S \cup T$ (if necessary), we can assume that $m$ is divisible only by primes in $S \cup T$. (Note that inflation to allow additional primes of ramification still has image in the corresponding Shafarevich-Tate group, since $\Gamma_{F_{v}} / I_{F_{v}}$ has cohomological dimension one for all finite places $v$.) Thus, after these uniform enlargements of $m$ and $S \cup T$ (which we do not reflect in the notation), we have $\mathcal{O}\left(\rho_{\lambda}, \psi_{\lambda}\right) \in \amalg_{S \cup T \cup S_{\lambda}}^{2}(F, C[m])$.

We are now in a position to apply global duality to analyze the cohomology group $\amalg_{S \cup T \cup S_{\lambda}}^{2}(F, C[m])$. We will need, however, to allow still more primes of ramification in order to kill the class $\mathcal{O}\left(\rho_{\lambda}, \psi_{\lambda}\right)$; the following crucial lemma allows us to do this in a way that does not depend on $\lambda$, but before stating the lemma, we have to recall the Grunwald-Wang theorem (in a somewhat specialized form).

Theorem 3.6 (Grunwald-Wang; see Theorem X. 1 of [Artin and Tate 1968]). Let F be a number field, and let $m$ be a positive integer. Then an element $x \in F^{\times}$belongs to $\left(F^{\times}\right)^{m}$ if and only if $x$ is in $\left(F_{v}^{\times}\right)^{m}$ for all places $v$ of $F$, except when all three of the following conditions, referred to as the special case, hold for the pair $(F, m)$ :

- Let $s_{F}$ denote the largest integer $r$ such that $\eta_{r}=\zeta_{2^{r}}+\zeta_{2^{r}}^{-1}$ is an element of $F$ (here $\zeta_{2^{r}}$ denotes a primitive $2^{r}$-th root of unity). Then $-1,2+\eta_{s_{F}}$, and $-\left(2+\eta_{s_{F}}\right)$ are nonsquares in $F$.
- $\operatorname{ord}_{2}(m)>s_{F}$.
- The set of 2-adic places of $F$ at which $-1,2+\eta_{s_{F}}$, and $-\left(2+\eta_{s_{F}}\right)$ are nonsquares in $F$ is empty.

In the special case, the element $\left(2+\eta_{s_{F}}\right)^{m / 2}$ is the unique (up to $\left(F^{\times}\right)^{m}$-multiple) counterexample to the local-global principle for $m$-th powers in $F^{\times}$.

Here is the lemma:
Lemma 3.7. Recall that $S \cup T$ is a fixed finite set of places of the number field $F$, and that $m$ is a fixed integer. Let $V$ be a finite set of finite places of $F$ such that

- all elements of $V$ are unramified in $F\left(\mu_{m}\right)$,
- the places of $F\left(\mu_{m}\right)$ lying above $V$ generate the class group of $F\left(\mu_{m}\right)$, and
- every element of $\operatorname{Gal}\left(F\left(\mu_{m}\right) / F\right)$ is equal to a (geometric, say) frobenius element at $v$ for some $v \in V$.

Then for all places $\lambda$ of $E$ we can deduce:
(1) If $(F, m)$ is not in the Grunwald-Wang special case, $\amalg_{S \cup T \cup V \cup S_{\lambda}}^{2}(F, C[m])$ is trivial.
(2) If $(F, m)$ is in the Grunwald-Wang special case, then the image of the canonical map

$$
Ш_{S \cup T \cup V \cup S_{\lambda}}^{2}(F, C[m]) \rightarrow Ш_{S \cup T \cup V \cup S_{\lambda}}^{2}(F, C[2 m])
$$

is trivial.
Proof. First note that such sets $V$ exist, by finiteness of the class number and the Čebotarev density theorem. Since (all places of $F$ above) the primes dividing $m$ are contained in $S \cup T$, an application of Poitou-Tate duality immediately reduces us to showing (as a Galois module, $C[m]$ is $\operatorname{dim}(C)$ copies of $\mathbb{Z} / m$ ) the following cases:
(1) If $(F, m)$ is not in the Grunwald-Wang special case, then

$$
\amalg_{S \cup T \cup V \cup S_{\lambda}}^{1}\left(F, \mu_{m}\right)=0
$$

(2) If $(F, m)$ is in the special case, then the map

$$
\amalg_{S \cup T \cup V \cup S_{\lambda}}^{1}\left(F, \mu_{2 m}\right) \rightarrow \amalg_{S \cup T \cup V \cup S_{\lambda}}^{1}\left(F, \mu_{m}\right)
$$

induced by $\mu_{2 m} \xrightarrow{2} \mu_{m}$ is zero.
We first restrict to $\Gamma_{F\left(\mu_{m}\right), S \cup T \cup V \cup S_{\lambda}}$ (note that this is actually restriction to a subgroup, since $F\left(\mu_{m}\right) / F$ is ramified only at primes in $S \cup T$ ), obtaining an element of $\amalg_{S \cup T \cup V \cup S_{\lambda}}^{1}\left(F\left(\mu_{m}\right), \mu_{m}\right)$. After this restriction, as we will see, the GrunwaldWang theorem does not intervene.

To lighten the notation in the rest of the proof, we define $L=F\left(\mu_{m}\right)$ and $Q_{\lambda}=S \cup T \cup V \cup S_{\lambda}$. We also refer the reader to the notation established at the beginning of Section 3. Recall that $F\left(Q_{\lambda}\right)$ denotes the maximal extension of $F$ inside $\bar{F}$ that is unramified outside $Q_{\lambda}$; it contains $L$. Let $\mathcal{O}_{F\left(Q_{\lambda}\right)}$ denote the ring
of $Q_{\lambda}$-integers in $F\left(Q_{\lambda}\right)$ (i.e., the elements of $F\left(Q_{\lambda}\right)$ that are integral outside of places above $Q_{\lambda}$ ). We then have an exact (Kummer theory) sequence

$$
1 \rightarrow \mu_{m} \rightarrow \mathcal{O}_{F\left(Q_{\lambda}\right)}^{\times} \xrightarrow{m} \mathcal{O}_{F\left(Q_{\lambda}\right)}^{\times} \rightarrow 1,
$$

and the corresponding long exact sequence in $\Gamma_{L, Q_{\lambda}}$-cohomology yields an isomorphism

$$
\mathcal{O}_{L}\left[\frac{1}{Q_{\lambda}}\right]^{\times} /\left(\mathcal{O}_{L}\left[\frac{1}{Q_{\lambda}}\right]^{\times}\right)^{m} \xrightarrow{\sim} H^{1}\left(\Gamma_{L, Q_{\lambda}}, \mu_{m}\right)
$$

critically, surjectivity here follows from the vanishing of $H^{1}\left(\Gamma_{L, Q_{\lambda}}, \mathcal{O}_{F\left(Q_{\lambda}\right)}^{\times}\right)$, which itself is a consequence of the natural isomorphism $\mathrm{Cl}_{Q_{\lambda}}(L) \cong H^{1}\left(\Gamma_{L, Q_{\lambda}}, \mathcal{O}_{F\left(Q_{\lambda}\right)}^{\times}\right)$ [Neukirch et al. 2000, Proposition 8.3.11(ii)] and our assumption that $V$ (and hence $Q_{\lambda}$ ) generates the class group of $L$. Restricting the Kummer theory isomorphism to classes that are locally trivial at each place of $Q_{\lambda}$, we also obtain the isomorphism

$$
\left(\mathcal{O}_{L}\left[\frac{1}{Q_{\lambda}}\right]^{\times} \cap\left(L_{Q_{\lambda}}^{\times}\right)^{m}\right) /\left(\mathcal{O}_{L}\left[\frac{1}{Q_{\lambda}}\right]^{\times}\right)^{m} \xrightarrow{\sim} \amalg_{Q_{\lambda}}^{1}\left(L, \mu_{m}\right)
$$

We claim these groups are trivial. Indeed, take $\alpha \in \mathcal{O}_{L}\left[1 / Q_{\lambda}\right]^{\times} \cap\left(L_{Q_{\lambda}}\right)^{m}$, and consider the (abelian) extension $L\left(\alpha^{1 / m}\right) / L$. Global class field theory yields the reciprocity isomorphism

$$
\mathbb{A}_{L}^{\times} /\left(L^{\times} N_{L\left(\alpha^{1 / m}\right) / L}\left(\mathbb{A}_{L\left(\alpha^{1 / m}\right)}^{\times}\right)\right) \xrightarrow{\sim} \operatorname{Gal}\left(L\left(\alpha^{1 / m}\right) / L\right),
$$

but by assumption the source of this map admits a surjection

$$
\mathbb{A}_{L}^{\times} /\left(L^{\times} L_{\infty}^{\times} L_{Q_{\lambda}}^{\times} \prod_{w \notin Q_{\lambda}} \mathcal{O}_{L_{w}}^{\times}\right) \rightarrow \mathbb{A}_{L}^{\times} /\left(L^{\times} N_{L\left(\alpha^{1 / m}\right) / L}\left(\mathbb{A}_{L\left(\alpha^{1 / m}\right)}^{\times}\right)\right)
$$

(At unramified places, the image of the norm map contains the local units; and at places in $Q_{\lambda}, L\left(\alpha^{1 / m}\right) / L$ is split.) By assumption $\left(\mathrm{Cl}_{Q_{\lambda}}(L)=0\right)$, the source of this surjection is trivial, so $L\left(\alpha^{1 / m}\right)=L$, and we deduce that $\amalg_{Q_{\lambda}}^{1}\left(L, \mu_{m}\right)=0$.

It follows that inflation identifies the group $\amalg_{Q_{\lambda}}^{1}\left(F, \mu_{m}\right)$ with the classes in $H^{1}\left(\operatorname{Gal}(L / F), \mu_{m}\right)$ that are trivial upon restriction to $Q_{\lambda}$. Since every element of $\operatorname{Gal}(L / F)$ is a frobenius element at some prime in $V \subset Q_{\lambda}, \amalg_{Q_{\lambda}}^{1}\left(F, \mu_{m}\right)$ is actually equal to the set of everywhere locally trivial classes

$$
\amalg_{|F|}^{1}\left(F, \mu_{m}\right):=\operatorname{ker}\left(H^{1}\left(\Gamma_{F}, \mu_{m}\right) \rightarrow \prod_{v \in|F|} H^{1}\left(\Gamma_{F_{v}}, \mu_{m}\right)\right)
$$

where $|F|$ denotes the set of all places of $F$. This is precisely the subject of the Grunwald-Wang theorem, and it is zero if $(F, m)$ is not in the special case. Thus, we need only consider the possibility that $(F, m)$ is in the special case, where $\amalg_{|F|}^{1}\left(F, \mu_{m}\right)$ has order two, and a representative of the nontrivial class is the (image under the Kummer map of the) element $\left(2+\eta_{S_{F}}\right)^{m / 2}$ of $\left(F^{\times}\right)^{m / 2}$. This description
holds regardless of $m$, so in particular the nontrivial class of $\amalg_{|F|}^{1}\left(F, \mu_{2 m}\right)$ is represented by $\left(2+\eta_{S_{F}}\right)^{m}$. Its image under $\mu_{2 m} \xrightarrow{2} \mu_{m}$, which via Kummer theory is induced by the identity map $F^{\times} \rightarrow F^{\times}$, is again $\left(2+\eta_{s_{F}}\right)^{m}$, which is now visibly an $m$-th power, completing the proof. ${ }^{5}$

We summarize our conclusion, noting that the value of $m$ in the following corollary may be $2 m$ in the earlier notation:

Corollary 3.8. There is an integer $m$ and a finite set of places $Q \supset S \cup T$, both independent of $\lambda$ and of the choice of $\rho_{\lambda}$ having good reduction outside $S$ and the prescribed Hodge-Tate cocharacters $\left\{\mu_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$, such that the image of $\mathcal{O}\left(\rho_{\lambda}, \psi_{\lambda}\right)$ in $H^{2}\left(\Gamma_{F, Q \cup S_{\lambda}}, C[m]\right)$ is zero.

Before proceeding, it is worth noting that the argument just given yields a novel proof of the global version of Tate's vanishing theorem (taking as input the much easier local theorem); it is also a stronger proof, yielding an explicit upper bound on how much ramification has to be allowed, and how much the coefficients need to be enlarged, in order to kill a cohomology class in $H^{2}\left(\Gamma_{F, V}, \mathbb{Z} / N\right)$ for some finite set of places $V$ and integer $N$.

Corollary 3.9. Let $V$ be a finite set of places of $F$, and let $N$ be an integer. Then the image of $H^{2}\left(\Gamma_{F, V}, \mathbb{Z} / N\right)$ in $H^{2}\left(\Gamma_{F, V \cup W}, \mathbb{Z} / 2 N M\right)$ is trivial, where

- $M$ is large enough that for all $v \in V$, the image of

$$
H^{2}\left(\Gamma_{F_{v}}, \mathbb{Z} / N\right) \rightarrow H^{2}\left(\Gamma_{F_{v}}, \mathbb{Z} / N M\right)
$$

is zero, ${ }^{6}$ and

- once $M$ is fixed as above, $W$ is large enough that
- $V \cup W$ contains (all places above) $2 N M$,
$-\mathrm{Cl}_{V \cup W}\left(F\left(\mu_{N M}\right)\right)=0$, and
- each element of $\operatorname{Gal}\left(F\left(\mu_{N M}\right) / F\right)$ is equal to a frobenius element at $w$ for some $w \in W$.
(The factors of two are only necessary in the Grunwald-Wang special case.) In particular, $H^{2}\left(\Gamma_{F}, \mathbb{Q} / \mathbb{Z}\right)=0$.
Remark 3.10. A different proof of Tate's theorem (without arithmetic duality theorems, but instead relying on a finer study of Hecke characters of $F$ ) is given in [Serre 1977, §6.5]. There Serre remarks that Tate originally proved the vanishing theorem using global duality, but further assuming Leopoldt's conjecture; we have of course circumvented Leopoldt here.

[^23]Now we return to the conclusion of Corollary 3.8. Let $b_{\lambda}: \Gamma_{F, Q \cup S_{\lambda}} \rightarrow C[m]$ be a cochain trivializing $\mathcal{O}\left(\rho_{\lambda}, \psi_{\lambda}\right)$. Then

$$
\tilde{\rho}_{\lambda}=\rho_{\lambda}^{\prime} \cdot \psi_{\lambda}^{\prime} \cdot b_{\lambda}: \Gamma_{F, Q \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)
$$

is a homomorphism lifting $\rho_{\lambda}$. We claim that $\tilde{\rho}_{\lambda}$ is, moreover, de Rham at all places in $S_{\lambda}$. To see this, note that under the isogeny $\widetilde{H} \rightarrow H \times \widetilde{Z}^{\vee}$, $\tilde{\rho}_{\lambda}$ pushes forward to $\left(\rho_{\lambda}, \psi_{\lambda} \xi\left(b_{\lambda}\right)\right)$, the second coordinate being a finite-order twist of $\psi_{\lambda}$ (and in particular, de Rham). But now we can invoke the local results of Wintenberger [1995, §1] and Conrad [2011, Theorem 6.2], asserting that a lift of a de Rham representation through an isogeny is de Rham if and only if the Hodge-Tate cocharacter lifts through the isogeny (which is obviously the case here, as $\psi$ was constructed to ensure this).

Finally, we can refine this to the statement that $\rho_{\lambda}$ admits a geometric lift that is moreover crystalline at all places of $S_{\lambda}$, provided $S_{\lambda}$ does not intersect a certain finite set of primes that is independent of $\lambda$ and $\rho_{\lambda}$ (but somewhat larger than the set $Q$ we have thus far constructed). This will complete the proof of Theorem 1.3.

Proof of Theorem 1.3. We resume the above discussion. So far we have a constructed geometric lifts $\tilde{\rho}_{\lambda}: \Gamma_{F, Q \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$, where $Q$ contains $S \cup T$ and whatever other additional primes are needed for the conclusion of Corollary 3.8. The only remaining task is to show that for some (independent of $\rho_{\lambda}$ ) set $P$, we can modify the initial lift (by a finite-order twist) to guarantee that it has good reduction outside $P$. Under the isogeny $\widetilde{H} \rightarrow H \times \widetilde{Z}^{\vee} / \xi(C[m]), \tilde{\rho}_{\lambda}$ pushes forward to

$$
\tau_{\lambda}:=\left(\rho_{\lambda}, \psi_{\lambda} \bmod \xi(C[m])\right)
$$

which is crystalline for all $v$ in $S_{\lambda}$ but not in $S \cup T$. For all $v \in S_{\lambda} \backslash\left(S_{\lambda} \cap(S \cup T)\right)$, $\left.\tilde{\rho}_{\lambda}\right|_{\Gamma_{v}}$ is of course a de Rham lift of $\tau_{\lambda}$, so [Conrad 2011, Theorem 6.2 and Corollary 6.7] (building on [Wintenberger 1995]) shows that $\left.\tau_{\lambda}\right|_{\Gamma_{F_{v}}}$ admits some crystalline lift $\tilde{\tau}_{\lambda, v}: \Gamma_{F_{v}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$, and therefore there are finite-order characters $\chi_{\lambda, v}: \Gamma_{F_{v}} \rightarrow C[m]$ such that each $\left.\tilde{\rho}_{\lambda}\right|_{\Gamma_{F_{v}}} \cdot \chi_{\lambda, v}$ is crystalline. We wish to glue the inertial restrictions $\left.\chi_{\lambda, v}\right|_{I_{F_{v}}}$ together into a global character, with an independentof $-\lambda$ control on the ramification. The cokernel of the restriction map

$$
\begin{equation*}
\operatorname{Hom}\left(\Gamma_{F, Q \cup S_{\lambda}}, C[m]\right) \rightarrow \bigoplus_{v \in S_{\lambda}} \operatorname{Hom}\left(I_{F_{v}}, C[m]\right)^{\Gamma_{F_{v}} / I_{F_{v}}} \tag{3}
\end{equation*}
$$

may be nontrivial; ${ }^{7}$ but we will show that any element of the cokernel is annihilated by appropriate enlargements of $Q$ and $m$.

[^24]By the congruence subgroup property for $\mathrm{GL}_{1}$ (a theorem of Chevalley [1951]), there is an ideal $\mathfrak{n}$ of $\mathcal{O}_{F}$ such that

$$
\left\{x \in \mathcal{O}_{F}^{\times}: x \equiv 1(\bmod \mathfrak{n})\right\} \subseteq\left(\mathcal{O}_{F}^{\times}\right)^{m}
$$

Let $R$ be the set of primes supporting $\mathfrak{n}$ (note that $\mathfrak{n}$ and $R$ are independent of $\rho_{\lambda}!$ ), and set

$$
U_{R}=\left\{\left(x_{v}\right)_{v \in R} \in \prod_{v \in R} \mathcal{O}_{F_{v}}^{\times}: x_{v} \equiv 1(\bmod \mathfrak{n}) \text { for all } v \in R\right\}
$$

Then whenever $S_{\lambda} \cap R=\varnothing$ (so, excluding a finite number of bad $\lambda$ ), consider the character (here and in what follows, we suppress the class field theory identifications)

$$
\left(\chi_{\lambda, v}\right)_{v \in S_{\lambda}} \times 1 \times 1 \times 1: \prod_{v \in S_{\lambda}} \mathcal{O}_{F_{v}}^{\times} \times \prod_{v \in R} U_{R} \times \prod_{v \notin R \cup S_{\lambda}} \mathcal{O}_{F_{v}}^{\times} \times F_{\infty}^{\times} \rightarrow C[m]
$$

which extends by 1 to a character

$$
\left(\prod_{v \in S_{\lambda}} \mathcal{O}_{F_{v}}^{\times} \times \prod_{v \in R} U_{R} \times \prod_{v \notin R \cup S_{\lambda}} \mathcal{O}_{F_{v}}^{\times} \times F_{\infty}^{\times}\right) \cdot F^{\times} \rightarrow C[m]
$$

(an element of the intersection is a global unit congruent to 1 modulo $\mathfrak{n}$, hence is contained in $\left(\mathcal{O}_{F}^{\times}\right)^{m}$, where $\chi_{\lambda, v}$ is obviously trivial). We can then extend from this finite-index subgroup of $\mathbb{A}_{F}^{\times}$to a character $\chi_{\lambda}: \mathbb{A}_{F}^{\times} / F^{\times} \rightarrow \mu_{\infty}(C)$. In fact, we see that $\chi_{\lambda}$ can be chosen to be valued in $C[M]$ for $m$ sufficiently large but independent of $\lambda: m$ can be quantified in terms of the generalized class group of level $U_{R}$, but the details do not concern us.

Replacing $\tilde{\rho}_{\lambda}$ by its finite-order twist

$$
\tilde{\rho}_{\lambda} \cdot \chi_{\lambda}: \Gamma_{F, Q \cup R \cup S_{\lambda}} \rightarrow \tilde{H}\left(\bar{E}_{\lambda}\right)
$$

we have achieved geometric lifts of $\rho_{\lambda}$ with compatible Hodge-Tate cocharacters, and which are crystalline at all places in $S_{\lambda}$ outside of $R \cup S \cup T$.
Remark 3.11. Contrast the final step [Wintenberger 1995, §2.3.5] of Wintenberger's main theorem, where to ensure crystallinity of the lifts he makes a further finite base change on $F$ (having already made several such in order to show lifts exist, as is necessary in his isogeny setup), adding appropriate roots of unity and then passing to a Hilbert class field to kill a cokernel analogous to that of (3). As elsewhere, our argument is orthogonal to Wintenberger's, in allowing additional ramification and larger coefficients, rather than passing to a finite extension of $F$.

We now deduce some corollaries on finding lifts of ramification-compatible systems whose "similitude characters" (determinant, Clifford norm, etc.) form strongly compatible systems, in the sense that at all finite places their associated Weil group representations are isomorphic (see, e.g., [Barnet-Lamb et al. 2014,
§5.1], where these are called strictly compatible). As with Theorem 1.3 and the preceding results, we show a somewhat stronger finiteness result, which applies to all representations with good reduction outside a fixed finite set $S$. These corollaries will follow from the above results and the Hermite-Minkowski finiteness theorem.

Corollary 3.12. Let $F, S$, and $\left\{\mu_{\tau}\right\}$ be as in the statement of Theorem 1.3, except now $F$ may be any number field. Then there exist a finite set of places $P \supset S$ and a finite extension $F^{\prime} / F$ such that any geometric $\rho_{\lambda}$ with good reduction outside $S$, and with Hodge-Tate cocharacters arising from $\left\{\mu_{\tau}\right\}$ via an embedding $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$, admits a geometric lift $\tilde{\rho}_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ such that the restrictions

$$
\xi\left(\tilde{\rho}_{\lambda}\right): \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow \widetilde{Z}^{\vee}\left(\bar{E}_{\lambda}\right)
$$

are equal to the $\iota_{\lambda}$-adic realizations of the single (independent of $\lambda$ and $\rho_{\lambda}$ ) algebraic Hecke character $\psi$ of $\widetilde{Z}\left(\mathbb{A}_{F}\right)$.

In particular, let $\left\{\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)\right\}_{\lambda}$ be a ramification-compatible system. Then there exist a ramification-compatible system of lifts $\left\{\tilde{\rho}_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)\right\}_{\lambda}$, and a finite, independent-of- $\lambda$ extension $F^{\prime} / F$ such that the restrictions

$$
\xi\left(\tilde{\rho}_{\lambda}\right): \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow \widetilde{Z}^{\vee}\left(\bar{E}_{\lambda}\right)
$$

form a strongly compatible system.
Proof. We may assume that the number field $F$ is totally imaginary. Consider the lifts $\tilde{\rho}_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ produced by Theorem 1.3. We write $\xi\left(\tilde{\rho}_{\lambda}\right)=\psi_{\lambda} \cdot \eta_{\lambda}$, where $\eta_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{Z}^{\vee}[M]$ is a finite-order character; the independent-of- $\lambda$ bound on the order was established within the proof of Theorem 1.3. Moreover, for all $v \in S_{\lambda} \backslash\left(S_{\lambda} \cap P\right), \tilde{\rho}_{\lambda}$ and $\psi_{\lambda}$ are crystalline at $v$, so as long as $S_{\lambda} \cap P$ is empty, $\eta_{\lambda}$ factors through $\Gamma_{F, P} \rightarrow \widetilde{Z}^{\vee}[M]$ (we again use that a finite-order crystalline character is unramified). By the Hermite-Minkowski theorem, there are a finite number of such characters $\eta_{\lambda}$. For the finite number of bad $\lambda$ (at which $S_{\lambda} \cap P \neq \varnothing$ ), the same finiteness assertion holds. Thus, after a finite base change $F^{\prime} / F$, trivializing this finite collection of possible characters $\eta_{\lambda}$, we see that $\left.\xi\left(\tilde{\rho}_{\lambda}\right)\right|_{\Gamma_{F^{\prime}, P \cup S_{\lambda}}}=\left.\psi_{\lambda}\right|_{\Gamma_{F^{\prime}, P \cup S_{\lambda}}}$ for all $\lambda$. The second part of the corollary follows since the $\lambda$-adic realizations of an abelian $L$-algebraic representation form a strongly compatible system, as is evident from the construction of $\psi_{\lambda}$, as in, e.g., [Serre 1968].

We would like to upgrade this to a compatibility statement not just for the pushforwards $\xi\left(\tilde{\rho}_{\lambda}\right)$, but for the full abelianizations $\tilde{\rho}^{\mathrm{ab}}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{H}^{\text {ab }}\left(\bar{E}_{\lambda}\right)$. Of course, such a result requires first (taking $\widetilde{H}=H$ ) having the corresponding assertion for the abelianizations $\rho_{\lambda}^{\mathrm{ab}}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H^{\mathrm{ab}}\left(\bar{E}_{\lambda}\right)$. Here, however, it is of course false without imposing further conditions on the system $\left\{\rho_{\lambda}\right\}_{\lambda}$ (see Remark 3.15). There are various conditions we might impose on the $\rho_{\lambda}$ to ensure (potential) compatibility
of the $\rho_{\lambda}^{\mathrm{ab}}$. Perhaps most interesting is to restrict the coefficients of $\rho_{\lambda}^{\mathrm{ab}}$. To that end, we first prove a finiteness result for Galois characters:

Lemma 3.13. Let $F$ be a number field, and let $S$ be a finite set of places of $F$. Fix a finite extension $E^{\prime} / E$ (inside $\bar{E}$ ), a set $\left\{m_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$ of integers satisfying the Hodge-symmetry condition of Definition 1.1, and an embedding $\iota_{\infty}: \bar{E} \hookrightarrow \mathbb{C}$. Then there exist a finite extension $F^{\prime} / F$, and an algebraic Hecke character $\alpha$ of $\mathbb{A}_{F^{\prime}}$, such that any geometric character $\omega_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow \bar{E}_{\lambda}^{\times}$

- having good reduction outside $S$;
- having labeled Hodge-Tate weights corresponding to $\left\{m_{\tau}\right\}$ via some embed$\operatorname{ding} \iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda} ;$ and
- for which $\omega_{\lambda}\left(f r_{v}\right)$ belongs to $\left(E^{\prime}\right)^{\times}$for a density-one set of places $v$ of $F$;
will upon restriction $\left.\omega_{\lambda}\right|_{\Gamma_{F^{\prime}}}$, become isomorphic to the $\iota_{\lambda}$-adic realization of $\alpha$.
Proof. We may assume $F$ is totally imaginary. Invoking the Hodge symmetry hypothesis, we apply Lemma 3.1 to produce an algebraic Hecke character $\alpha$ of $F$ whose archimedean components are given in terms of the $m_{\tau}$, exactly as in Lemma 3.1 (with $\xi\left(\tilde{\mu}_{\tau}\right)=m_{\tau}$ ). Let $T$ denote the finite set of ramified places of $\alpha$, and let $\mathbb{Q}(\alpha)$ denote the field of coefficients of $\alpha$ (by definition the fixed field of all automorphisms of $\mathbb{C}$ that preserve the nonarchimedean component of $\alpha$; we will regard $\mathbb{Q}(\alpha)$ as a subfield of $\bar{E}$ via our fixed $\left.\iota_{\infty}\right)$. Thus the $\iota_{\lambda}$-adic realizations $\alpha_{\lambda}: \Gamma_{F, T \cup S_{\lambda}} \rightarrow \bar{E}_{\lambda}^{\times}$have labeled Hodge-Tate weights matching those of $\omega_{\lambda}$. Since $\mathbb{Q}(\alpha)$ contains the values $\alpha_{\lambda}\left(f r_{v}\right)$ for all $v \notin T \cup S_{\lambda}$, and $\omega_{\lambda} \alpha_{\lambda}^{-1}: \Gamma_{F, S \cup T \cup S_{\lambda}} \rightarrow \bar{E}_{\lambda}^{\times}$ is finite-order (all of its Hodge-Tate weights are zero), we see that $\left(\omega_{\lambda} \alpha_{\lambda}^{-1}\right)\left(f r_{v}\right)$ belongs to the finite (independent of $\lambda$ ) set $\mu_{\infty}\left(E^{\prime} \mathbb{Q}(\alpha)\right)$ for a density-one set of $v$. By Čebotarev, the character $\omega_{\lambda} \alpha_{\lambda}^{-1}$ takes all of its values in $\mu_{\infty}\left(E^{\prime} \mathbb{Q}(\alpha)\right)$. As long as $S_{\lambda} \cap(S \cup T)$ is empty, $\omega_{\lambda} \alpha_{\lambda}^{-1}$ is moreover unramified at $S_{\lambda}$ (because it is crystalline of finite order), so as in Corollary 3.12, there are (again by Hermite-Minkowski) a finite number of such characters $\omega_{\lambda} \alpha_{\lambda}^{-1}$. We deduce the existence of a single number field $F^{\prime}$ over which $\left.\omega_{\lambda}\right|_{F^{\prime}}=\left.\alpha_{\lambda}\right|_{\Gamma^{\prime}}$, for any $\lambda$ and any $\omega_{\lambda}$ as in the statement of the lemma.

We deduce a potential compatibility statement for the full abelianizations $\tilde{\rho}_{\lambda}^{\mathrm{ab}}$ :
Corollary 3.14. For simplicity, assume that $H^{\mathrm{ab}}$ is of multiplicative type. Let $F$, $S$, and $\left\{\mu_{\tau}\right\}$ be as in the statement of Theorem 1.3, except with $F$ now allowed to be any number field. Also fix a finite extension $E^{\prime}$ of $E$. Then there exist a finite set of places $P \supset S$, a finite extension $F^{\prime} / F$, and an algebraic Hecke character $\beta$ of the split group $\widetilde{D}$ over $F^{\prime}$ whose dual group over $\bar{E}$ is isomorphic to $\left(H^{\mathrm{ab}}\right)^{0} \otimes_{E} \bar{E}$ (and we fix such an isomorphism), satisfying the following: if a geometric $\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)$

- has good reduction outside $S$;
- has Hodge-Tate cocharacters arising from $\left\{\mu_{\tau}\right\}$ via $\iota_{\lambda}: \bar{E} \hookrightarrow \bar{E}_{\lambda}$;
- and admits, for some faithful representation $r$ of $H^{\mathrm{ab}}$, a density-one set of places $v$ of $F$ such that the characteristic polynomial $\operatorname{ch}\left(r \circ \tilde{\rho}_{\lambda}^{\mathrm{ab}}\right)\left(f r_{v}\right)$ has coefficients in $E^{\prime}$;
then there is a geometric lift $\tilde{\rho}_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ having good reduction outside $P$, such that the restriction $\tilde{\rho}_{\lambda}^{\mathrm{ab}}: \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow \widetilde{H}^{\text {ab }}\left(\bar{E}_{\lambda}\right)$ is equal to the $\iota_{\lambda}$-adic realization $\beta_{\lambda}$ of $\beta$.

In particular, let $\left\{\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)\right\}_{\lambda}$ be a ramification-compatible system, and assume that for some faithful representation $r$ of $H^{\mathrm{ab}}$, some number field $E^{\prime}$, and for almost all $\lambda$, there is a density-one set of places $v$ of $F$ such that the characteristic polynomial $\operatorname{ch}\left(r \circ \tilde{\rho}_{\lambda}^{\mathrm{ab}}\right)\left(f r_{v}\right)$ has coefficients in $E^{\prime}$. Then there is a ramification-compatible system $\tilde{\rho}_{\lambda}: \Gamma_{F, P \cup S_{\lambda}} \rightarrow \widetilde{H}\left(\bar{E}_{\lambda}\right)$ lifting $\rho_{\lambda}$, and a finite extension $F^{\prime} / F$ such that

$$
\left.\tilde{\rho}_{\lambda}^{\mathrm{ab}}\right|_{\Gamma_{F^{\prime}}}: \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow \widetilde{H}^{\mathrm{ab}}\left(\bar{E}_{\lambda}\right)
$$

forms a strongly compatible system.
Proof. The proof follows familiar lines. Since $C$ is central, the abelianization $\widetilde{H}^{\text {ab }}$ is simply $\widetilde{H} / H_{1}^{\text {der }}$, so there is a natural map

$$
f: \widetilde{H}^{\mathrm{ab}} \rightarrow \widetilde{H} / H_{1} \times H / \operatorname{im}\left(H_{1}^{\mathrm{der}}\right)=\widetilde{Z}^{\vee} \times H^{\mathrm{ab}}
$$

under which $\tilde{\rho}_{\lambda}^{\mathrm{ab}}$ pushes forward to $\left(\xi\left(\tilde{\rho}_{\lambda}\right), \rho_{\lambda}^{\mathrm{ab}}\right)$. (We have chosen $\left\{\tilde{\rho}_{\lambda}\right\}_{\lambda}$ as in Theorem 1.3 and Corollary 3.12, of course.) First we claim that a conclusion analogous to that of the corollary holds for the pair $\left(\xi\left(\tilde{\rho}_{\lambda}\right), \rho_{\lambda}^{\mathrm{ab}}\right)$, and certainly it suffices to check this independently for the two components. The assertion for $\xi\left(\tilde{\rho}_{\lambda}\right)$ is Corollary 3.12, and for $\rho_{\lambda}^{\mathrm{ab}}$ it follows easily from Lemma 3.13 (first reduce, by a finite base change, to the case where $H^{\text {ab }}$ is connected, using the fact that $\pi_{0}\left(H^{\mathrm{ab}}\right)$ is of course finite and independent of $\left.\lambda\right)$. Thus, letting $D$ denote a split torus whose dual group is identified with $\left(H^{\mathrm{ab}}\right)^{0}$, there exists a finite extension $F_{1} / F$ such that $\left.f\left(\tilde{\rho}_{\lambda}^{\mathrm{ab}}\right)\right|_{\Gamma_{F_{1}, P \cup S_{\lambda}}}$ is the $\iota_{\lambda}$-adic realization of a Hecke character (which does not depend on $\lambda$ or $\rho_{\lambda}$ ) of $\widetilde{Z} \times D$.

Now suppose that $\beta$ is a Hecke character of $\widetilde{D}$ for which the $\iota_{\lambda}$-adic realization $\beta_{\lambda}: \Gamma_{F, T \cup S_{\lambda}} \rightarrow\left(\tilde{H}^{\text {ab }}\right)^{0}\left(\bar{E}_{\lambda}\right)$ has labeled Hodge-Tate cocharacters matching those of $\tilde{\rho}_{\lambda}^{\mathrm{ab}}$. Since

$$
\left.f\left(\tilde{\rho}_{\lambda}^{\mathrm{ab}} \cdot \beta_{\lambda}^{-1}\right)\right|_{\Gamma_{F_{1}, P \cup T \cup S_{\lambda}}}
$$

is automorphic (independently of $\lambda, \rho_{\lambda}$ ) of finite order, it is trivial after a finite base change $F_{2} / F_{1}$. Now observe that the kernel of $f$ is finite, so $\left.\left(\tilde{\rho}_{\lambda}^{\text {ab }} \cdot \beta_{\lambda}^{-1}\right)\right|_{\Gamma_{F_{2}, P \cup T \cup S_{\lambda}}}$ has finite order, bounded only in terms of \# $\operatorname{ker}(f)$, and is crystalline away from $P \cup T$;
as before, we find a further finite extension $F_{3} / F_{2}$ such that $\left.\left(\tilde{\rho}_{\lambda}^{\mathrm{ab}} \cdot \beta_{\lambda}^{-1}\right)\right|_{\Gamma_{F_{3}}}=1$. The conclusion of the corollary then holds with $F^{\prime}=F_{3}$.

Remark 3.15. - It does not suffice to ask for a fixed number field $E$ such that all $\rho_{\lambda}^{\mathrm{ab}}$ are valued in $H^{\mathrm{ab}}\left(E_{\lambda}\right)$. For instance, taking $F=\mathbb{Q}$ and $S=\{p\}$, and for all $n$ choosing a prime $\ell_{n} \equiv 1\left(\bmod \varphi\left(p^{n}\right)\right)$, we can define $\rho_{\ell_{n}}$ : $\Gamma_{\mathbb{Q},\{p\}} \rightarrow \mathbb{Q}_{\ell_{n}}^{\times}$as the composition of the mod $p^{n}$ cyclotomic character with an inclusion $\left(\mathbb{Z} / p^{n} \mathbb{Z}\right)^{\times} \hookrightarrow \mu_{\ell_{n}-1} \hookrightarrow \mathbb{Q}_{\ell_{n}}^{\times}$, and for all $\ell \notin\left\{\ell_{n}\right\}_{n}$ we can take $\rho_{\ell}$ to be the trivial character. Then $\left\{\rho_{\ell}: \Gamma_{\mathbb{Q},\{p\}} \rightarrow \mathbb{Q}_{\ell}^{\times}\right\}_{\ell}$ is an abelian, ramificationcompatible system that does not become a strongly compatible system after any finite base change.

- Having only hypothesized ramification-compatibility for the $\left\{\rho_{\lambda}\right\}_{\lambda}$, we cannot hope for the stronger conclusion that the $\left\{\tilde{\rho}_{\lambda}^{\mathrm{ab}}\right\}_{\lambda}$ form a strongly compatible system over $F$ itself.

General multiplicative-type quotients. In fact, the argument of Theorem 1.3 directly implies the main theorem of [Wintenberger 1995], as well as a generalization to lifting through quotients where the kernel is central of multiplicative type. We thus obtain an essentially different proof (and generalization) of Wintenberger's result. In this section, we briefly describe how this works.

Corollary 3.16 (Wintenberger). Let $H_{1} \rightarrow H$ be a central isogeny of linear algebraic groups over $E$, and let $S$ be a finite set of places of $F$. Then there exist a finite extension $F^{\prime} / F$ and a finite set of places $P \supset S$ of $F$ such that any geometric representation

$$
\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(E_{\lambda}\right)
$$

having

- good reduction outside S, and
- labeled Hodge-Tate cocharacters that lift to $H_{1}$,
lifts to a geometric representation $\rho_{\lambda}^{\prime}: \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow H_{1}\left(E_{\lambda}\right)$, which moreover has good reduction outside $P$.

Proof. We begin by replacing $F$ by a finite extension $F_{0}$ such that image of $\left.\rho_{\lambda}\right|_{\Gamma_{0}}$ is contained in the image of $H_{1}\left(E_{\lambda}\right) \rightarrow H\left(E_{\lambda}\right)$. That such an extension, depending only on $H_{1} \rightarrow H, S$, and $F$, exists follows as in [Wintenberger 1995, 2.3.2], and we do not repeat the argument. We note, though, that making this construction in an independent-of- $\lambda$ fashion already uses liftability of the HodgeTate cocharacters. (If we were not concerned with preserving $E_{\lambda}$-rationality of the lift, then we could skip this step.) It is then possible to build an obstruction class $\mathcal{O}\left(\rho_{\lambda}\right) \in H^{2}\left(\Gamma_{F_{0}, S \cup S_{\lambda}}, \operatorname{ker}\left(H_{1} \rightarrow H\right)\left(E_{\lambda}\right)\right)$ via a topological lift $\rho_{\lambda}^{\prime}$ to $H_{1}\left(E_{\lambda}\right)$.

Embed $\operatorname{ker}\left(H_{1} \rightarrow H\right) \otimes_{E} \bar{E}$ into a torus $C$, and with the kernel embedded antidiagonally, form the new group $\widetilde{H}=\left(H_{1} \times C\right) / \operatorname{ker}\left(H_{1} \rightarrow H\right)$. The surjection $\widetilde{H} \rightarrow H$ now has kernel equal to a central torus $C$, and as before we let $\widetilde{Z}^{\vee}$ be the (torus) quotient $\tilde{H} / H_{1}$. By hypothesis, we can lift the Hodge-Tate cocharacters of $\rho_{\lambda}$ to $H_{1}$; when pushed forward to $\widetilde{Z}^{\vee}$, these lifts are of course trivial. Thus, in the notation of Lemma 3.1, we may take the trivial Hecke character $\psi=1$ of $\widetilde{Z}\left(\mathbb{A}_{F_{0}}\right)$. For topological lifts $\psi_{\lambda}^{\prime}$ to $C\left(E_{\lambda}\right)$ (as in Lemma 3.3) of the (trivial) $\lambda$-adic realizations $\psi_{\lambda}$, we may of course also take $\psi_{\lambda}^{\prime}=1$. Theorem 1.3 then produces a finite set of primes $P \supset S$ and an integer $M$, both only depending on $H_{1} \rightarrow H, S$, and $F$, and a geometric lift $\tilde{\rho}_{\lambda}: \Gamma_{F_{0}, P \cup S_{\lambda}} \rightarrow H_{1}\left(E_{\lambda}\right) \cdot C[M]$ such that $\tilde{\rho}_{\lambda}$ has good reduction outside $P$. (The assertion that $\tilde{\rho}_{\lambda}$ is valued in the subset $H_{1}\left(E_{\lambda}\right) \cdot C[M]$ of $\widetilde{H}\left(\bar{E}_{\lambda}\right)$ follows from the explicit description of $\tilde{\rho}_{\lambda}$, since $\rho_{\lambda}^{\prime}$ lands in $H_{1}\left(E_{\lambda}\right)$, and $\psi_{\lambda}^{\prime}$ is trivial.) For all $\lambda$ for which $S_{\lambda} \cap P=\varnothing, \xi\left(\tilde{\rho}_{\lambda}\right): \Gamma_{F_{0}, P \cup S_{\lambda}} \rightarrow \widetilde{Z}^{\vee}[M]$ is also unramified at $S_{\lambda}$, and all such characters are trivialized by a common finite extension $F_{1} / F_{0}$. For the finite number of $\lambda$ such that $S_{\lambda} \cap P$ is nonempty, we can again trivialize the possible $\xi\left(\tilde{\rho}_{\lambda}\right)$ by restricting to a common finite extension $F_{2} / F_{0}$. Taking $F^{\prime}=F_{1} F_{2}$, all $\left.\tilde{\rho}_{\lambda}\right|_{\Gamma_{F^{\prime}, P \cup S_{\lambda}}}$ land in $H_{1}\left(E_{\lambda}\right)$, proving the corollary.
Remark 3.17. For Wintenberger's result, take $E=\mathbb{Q}$. He also shows [1995, 2.3.6] that there is a second finite extension $F^{\prime \prime} / F^{\prime}$ (only depending on $H_{1} \rightarrow H, F$, and $S$ ) such that any two lifts $\rho_{\lambda}^{\prime}$ as in the corollary become equal after restriction to $\Gamma_{F^{\prime \prime}}$. This refinement similarly follows in our setup, but there is no need to repeat Wintenberger's argument.

Here is the more general version with multiplicative-type kernels. Note that, as with Theorem 1.3, but unlike Corollary 3.16, it makes use of a "Hodge symmetry" hypothesis.

Corollary 3.18. Let $H^{\prime} \rightarrow H$ be a surjection of linear algebraic groups over $E$ whose kernel is central and of multiplicative type. Let $F$ be a number field, and let $S$ be a finite set of places of $F$ containing the archimedean places. Fix a set of cocharacters $\left\{\mu_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$ as in part (2) of Definition 1.1, and moreover, assume that each $\mu_{\tau}$ lifts to a cocharacter of $H^{\prime}$.

Then there exist a finite set of places $P \supset S$, and a finite extension $F^{\prime} / F$, such that any geometric representation $\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)$ having good reduction outside $S$, and whose Hodge-Tate cocharacters arise from the set $\left\{\mu_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$ via some embedding $\bar{E} \hookrightarrow \bar{E}_{\lambda}$, admits a geometric lift $\tilde{\rho}_{\lambda}: \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow H^{\prime}\left(\bar{E}_{\lambda}\right)$ having good reduction outside $P$.

In particular, if $\left\{\rho_{\lambda}: \Gamma_{F, S \cup S_{\lambda}} \rightarrow H\left(\bar{E}_{\lambda}\right)\right\}_{\lambda}$ is a ramification-compatible system with Hodge cocharacter $\left\{\mu_{\tau}\right\}_{\tau: F \hookrightarrow \bar{E}}$, then there exist a finite set of places $P \supset S$, a finite extension $F^{\prime} / F$, and lifts $\tilde{\rho}_{\lambda}: \Gamma_{F^{\prime}, P \cup S_{\lambda}} \rightarrow H^{\prime}\left(\bar{E}_{\lambda}\right)$ such that $\left\{\tilde{\rho}_{\lambda}\right\}_{\lambda}$ is a ramification-compatible system.

Proof. As in the proofs of Theorem 1.3 and Corollary 3.16, we construct an isogeny complement $H_{1} \subset H$ to $\operatorname{ker}\left(H^{\prime} \rightarrow H\right)$, as well as an enlargement $\widetilde{H} \supset H^{\prime}$ surjecting onto $H$ with a central torus kernel. We then run the argument of Theorem 1.3, starting from lifts $\left\{\mu_{\tau}^{\prime}\right\}$ to $H^{\prime}$ of the Hodge cocharacters: the Hecke character $\psi$ (in the notation of that proof) then constructed has $\lambda$-adic realizations that push-forward to finite-order characters $\psi_{\lambda}: \Gamma_{F, T \cup S_{\lambda}} \rightarrow \widetilde{H} / H^{\prime}\left(\bar{E}_{\lambda}\right)$, and from here it is easy to proceed; we omit the details, since the argument will by now be familiar.

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# Remarks on the arithmetic fundamental lemma 

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#### Abstract

W. Zhang's arithmetic fundamental lemma (AFL) is a conjectural identity between the derivative of an orbital integral on a symmetric space with an arithmetic intersection number on a unitary Rapoport-Zink space. In the minuscule case, Rapoport, Terstiege and Zhang have verified the AFL conjecture via explicit evaluation of both sides of the identity. We present a simpler way for evaluating the arithmetic intersection number, thereby providing a new proof of the AFL conjecture in the minuscule case.


## 1. Introduction

1.1. Zhang's arithmetic fundamental lemma. The arithmetic Gan-Gross-Prasad conjectures (arithmetic GGP) generalize the celebrated Gross-Zagier formula to higher-dimensional Shimura varieties [Gan et al. 2012, §27; Zhang 2012, §3.2]. The arithmetic fundamental lemma (AFL) conjecture arises from Zhang's relative trace formula approach for establishing the arithmetic GGP for the group $U(1, n-2) \times$ $U(1, n-1)$. It relates a derivative of orbital integrals on symmetric spaces to an arithmetic intersection number of cycles on unitary Rapoport-Zink spaces,

$$
\begin{equation*}
O^{\prime}\left(\gamma, \mathbf{1}_{S_{n}\left(\mathbb{Z}_{p}\right)}\right)=-\omega(\gamma)\left\langle\Delta\left(\mathcal{N}_{n-1}\right),(\mathrm{id} \times g) \Delta\left(\mathcal{N}_{n-1}\right)\right\rangle . \tag{1.1.0.1}
\end{equation*}
$$

For the precise definitions of quantities appearing in the identity, see [Rapoport et al. 2013, Conjecture 1.2]. The left-hand side of (1.1.0.1) is known as the analytic side and the right-hand side is known as the arithmetic-geometric side. The AFL conjecture has been verified for $n=2,3$ [Zhang 2012], and for general $n$ in the minuscule case (in the sense that $g$ satisfies a certain minuscule condition) by Rapoport, Terstiege and Zhang [2013]. In all these cases, the identity (1.1.0.1) is proved via explicit evaluation of both sides. When $g$ satisfies a certain inductive condition, Mihatsch [2016] has recently developed a recursive algorithm which

[^25]reduces the identity (1.1.0.1) to smaller $n$, thus establishing some new cases of the AFL conjecture.

In the minuscule case, the evaluation of the analytic side is relatively straightforward. The bulk of [Rapoport et al. 2013] is devoted to a highly nontrivial evaluation of the arithmetic-geometric side, which is truly a tour de force. Our main goal in this short note is to present a new (and arguably simpler) way to evaluate the arithmetic-geometric side in [Rapoport et al. 2013], henceforth abbreviated [RTZ].
1.2. The main results. Let $p$ be an odd prime. Let $F=\mathbb{Q}_{p}, k=\overline{\mathbb{F}}_{p}, W=W(k)$ and $K=W[1 / p]$. Let $\sigma$ be the $p$-Frobenius acting on $\overline{\mathbb{F}}_{p}, W$ and $K$. Let $E=\mathbb{Q}_{p^{2}}$ be the unramified quadratic extension of $F$. The unitary Rapoport-Zink space $\mathcal{N}_{n}$ is the formal scheme over $\operatorname{Spf} W$ parametrizing deformations up to quasi-isogeny of height 0 of unitary $p$-divisible groups of signature $(1, n-1)$ (definitions recalled in Section 2.1). Fix $n \geq 2$ and write $\mathcal{N}=\mathcal{N}_{n}$ and $\mathcal{M}=\mathcal{N}_{n-1}$ for short. There is a natural closed immersion $\delta: \mathcal{M} \rightarrow \mathcal{N}$. Denote by $\Delta \subseteq \mathcal{M} \times{ }_{W} \mathcal{N}$ the image of (id, $\delta$ ): $\mathcal{M} \rightarrow \mathcal{M} \times{ }_{W} \mathcal{N}$, known as the (local) diagonal cycle or $G G P$ cycle on $\mathcal{M} \times{ }_{W} \mathcal{N}$.

Let $C_{n-1}$ be a nonsplit $\sigma$-Hermitian $E$-space of dimension $n-1$. Let $C_{n}=$ $C_{n-1} \oplus E u$ (where the direct sum is orthogonal and $u$ has norm 1) be a nonsplit $\sigma$-Hermitian $E$-space of dimension $n$. The unitary group $J=\mathrm{U}\left(C_{n}\right)$ acts on $\mathcal{N}$ in a natural way (see Section 2.2). Let $g \in J\left(\mathbb{Q}_{p}\right)$. The arithmetic-geometric side of the AFL conjecture (1.1.0.1) concerns the arithmetic intersection number of the diagonal cycle $\Delta$ and its translate by id $\times g$, defined as

$$
\langle\Delta,(\mathrm{id} \times g) \Delta\rangle:=\log p \cdot \chi\left(\mathcal{M} \times_{W} \mathcal{N}, \mathcal{O}_{\Delta} \otimes^{\mathbb{L}} \mathcal{O}_{(\mathrm{id} \times g) \Delta}\right)
$$

When $\Delta$ and $(\mathrm{id} \times g) \Delta$ intersect properly, namely when the formal scheme

$$
\begin{equation*}
\Delta \cap(\mathrm{id} \times g) \Delta \cong \delta(\mathcal{M}) \cap \mathcal{N}^{g} \tag{1.2.0.1}
\end{equation*}
$$

is an Artinian scheme (where $\mathcal{N}^{g}$ denotes the fixed points of $g$ ), the intersection number is simply $\log p$ times the $W$-length of the Artinian scheme (1.2.0.1).

Recall that $g \in J\left(\mathbb{Q}_{p}\right)$ is called regular semisimple if

$$
L(g):=\mathcal{O}_{E} \cdot u+\mathcal{O}_{E} \cdot g u+\cdots+\mathcal{O}_{E} \cdot g^{n-1} u
$$

is an $\mathcal{O}_{E}$-lattice in $C_{n}$. In this case, the invariant of $g$ is the unique sequence of integers

$$
\operatorname{inv}(g):=\left(r_{1} \geq r_{2} \geq \cdots \geq r_{n}\right)
$$

characterized by the condition that there exists a basis $\left\{e_{i}\right\}$ of the lattice $L(g)$ such that $\left\{p^{-r_{i}} e_{i}\right\}$ is a basis of the dual lattice $L(g)^{\vee}$. It turns out that the "bigger" $\operatorname{inv}(g)$ is, the more difficult it is to compute the intersection. With this in mind, recall that a regular semisimple element $g$ is called minuscule if $r_{1}=1$ and $r_{n} \geq 0$ (equivalently, $\left.p L(g)^{\vee} \subseteq L(g) \subseteq L(g)^{\vee}\right)$. In this minuscule case, the intersection turns out to be
proper, and one of the main results of [RTZ] is an explicit formula for the $W$-length of (1.2.0.1) at each of its $k$-point.

To state the formula, assume $g$ is regular semisimple and minuscule, and suppose $\mathcal{N}^{g}$ is nonempty. Then $g$ stabilizes both $L(g)^{\vee}$ and $L(g)$ and thus acts on the $\mathbb{F}_{p^{2}}$-vector space $L(g)^{\vee} / L(g)$. Let $P(T) \in \mathbb{F}_{p^{2}}[T]$ be the characteristic polynomial of $g$ acting on $L(g)^{\vee} / L(g)$. For any irreducible polynomial $R(T) \in \mathbb{F}_{p^{2}}[T]$, we denote its multiplicity in $P(T)$ by $m(R(T))$ and define its reciprocal by

$$
R^{*}(T):=T^{\operatorname{deg} R(T)} \cdot \sigma\left(R\left(\frac{1}{T}\right)\right)
$$

We say $R(T)$ is self-reciprocal if $R(T)=R^{*}(T)$. By [RTZ, 8.1], if $\left(\delta(M) \cap \mathcal{N}^{g}\right)(k)$ is nonempty, then $P(T)$ has a unique self-reciprocal monic irreducible factor $Q(T) \mid P(T)$ such that $m(Q(T))$ is odd. We denote

$$
c:=\frac{1}{2}(m(Q(T))+1) .
$$

Then $1 \leq c \leq \frac{1}{2}(n+1)$. Now we are ready to state the intersection length formula.
Theorem A [RTZ, Theorem 9.5]. Assume $g$ is regular semisimple and minuscule. Assume $p>c$. Then for any $x \in\left(\delta(M) \cap \mathcal{N}^{g}\right)(k)$, the complete local ring of $\delta(M) \cap \mathcal{N}^{g}$ at $x$ is isomorphic to $k[X] / X^{c}$, and hence has $W$-length equal to $c$.

We will present a simpler proof of Theorem A in Theorem 4.3.5. Along the way, we will also give a simpler proof of the following Theorem B in Corollary 3.2.3, which concerns minuscule special cycles (recalled in Section 2.10) on unitary Rapoport-Zink spaces and may be of independent interest.
Theorem B [RTZ, Theorems 9.4 and 10.1]. Let $\boldsymbol{v}=\left(v_{1}, \ldots, v_{n}\right)$ be an $n$-tuple of vectors in $C_{n}$. Assume it is minuscule in the sense that $L(\boldsymbol{v}):=\operatorname{span}_{\mathcal{O}_{E}} \boldsymbol{v}$ is an $\mathcal{O}_{E}$-lattice in $C_{n}$ satisfying $p L(\boldsymbol{v})^{\vee} \subseteq L(\boldsymbol{v}) \subseteq L(\boldsymbol{v})^{\vee}$. Let $\mathcal{Z}(\boldsymbol{v}) \subseteq \mathcal{N}$ be the associated special cycle. Then $\mathcal{Z}(\boldsymbol{v})$ is a reduced $k$-scheme.
1.3. Novelty of the proof. The original proofs of Theorems A and B form the technical heart of [RTZ] and occupy its two sections $\S 10-\S 11$. As explained below, our new proofs presented here have the merit of being much shorter and more conceptual.
1.3.1. Theorem A. The original proof of Theorem A uses Zink's theory of windows to compute the local equations of (1.2.0.1). It requires explicitly writing down the window of the universal deformation of $p$-divisible groups and solving quite involved linear algebra problems. Theorem B ensures that the intersection is entirely concentrated in the special fiber so that each local ring has the form $k[X] / X^{\ell}$. The assumption $p>c$ ensures $\ell<p$ so that the ideal of local equations is admissible (see the last paragraph of [RTZ, p. 1661]), which is crucial in order to construct the frames for the relevant windows needed in Zink's theory.

Our new proof of Theorem A does not use Zink's theory and involves little explicit computation. Our key observation is that Theorem B indeed allows us to identify the intersection (1.2.0.1) as the fixed point scheme $\mathcal{V}(\Lambda)^{\bar{g}}$ of a finite order automorphism $\bar{g}$ on a generalized Deligne-Lusztig variety $\mathcal{V}(\Lambda)$ (Section 4.1), which becomes purely an algebraic geometry problem over the residue field $k$. When $p>c$, it further simplifies to a more elementary problem of determining the fixed point scheme of a finite order automorphism $\bar{g} \in \mathrm{GL}_{d+1}(k)$ on a projective space $\mathbb{P}^{d}$ over $k$ (Section 4.2). This elementary problem has an answer in terms of the sizes of the Jordan blocks of $\bar{g}$ (Lemma 4.3.4), which explains conceptually why the intersection multiplicity should be equal to $c$. Notice that our method completely avoids computation within Zink's theory, and it would be interesting to explore the possibility of removing the assumption $p>c$ using this method.
1.3.2. Theorem $B$. The original proof of Theorem $B$ relies on showing two things (by [RTZ, Lemma 10.2]): (1) the minuscule special cycle $\mathcal{Z}(\boldsymbol{v})$ has no $W / p^{2}$ points and (2) its special fiber $\mathcal{Z}(\boldsymbol{v})_{k}$ is regular. Step (1) is relatively easy using Grothendieck-Messing theory. Step (2) is more difficult: for super-general points $x$ on $\mathcal{Z}(\boldsymbol{v})_{k}$, the regularity is shown by explicitly computing the local equation of $\mathcal{Z}(\boldsymbol{v})_{k}$ at $x$ using Zink's theory; for points which are not super-general, the regularity is shown using induction and reduces to the regularity of certain special divisors, whose local equations can again be explicitly computed using Zink's theory.

Our new proof of Theorem B does not use Zink's theory either and involves little explicit computation. Our key observation is that to show both (1) and (2) it suffices to consider the thickenings $\mathcal{O}$ of $k$ which are objects of the crystalline site of $k$. These $\mathcal{O}$-points of $\mathcal{Z}(\boldsymbol{v})$ then can be understood using only GrothendieckMessing theory (Theorem 3.1.3). We prove a slight generalization of (1) which applies to possibly nonminuscule special cycles (Corollary 3.2.1). We then prove the tangent space of the minuscule special cycle $\mathcal{Z}(\boldsymbol{v})_{k}$ has the expected dimension (Corollary 3.2.2). The desired regularity (2) follows immediately.
1.3.3. Our new proofs are largely inspired by our previous work on arithmetic intersections on GSpin Rapoport-Zink spaces [Li and Zhu 2017]. The GSpin Rapoport-Zink spaces considered there are not of PEL type, which makes them technically more complicated. So the unitary case treated here can serve as a guide to [ Li and Zhu 2017]. We have tried to indicate similarities between certain statements and proofs, for both clarity and the convenience of the readers.
1.4. Structure of the paper. In Section 2, we recall necessary backgrounds on unitary Rapoport-Zink spaces and the formulation of the arithmetic intersection problem. In Section 3, we study the local structure of the minuscule special cycles and prove Theorem B. In Section 4, we provide an alternative moduli interpretation of the generalized Deligne-Lusztig variety $\mathcal{V}(\Lambda)$ and prove Theorem A.

## 2. Unitary Rapoport-Zink spaces

In this section we review the structure of unitary Rapoport-Zink spaces. We refer to [Vollaard 2010; Vollaard and Wedhorn 2011; Kudla and Rapoport 2011] for the proofs of these facts.
2.1. Unitary Rapoport-Zink spaces. Let $p$ be an odd prime. Let $F=\mathbb{Q}_{p}, k=\overline{\mathbb{F}}_{p}$, $W=W(k)$ and $K=W[1 / p]$. Let $\sigma$ be the $p$-Frobenius acting on $\overline{\mathbb{F}}_{p}$, and we also denote by $\sigma$ the canonical lift of the $p$-Frobenius to $W$ and $K$. For any $\mathbb{F}_{p}$-algebra $R$, we also denote by $\sigma$ the Frobenius $x \mapsto x^{p}$ on $R$.

Let $E=\mathbb{Q}_{p^{2}}$ be the unramified quadratic extension of $F$. Fix a $\mathbb{Q}_{p}$-algebra embedding $\phi_{0}: \mathcal{O}_{E} \hookrightarrow W$ and denote by $\phi_{1}$ the embedding $\sigma \circ \phi_{0}: \mathcal{O}_{E} \hookrightarrow W$. The embedding $\phi_{0}$ induces an embedding between the residue fields $\mathbb{F}_{p^{2}} \hookrightarrow k$, which we shall think of as the natural embedding. For any $\mathcal{O}_{E}$-module $\Lambda$ we shall write $\Lambda_{W}$ for $\Lambda \otimes_{\mathcal{O}_{E}, \phi_{0}} W$.

Let $r$ and $s$ be positive integers and let $n=r+s$. We denote by $\mathcal{N}_{r, s}$ the unitary Rapoport-Zink spaces of signature $(r, s)$, a formally smooth formal $W$-scheme, parametrizing deformations up to quasi-isogeny of height 0 of unitary $p$-divisible groups of signature ( $r, s$ ). More precisely, for a $W$-scheme $S$, a unitary $p$-divisible groups of signature $(r, s)$ over $S$ is a triple $(X, \iota, \lambda)$, where
(1) $X$ is a $p$-divisible group of dimension $n$ and height $2 n$ over $S$,
(2) $\iota: \mathcal{O}_{E} \rightarrow \operatorname{End}(X)$ is an action satisfying the signature $(r, s)$ condition, i.e., for $\alpha \in \mathcal{O}_{E}$,

$$
\operatorname{char}(\iota(\alpha): \operatorname{Lie} X)(T)=\left(T-\phi_{0}(\alpha)\right)^{r}\left(T-\phi_{1}(\sigma)\right)^{s} \in \mathcal{O}_{S}[T]
$$

(3) $\lambda: X \rightarrow X^{\vee}$ is a principal polarization such that the associated Rosati involution induces $\alpha \mapsto \sigma(\alpha)$ on $\mathcal{O}_{E}$ via $\iota$.

Over $k$, there is a unique such triple $(X, \iota, \lambda)$ such that $X$ is supersingular, up to $\mathcal{O}_{E}$-linear isogeny preserving the polarization up to scalars. Fix such a framing triple and denote it by $\left(\mathbb{X}, \mathfrak{l}_{\mathbb{X}}, \lambda_{\mathbb{X}}\right)$.

Let $\mathrm{Nilp}_{W}$ be the category of $W$-schemes on which $p$ is locally nilpotent. Then the unitary Rapoport-Zink space $\mathcal{N}_{r, s}$ represents the functor Nilp ${ }_{W} \rightarrow$ Sets which sends $S \in \operatorname{Nilp}_{W}$ to the set of isomorphism classes of quadruples $(X, \iota, \lambda, \rho)$, where $(X, \iota, \lambda)$ is a unitary $p$-divisible group over $S$ of signature $(r, s)$ and $\rho: X \times{ }_{S} S_{k} \rightarrow$ $\mathbb{X} \times_{k} S_{k}$ is an $\mathcal{O}_{E}$-linear quasi-isogeny of height zero which respects $\lambda$ and $\lambda_{\mathbb{X}}$ up to a scalar $c(\rho) \in \mathcal{O}_{F}^{\times}=\mathbb{Z}_{p}^{\times}$(i.e., $\rho^{\vee} \circ \lambda_{X} \circ \rho=c(\rho) \cdot \lambda$ ).

In the following we denote $\mathcal{N}:=\mathcal{N}_{1, n-1}, \mathcal{M}:=\mathcal{N}_{1, n-2}$ and $\overline{\mathcal{N}}_{0}:=\mathcal{N}_{0,1} \cong \operatorname{Spf} W$. They have relative dimension $n-1, n$ and 0 over $\operatorname{Spf} W$ respectively. We denote by $\overline{\mathbb{Y}}=\left(\overline{\mathbb{Y}}, \iota_{\mathbb{Y}}, \lambda_{\bar{Y}}\right)$ the framing object for $\overline{\mathcal{N}}_{0}$ and denote by $\bar{Y}=\left(\bar{Y}, \iota_{\bar{Y}}, \lambda_{\bar{Y}}\right)$ the
universal $p$-divisible group over $\overline{\mathcal{N}}_{0}$. We may and shall choose framing objects $\mathbb{X}=\left(\mathbb{X}, \mathscr{X}_{\mathbb{X}}, \lambda_{\mathbb{X}}\right)$ and $\mathbb{X}^{b}=\left(\mathbb{X}^{b}, \iota_{\mathbb{X}^{b}}, \lambda_{\mathbb{X}^{b}}\right)$ for $\mathcal{N}$ and $\mathcal{M}$ respectively such that

$$
\mathbb{X}=\mathbb{X}^{b} \times \overline{\mathbb{Y}}
$$

as unitary $p$-divisible groups.
2.2. The group $J$. The covariant Dieudonné module $M=\mathbb{D}(\mathbb{X})$ of the framing unitary $p$-divisible group is a free $W$-module of rank $2 n$ together with an $\mathcal{O}_{E}$-action (induced by $\iota$ ) and a perfect symplectic $W$-bilinear form $\langle\cdot, \cdot\rangle: M \times M \rightarrow W$ (induced by $\lambda$ ), see [Vollaard and Wedhorn 2011, §2.3]. Let $N=M \otimes_{W} K$ be the associated isocrystal and extend $\langle\cdot, \cdot\rangle$ to $N$ bilinearly. Let $F$ and $V$ be the usual operators on $N$. We have

$$
\begin{equation*}
\langle F x, F y\rangle=p \sigma(\langle x, y\rangle), \quad \forall x, y \in N \tag{2.2.0.1}
\end{equation*}
$$

The $E$-action decomposes $N$ into a direct sum of two $K$-vector spaces of dimension $n$,

$$
\begin{equation*}
N=N_{0} \oplus N_{1} \tag{2.2.0.2}
\end{equation*}
$$

where the action of $E$ on $N_{i}$ is induced by the embedding $\phi_{i}$. Both $N_{0}$ and $N_{1}$ are totally isotropic under the symplectic form. The operator $F$ is of degree one and induces a $\sigma$-linear bijection $N_{0} \xrightarrow{\sim} N_{1}$. Since the isocrystal $N$ is supersingular, the degree 0 and $\sigma^{2}$-linear operator

$$
\Phi=V^{-1} F=p^{-1} F^{2}
$$

has all slopes zero [Kudla and Rapoport 2011, §2.1]. We have a $K$-vector space $N_{0}$ together with a $\sigma^{2}$-linear automorphism $\Phi .{ }^{1}$ The space of fixed points

$$
C=N_{0}^{\Phi}
$$

is an $E$-vector space of dimension $n$ and $N_{0}=C \otimes_{E, \phi_{0}} K$. Fix $\delta \in \mathcal{O}_{E}^{\times}$such that $\sigma(\delta)=-\delta$. Define a nondegenerate $\sigma$-sesquilinear form on $N_{0}$ by

$$
\begin{equation*}
\{x, y\}:=(p \delta)^{-1}\langle x, F y\rangle . \tag{2.2.0.3}
\end{equation*}
$$

Using (2.2.0.1) it is easy to see that

$$
\begin{equation*}
\sigma(\{x, y\})=\{\Phi y, x\}, \quad \forall x, y \in N_{0} \tag{2.2.0.4}
\end{equation*}
$$

In particular, when restricted to $C$, the form $\{\cdot, \cdot\}$ is $\sigma$-Hermitian, namely

$$
\begin{equation*}
\sigma(\{x, y\})=\{y, x\}, \quad \forall x, y \in C \tag{2.2.0.5}
\end{equation*}
$$

[^26]In fact, $(C,\{\cdot, \cdot\})$ is the unique (up to isomorphism) nondegenerate nonsplit $\sigma-$ Hermitian $E$-space of dimension $n$. Let $J=\mathrm{U}(C)$ be the unitary group of $(C,\{\cdot, \cdot\})$. It is an algebraic group over $F=\mathbb{Q}_{p}$. By Dieudonné theory, the group $J\left(\mathbb{Q}_{p}\right)$ can be identified with the automorphism group of the framing unitary $p$-divisible group $\left(\mathbb{X}, \mathscr{X}_{\mathbb{X}}, \lambda_{\mathbb{X}}\right)$ and hence acts on the Rapoport-Zink space $\mathcal{N}$.
2.3. Special homomorphisms. By definition, the space of special homomorphisms is the $\mathcal{O}_{E}$-module $\operatorname{Hom}_{\mathcal{O}_{E}}(\overline{\mathbb{Y}}, \mathbb{X})$. There is a natural $\mathcal{O}_{E}$-valued $\sigma$-Hermitian form on $\operatorname{Hom}_{\mathcal{O}_{E}}(\mathbb{Y}, \mathbb{X})$ given by

$$
(x, y) \mapsto \lambda_{\overline{\mathbb{Y}}}^{-1} \circ \hat{y} \circ \lambda_{\mathbb{X}} \circ x \in \operatorname{End}_{\mathcal{O}_{E}}(\overline{\mathbb{Y}}) \xrightarrow{\sim} \mathcal{O}_{E}
$$

By [Kudla and Rapoport 2011, Lemma 3.9], there is an isomorphism of $\sigma$-Hermitian $E$-spaces

$$
\begin{equation*}
\operatorname{Hom}_{\mathcal{O}_{E}}(\overline{\mathbb{Y}}, \mathbb{X}) \otimes_{\mathcal{O}_{E}} E \xrightarrow{\sim} C \tag{2.3.0.1}
\end{equation*}
$$

Therefore we may view elements of $C$ as special quasi-homomorphisms.
2.4. Vertex lattices. For any $\mathcal{O}_{E}$-lattice $\Lambda \subset C$, we define the dual lattice $\Lambda^{\vee}:=$ $\left\{x \in C:\{x, \Lambda\} \subseteq \mathcal{O}_{E}\right\}$. It follows from the $\sigma$-Hermitian property (2.2.0.5) that we have $\left(\Lambda^{\vee}\right)^{\vee}=\Lambda$.

A vertex lattice is an $\mathcal{O}_{E}$-lattice $\Lambda \subseteq C$ such that $p \Lambda \subseteq \Lambda^{\vee} \subseteq \Lambda$. Such lattices correspond to the vertices of the Bruhat-Tits building of the unitary group $U(C)$. Fix a vertex lattice $\Lambda$. The type of $\Lambda$ is defined to be $t_{\Lambda}:=\operatorname{dim}_{\mathbb{F}_{p^{2}}} \Lambda / \Lambda^{\vee}$, which is always an odd integer such that $1 \leq t_{\Lambda} \leq n$ (see [Vollaard 2010, Remark 2.3]).

We define $\Omega_{0}(\Lambda):=\Lambda / \Lambda^{\vee}$ and equip it with the perfect $\sigma$-Hermitian form

$$
(\cdot, \cdot): \Omega_{0}(\Lambda) \times \Omega_{0}(\Lambda) \rightarrow \mathbb{F}_{p^{2}}, \quad(x, y):=p\{\tilde{x}, \tilde{y}\} \bmod p
$$

where $\{\cdot, \cdot\}$ is the Hermitian form on $C$ defined in (2.2.0.3), and $\tilde{x}, \tilde{y} \in \Lambda$ are lifts of $x$ and $y$.

We define

$$
\Omega(\Lambda):=\Omega_{0}(\Lambda) \otimes_{F_{p^{2}}} k
$$

Remark 2.4.1. Our $\Omega_{0}(\Lambda)$ is the space $V$ in [Vollaard 2010, (2.11)], and our pairing $(\cdot, \cdot)$ differs from the pairing $(\cdot, \cdot)$ defined in [loc. cit.] by a factor of the reduction $\bar{\delta} \in \mathbb{F}_{p^{2}}^{\times}$of $\delta$.
2.5. The variety $\mathcal{V}(\boldsymbol{\Lambda})$. Let $\Lambda$ be a vertex lattice and let $\Omega_{0}=\Omega_{0}(\Lambda)$. Recall from Section 2.4 that $\Omega_{0}$ is an $\mathbb{F}_{p^{2}}$-vector space whose dimension is equal to the type $t=t_{\Lambda}$ of $\Lambda$, an odd number. Let $d:=(t-1) / 2$. We define $\mathcal{V}\left(\Omega_{0}\right)$ to be the closed $\mathbb{F}_{p^{2}}$-subscheme of the Grassmannian $\operatorname{Gr}_{d+1}\left(\Omega_{0}\right)$ (viewed as a scheme over $\mathbb{F}_{p^{2}}$ ) such
that for any $\mathbb{F}_{p^{2}}$-algebra $R$,
$\mathcal{V}\left(\Omega_{0}\right)(R)=\left\{R\right.$-module local direct summands $U \subseteq \Omega_{0}{\otimes \mathbb{F}_{p^{2}} R}$ :

$$
\begin{equation*}
\left.\operatorname{rank} U=d+1 \text { and } U^{\perp} \subseteq U\right\} \tag{2.5.0.1}
\end{equation*}
$$

Here $U^{\perp}$ is by definition $\left\{v \in \Omega_{0} \otimes R:(v, u)_{R}=0, \forall u \in U\right\}$, where $(\cdot, \cdot)_{R}$ is the $R$-sesquilinear form on $\Omega_{0} \otimes R$ obtained from $(\cdot, \cdot)$ by extension of scalars (linearly in the first variable and $\sigma$-linearly in the second variable). Then $\mathcal{V}\left(\Omega_{0}\right)$ is a smooth projective $\mathbb{F}_{p^{2}}$-scheme of dimension $d$ by [Vollaard 2010, Proposition 2.13] and Remark 2.4.1. In fact, $\mathcal{V}\left(\Omega_{0}\right)$ can be identified as a (generalized) DeligneLusztig variety, by [Vollaard and Wedhorn 2011, §4.5] (though we will not use this identification in the following).

We write $\mathcal{V}(\Lambda)$ for the base change of $\mathcal{V}\left(\Omega_{0}\right)$ from $\mathbb{F}_{p^{2}}$ to $k$.
2.6. Structure of the reduced scheme $\boldsymbol{\mathcal { N }}^{\text {red }}$. For each vertex lattice $\Lambda \subseteq C$, we define $\mathcal{N}_{\Lambda} \subseteq \mathcal{N}$ to be the locus where $\rho_{X}^{-1} \circ \Lambda^{\vee} \subseteq \operatorname{Hom}(\bar{Y}, X)$, i.e., where the quasi-homomorphisms $\rho^{-1} \circ v$ lift to actual homomorphisms for any $v \in \Lambda^{\vee}$. Then $\mathcal{N}_{\Lambda}$ is a closed formal subscheme by [Rapoport and Zink 1996, Proposition 2.9]. By [Vollaard and Wedhorn 2011, §4] we have an isomorphism of $k$-varieties

$$
\begin{equation*}
\mathcal{N}_{\Lambda}^{\text {red }} \xrightarrow{\sim} \mathcal{V}(\Lambda) . \tag{2.6.0.1}
\end{equation*}
$$

2.7. Some invariants associated to a $\boldsymbol{k}$-point of $\boldsymbol{\mathcal { N }}$. We follow [Kudla and Rapoport 2011, §2.1].

Let $x$ be a point in $\mathcal{N}(k)$. Then $x$ represents a tuple $(X, \iota, \lambda, \rho)$ over $k$ as recalled in Section 2.1. Via $\rho$, we view the Dieudonné module of $X$ as a $W$-lattice $M_{x}$ in $N$, which is stable under the operators $F$ and $V$. The endomorphism structure $\iota$ induces an action of $\mathcal{O}_{E} \otimes_{\mathbb{Z}_{p}} W \cong W \oplus W$ on $M_{x}$, which is equivalent to the structure of a $\mathbb{Z} / 2 \mathbb{Z}$-grading on $M_{x}$ (into $W$-modules). We denote this grading by

$$
M_{x}=\operatorname{gr}_{0} M_{x} \oplus \operatorname{gr}_{1} M_{x}
$$

This grading is compatible with (2.2.0.2) in the sense that

$$
\operatorname{gr}_{i} M_{x}=M_{x} \cap N_{i}, \quad i=0,1
$$

Moreover both $\mathrm{gr}_{0} M_{x}$ and $\mathrm{gr}_{1} M_{x}$ are free $W$-modules of rank $n$.
Consider the $k$-vector space $M_{x, k}:=M_{x} \otimes_{W} k$. It has an induced $\mathbb{Z} / 2 \mathbb{Z}$-grading, as well as a canonical filtration $\operatorname{Fil}^{1}\left(M_{x, k}\right) \subset M_{x, k}$. Explicitly, $\operatorname{Fil}^{1}\left(M_{x, k}\right)$ is the image of $V\left(M_{x}\right) \subseteq M_{x}$ under the reduction map $M_{x} \rightarrow M_{x, k}$. Define

$$
\operatorname{Fil}^{1}\left(\operatorname{gr}_{i} M_{x, k}\right):=\operatorname{Fil}^{1}\left(M_{x, k}\right) \cap \operatorname{gr}_{i} M_{x, k} .
$$

Then

$$
\operatorname{Fil}^{1}\left(M_{x, k}\right)=\operatorname{Fil}^{1}\left(\operatorname{gr}_{0} M_{x, k}\right) \oplus \operatorname{Fil}^{1}\left(\operatorname{gr}_{1} M_{x, k}\right),
$$

and by the signature $(1, n-1)$ condition we know that $\mathrm{Fil}^{1}\left(\mathrm{gr}_{0} M_{x, k}\right)$ is a hyperplane and $\mathrm{Fil}^{1}\left(\mathrm{gr}_{1} M_{x, k}\right)$ is a line in $\mathrm{gr}_{0} M_{x, k}$ and $\mathrm{gr}_{1} M_{x, k}$, respectively.

The symplectic form $\langle\cdot, \cdot\rangle$ on $N$ takes values in $W$ on $M_{x}$, and hence induces a symplectic form on $M_{x, k}$ by reduction. The latter restricts to a $k$-bilinear nondegenerate pairing

$$
\operatorname{gr}_{0} M_{x, k} \times \operatorname{gr}_{1} M_{x, k} \rightarrow k
$$

Under the above pairing, the spaces $\mathrm{Fil}^{1}\left(\mathrm{gr}_{0} M_{x, k}\right)$ and $\mathrm{Fil}^{1}\left(\mathrm{gr}_{1} M_{x, k}\right)$ are annihilators of each other. Equivalently, $\operatorname{Fil}^{1}\left(M_{x, k}\right)$ is a totally isotropic subspace of $M_{x, k}$.
2.8. Description of $\boldsymbol{k}$-points by special lattices. For a $W$-lattice $A$ in $N_{0}$, we define its dual lattice to be $A^{\vee}:=\left\{x \in N_{0}:\{x, A\} \subseteq W\right\}$. If $\Lambda$ is an $\mathcal{O}_{E}$-lattice in $C$, then we have $\left(\Lambda_{W}\right)^{\vee}=\left(\Lambda^{\vee}\right)_{W}$. In the following we denote both of them by $\Lambda_{W}^{\vee}$.
Definition 2.8.1. A special lattice is a $W$-lattice $A$ in $N_{0}$ such that

$$
A^{\vee} \subseteq A \subseteq p^{-1} A^{\vee}
$$

and such that $A / A^{\vee}$ is a one-dimensional $k$-vector space.
Remark 2.8.2. The apparent difference between the above definition and the condition in [Vollaard 2010, Proposition 1.10] (for $i=0$ ) is caused by the fact that we have normalized the pairing $\{\cdot, \cdot\}$ on $N_{0}$ differently from [loc. cit.], using an extra factor $(p \delta)^{-1}$ (see (2.2.0.3)). Our normalization is the same as that in [RTZ].

Recall the following result.
Proposition 2.8.3 [Vollaard 2010, Proposition 1.10]. There is a bijection from $\mathcal{N}(k)$ to the set of special lattices, sending a point $x$ to $\mathrm{gr}_{0} M_{x}$ considered in Section 2.7.

Remark 2.8.4. Let $x \in \mathcal{N}(k)$ and let $A$ be the special lattice associated to it by Proposition 2.8.3. Let $\Lambda$ be a vertex lattice. Then $x \in \mathcal{N}_{\Lambda}(k)$ if and only if $A \subseteq \Lambda_{W}$, if and only if $\Lambda_{W}^{\vee} \subseteq A^{\vee}$. (See also Remark 3.1.5 below.)
2.9. Filtrations. We introduce the following notation:

Definition 2.9.1. Let $A$ be a special lattice. Write $A_{k}:=A \otimes_{W} k$. Let $x \in \mathcal{N}(k)$ correspond to $A$ under Proposition 2.8.3. Thus $A_{k}=\operatorname{gr}_{0} M_{x, k}$. Define Fil ${ }^{1}\left(A_{k}\right):=$ $\mathrm{Fil}^{1}\left(\mathrm{gr}_{0} M_{x, k}\right)$ (see Section 2.7). It is a hyperplane in $A_{k}$.
Lemma 2.9.2. Let $A$ be a special lattice. Then $\Phi^{-1}\left(A^{\vee}\right)$ is contained inside $A$, and its image in $A_{k}$ is equal to $\operatorname{Fil}^{1}\left(A_{k}\right)$.

Proof. Let $A$ correspond to $x \in \mathcal{N}(k)$ under Proposition 2.8.3. Then $F$ and $V$ both preserve the $W$-lattice $M_{x}$ in $N$ (see Section 2.7). By definition, $\operatorname{Fil}^{1}\left(M_{x, k}\right)$ is the image of $V\left(M_{x}\right) \subseteq M_{x}$ under the reduction map $M_{x} \rightarrow M_{x, k}$. Since the operator $V$
is of degree 1 with respect to the $\mathbb{Z} / 2 \mathbb{Z}$-grading, we see that $\mathrm{Fil}^{1}\left(A_{k}\right)$ is the image of $V\left(\operatorname{gr}_{1} M_{x}\right) \subseteq A$ under $A \rightarrow A_{k}$. It suffices to prove that

$$
\begin{equation*}
\Phi^{-1}\left(A^{\vee}\right)=V\left(\operatorname{gr}_{1} M_{x}\right) \tag{2.9.2.1}
\end{equation*}
$$

By the proof of [Vollaard 2010, Proposition 1.10], we have $\operatorname{gr}_{1} M_{x}=F^{-1} A^{\vee}$. (Note that because of the difference of normalizations as discussed in Remark 2.8.2, what is denoted by $A^{\vee}$ here is denoted by $p A^{\vee}$ in [loc. cit.]. Also note that the integer $i$ appearing [loc. cit.] is 0 in our case.) Therefore $V\left(\operatorname{gr}_{1} M_{x}\right)=V\left(F^{-1} A^{\vee}\right)$. But $V F^{-1}=\left(V^{-1} F\right)^{-1}=\Phi^{-1}$ because $V F=F V=p$. Thus (2.9.2.1) holds as desired.
2.10. Special cycles. Let $\boldsymbol{v}$ be an arbitrary subset of $C$. We define the special cycle $\mathcal{Z}(\boldsymbol{v}) \subseteq \mathcal{N}$ to be the locus where $\rho^{-1} \circ v \in \operatorname{Hom}(\bar{Y}, X)$ for all $v \in \boldsymbol{v}$, i.e., all the quasi-homomorphisms $\rho^{-1} \circ v$ lift to actual homomorphisms. Note that $\mathcal{Z}(\boldsymbol{v})$ only depends on the $\mathcal{O}_{E}$-submodule $L(\boldsymbol{v})$ spanned by $\boldsymbol{v}$ in $C$, and we have $\mathcal{Z}(\boldsymbol{v})=\mathcal{Z}(L(\boldsymbol{v}))$.

We say $\boldsymbol{v}$ is minuscule if $L(\boldsymbol{v})$ is an $\mathcal{O}_{E}$-lattice in $C$ satisfying $p L(\boldsymbol{v})^{\vee} \subseteq L(\boldsymbol{v}) \subseteq$ $L(\boldsymbol{v})^{\vee}$, or equivalently, if $L(\boldsymbol{v})$ is the dual of a vertex lattice. When this is the case we have $\mathcal{Z}(\boldsymbol{v})=\mathcal{N}_{L(\boldsymbol{v})^{\vee}}$ by definition.
2.11. The intersection problem. Let $C^{b}$ be the analogue for $\mathcal{M}$ of the Hermitian space $C$. Then $C \cong C^{b} \oplus E u$ for some vector denoted by $u$ which is of norm 1 and orthogonal to $C^{b}$. We have a closed immersion

$$
\delta: \mathcal{M} \rightarrow \mathcal{N}
$$

sending $(X, \iota, \lambda, \rho)$ to $\left(X \times \bar{Y}, \iota \times \iota_{\bar{Y}}, \lambda \times \lambda_{\bar{Y}}, \rho \times \mathrm{id}\right)$. We have $\delta(\mathcal{M})=\mathcal{Z}(u)$. The closed immersion $\delta$ induces a closed immersion of formal schemes

$$
(\mathrm{id}, \delta): \mathcal{M} \rightarrow \mathcal{M} \times{ }_{W} \mathcal{N}
$$

Denote by $\Delta$ the image of (id, $\delta$ ), which we call the (local) GGP cycle. For any $g \in J\left(\mathbb{Q}_{p}\right)$, we obtain a formal subscheme

$$
(\mathrm{id} \times g) \Delta \subseteq \mathcal{M} \times{ }_{W} \mathcal{N}
$$

via the action of $g$ on $\mathcal{N}$. Let $g \in J\left(\mathbb{Q}_{p}\right)$ and let $\mathcal{N}^{g} \subseteq \mathcal{N}$ be the fixed locus of $g$. Then by definition we have

$$
\Delta \cap(\operatorname{id} \times g) \Delta \cong \delta(\mathcal{M}) \cap \mathcal{N}^{g}
$$

Our goal is to compute the arithmetic intersection number

$$
\langle\Delta,(\operatorname{id} \times g) \Delta\rangle
$$

when $g$ is regular semisimple and minuscule (as defined in the introduction). Notice that $g \in J\left(\mathbb{Q}_{p}\right)$ is regular semisimple if and only if $\boldsymbol{v}(g):=\left(u, g u, \ldots, g^{n-1} u\right)$ is an $E$-basis of $C$. Also notice that a regular semisimple element $g$ is minuscule if and only if $\boldsymbol{v}(\mathrm{g})$ is minuscule in the sense of Section 2.10.

## 3. Reducedness of minuscule special cycles

### 3.1. Local structure of special cycles.

Definition 3.1.1. Let $\mathscr{C}$ be the following category:

- Objects in $\mathscr{C}$ are triples $(\mathcal{O}, \mathcal{O} \rightarrow k, \delta)$, where $\mathcal{O}$ is a local Artinian $W$-algebra, $\mathcal{O} \rightarrow k$ is a $W$-algebra map, and $\delta$ is a nilpotent divided power structure on $\operatorname{ker}(\mathcal{O} \rightarrow k)$ (see [Berthelot and Ogus 1978, Definitions 3.1, 3.27]).
- Morphisms in $\mathscr{C}$ are $W$-algebra maps that are compatible with the structure maps to $k$ and the divided power structures.
3.1.2. Let $x \in \mathcal{N}(k)$ correspond to a special lattice $A$ under Proposition 2.8.3. Let $\mathcal{O} \in \mathscr{C}$. By a hyperplane in $A_{\mathcal{O}}:=A \otimes_{W} \mathcal{O}$ we mean a free direct summand of $A_{\mathcal{O}}$ of rank $n-1$. We define the $\mathbb{Z} / 2 \mathbb{Z}$-grading on $M_{x, \mathcal{O}}:=M_{x} \otimes_{W} \mathcal{O}$ by linearly extending that on $M_{x}$ (see Section 2.7). Denote by $\widehat{\mathcal{N}}_{x}$ the completion of $\mathcal{N}$ at $x$. For any $\tilde{x} \in \widehat{\mathcal{N}}_{x}(\mathcal{O})$, we have a unitary $p$-divisible group of signature $(1, n-1)$ over $\mathcal{O}$ deforming that over $k$ defined by $x$. By Grothendieck-Messing theory, we obtain the Hodge filtration $\operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}} \subseteq M_{x, \mathcal{O}}$. Define $f_{\mathcal{O}}(\tilde{x})$ to be the intersection

$$
\operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}} \cap \operatorname{gr}_{0} M_{x, \mathcal{O}}
$$

inside $M_{x, \mathcal{O}}$. By the signature $(1, n-1)$ condition, $f_{\mathcal{O}}(\tilde{x})$ is a hyperplane in $A_{\mathcal{O}}$. It also lifts Fil ${ }^{1} A_{k}$ (see Definition 2.9.1) by construction. Thus we have defined a map

$$
\begin{equation*}
\left.f_{\mathcal{O}}: \widehat{\mathcal{N}}_{x}(\mathcal{O}) \xrightarrow{\sim} \text { \{hyperplanes in } A_{\mathcal{O}} \text { lifting } \mathrm{Fil}^{1} A_{k}\right\} . \tag{3.1.2.1}
\end{equation*}
$$

By construction, $f_{\mathcal{O}}$ is functorial in $\mathcal{O}$ in the sense that the collection $\left(f_{\mathcal{O}}\right)_{\mathcal{O} \in \mathscr{C}}$ is a natural transformation between two set-valued functors on $\mathscr{C}$. Here we are viewing the right hand side of (3.1.2.1) as a functor in $\mathcal{O}$ using the base change maps.

The following result is the analogue of [ Li and Zhu 2017, Theorem 4.1.7]. As a direct consequence of the PEL moduli problem, it should be well known to the experts and is essentially proved in [Kudla and Rapoport 2011, Proposition 3.5].

Theorem 3.1.3. Keep the notations in Section 3.1.2:
(1) The natural transformation $\left(f_{\mathcal{O}}\right)_{\mathcal{O} \in \mathscr{C}}$ is an isomorphism.
(2) Let $\boldsymbol{v}$ be a subset of $C$. If $x \in \mathcal{Z}(\boldsymbol{v})(k)$, then $\boldsymbol{v} \subseteq$ A. Suppose $x \in \mathcal{Z}(\boldsymbol{v})(k)$. Then for any $\mathcal{O} \in \mathscr{C}$ the map $f_{\mathcal{O}}$ induces a bijection

$$
\widehat{\mathcal{Z}(\boldsymbol{v})_{x}}(\mathcal{O}) \xrightarrow{\sim}
$$

$\left\{\right.$ hyperplanes in $A_{\mathcal{O}}$ lifting $\mathrm{Fil}^{1} A_{k}$ and containing the image of $v$ in $\left.A_{\mathcal{O}}\right\}$.
Proof. (1) We need to check that for all $\mathcal{O} \in \mathscr{C}$ the map $f_{\mathcal{O}}$ is a bijection. Let $\tilde{x} \in \widehat{\mathcal{N}}_{x}(\mathcal{O})$. This represents a deformation over $\mathcal{O}$ of the $p$-divisible group at $x$. Similarly to the situation in Section 2.7, the compatibility with the endomorphism structure implies that

$$
\operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}}=\bigoplus_{i=0}^{1} \operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}} \cap \operatorname{gr}_{i} M_{x, \mathcal{O}}
$$

By the compatibility with the polarization, we know that $\operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}}$ is totally isotropic under the symplectic form on $M_{x, \mathcal{O}}$. It follows that the two modules $\mathrm{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}} \cap \mathrm{gr}_{1} M_{x, \mathcal{O}}$ and $\operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}} \cap \mathrm{gr}_{0} M_{x, \mathcal{O}}$ are annihilators of each other if we identify $\operatorname{gr}_{1} M_{x, \mathcal{O}}$ as the $\mathcal{O}$-linear dual of $\operatorname{gr}_{0} M_{x, \mathcal{O}}$ using the symplectic form on $M_{x, \mathcal{O}}$. Therefore, $\operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}}$ can be recovered from $f_{\mathcal{O}}(\tilde{x})$. This together with Grothendieck-Messing theory proves the injectivity of $f_{\mathcal{O}}$. The surjectivity of $f_{\mathcal{O}}$ also follows from Grothendieck-Messing theory and the above way of reconstructing $\operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}}$ from its intersection with $\operatorname{gr}_{0} M_{x, \mathcal{O}}$. Note that the unitary $p$-divisible groups reconstructed in this way do satisfy the signature condition because we have started with hyperplanes in $A_{\mathcal{O}}$.
(2) The statements follow from the proof of [Kudla and Rapoport 2011, Proposition 3.5] and the definition of (2.3.0.1) in [Kudla and Rapoport 2011, Lemma 3.9]. We briefly recall the arguments here. If $\phi \in \operatorname{Hom}_{\mathcal{O}_{E}}(\overline{\mathbb{Y}}, \mathbb{X}) \otimes_{\mathcal{O}_{E}} E$ is a special quasi-homomorphism, the element $v \in C$ corresponding to $\phi$ under (2.3.0.1) is by definition the projection to $N_{0}$ of $\phi_{*}\left(\overline{1}_{0}\right) \in N$, where $\phi_{*}$ is the map $\mathbb{D}(\overline{\mathbb{Y}}) \otimes_{W} K \rightarrow$ $\mathbb{D}(\mathbb{X}) \otimes_{W} K=N$ induced by $\phi$, and $\overline{1}_{0}$ is a certain fixed element in $\mathbb{D}(\overline{\mathbb{Y}})$. In fact, $\overline{1}_{0}$ is chosen such that

- $W \overline{1}_{0}=\operatorname{gr}_{0} \mathbb{D}(\overline{\mathbb{Y}})$, where the grading is with respect to the $\mathcal{O}_{E}$-action on $\overline{\mathbb{Y}}$,
- $W \overline{1}_{0}=\operatorname{Fil} \frac{1}{\bar{Y}} \mathbb{D}(\overline{\mathbb{Y}})$, the Hodge filtration for the deformation $\bar{Y}$ of $\overline{\mathbb{Y}}$ over $W$.

In particular $v$ and $\phi$ are related by the formula $v=\phi_{*}\left(\overline{1}_{0}\right)$, as the projection to $N_{0}$ is not needed.

From now on we assume without loss of generality that $\boldsymbol{v}=\{v\}$, with $v$ corresponding to $\phi$ as in the above paragraph. If $x \in \mathcal{Z}(v)(k)$, then $\phi_{*}$ has to map $\mathbb{D}(\overline{\mathbb{Y}})$ into $M_{x}$, so $v \in M_{x}$. Since $\phi_{*}$ is compatible with the $\mathbb{Z} / 2 \mathbb{Z}$-gradings, we further have $v \in A$. We have shown that if $x \in \mathcal{Z}(v)(k)$, then $v \in A$.

Now suppose $x \in \mathcal{Z}(\boldsymbol{v})(k)$. Let $\mathcal{O} \in \mathscr{C}$. Write $v_{\mathcal{O}}:=v \otimes 1 \in A_{\mathcal{O}} \subset M_{x, \mathcal{O}}$. For all $\tilde{x} \in \widehat{\mathcal{N}}_{x}(\mathcal{O})$, by Grothendieck-Messing theory we know that $\tilde{x} \in \widehat{\mathcal{Z}(v)_{x}}(\mathcal{O})$ if and only if the base change of $\phi_{*}$ to $\mathcal{O}$ (still denoted by $\phi_{*}$ ) preserves the Hodge filtrations, i.e.,

$$
\phi_{*}\left(\operatorname{Fil}_{\bar{Y}}^{1} \mathbb{D}(\overline{\mathbb{Y}})\right) \subseteq \operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}}
$$

Since $W \overline{1}_{0}=\operatorname{Fil} \frac{1}{\bar{Y}} \mathbb{D}(\overline{\mathbb{Y}})$, this last condition is equivalent to $v_{\mathcal{O}} \in \operatorname{Fil}_{\tilde{x}}^{1} M_{x, \mathcal{O}}$. Again, because $\phi_{*}$ is compatible with the $\mathbb{Z} / 2 \mathbb{Z}$-gradings, the last condition is equivalent to $v_{\mathcal{O}} \in f_{\mathcal{O}}(\tilde{x})$. In conclusion, we have shown that $\tilde{x} \in \widehat{\mathcal{N}}_{x}(\mathcal{O})$ is in $\widehat{\mathcal{Z}(v)_{x}}(\mathcal{O})$ if and only if $v_{\mathcal{O}} \in f_{\mathcal{O}}(\tilde{x})$, as desired.

Corollary 3.1.4. Let $x \in \mathcal{N}(k)$ correspond to the special lattice $A$. Let $\boldsymbol{v}$ be a subset of $C$. Then $x \in \mathcal{Z}(\boldsymbol{v})(k)$ if and only if $\boldsymbol{v} \subseteq A^{\vee}$.

Proof. By part (2) of Theorem 3.1.3 applied to $\mathcal{O}=k$, we see that $x \in \mathcal{Z}(\boldsymbol{v})(k)$ if and only if $\boldsymbol{v} \subseteq A$ and the image of $\boldsymbol{v}$ in $A_{k}$ is contained in $\operatorname{Fil}^{1}\left(A_{k}\right)$. The corollary follows from Lemma 2.9.2 and the $\Phi$-invariance of $\boldsymbol{v}$.

Remark 3.1.5. Note that Remark 2.8.4 is a special case of Corollary 3.1.4.

### 3.2. Proof of the reducedness.

Corollary 3.2.1. Let $\Lambda$ be an $\mathcal{O}_{E}$-lattice in $C$ with $p^{i} \Lambda \subseteq \Lambda^{\vee} \subseteq \Lambda$ for some $i \in \mathbb{Z}_{\geq 1}$. Then the special cycle $\mathcal{Z}\left(\Lambda^{\vee}\right)$ defined by $\Lambda^{\vee}$ has no $\left(W / p^{i+1}\right)$-points. In particular, taking $i=1$ we see that $\mathcal{N}_{\Lambda}\left(W / p^{2}\right)=\varnothing$ for any vertex lattice $\Lambda$.

Proof. Let $\mathcal{O}=W / p^{i+1}$, equipped with the reduction map $W / p^{i+1} \rightarrow k$ and the natural divided power structure on the kernel $p \mathcal{O}$. Then $\mathcal{O} \in \mathscr{C}$. Assume $\mathcal{Z}\left(\Lambda^{\vee}\right)$ has an $\mathcal{O}$-point $\tilde{x}$ reducing to a $k$-point $x$. Let $A$ be the special lattice corresponding to $x$ (see Section 2.8). By Theorem 3.1.3, there exists a hyperplane $P$ in $A_{\mathcal{O}}$ lifting $\operatorname{Fil}^{1}\left(A_{k}\right)$, such that $P \supseteq \Lambda^{\vee} \otimes_{\mathcal{O}_{E}} \mathcal{O}$. Since $P$ is a hyperplane in $A_{\mathcal{O}}$, there exists an element $l \in \operatorname{Hom}_{\mathcal{O}}\left(A_{\mathcal{O}}, \mathcal{O}\right)$ such that

$$
\begin{equation*}
l(P)=0 \quad \text { and } \quad l\left(A_{\mathcal{O}}\right)=\mathcal{O} \tag{3.2.1.1}
\end{equation*}
$$

We may find an element $\tilde{l} \in A^{\vee} \subseteq N_{0}$ to represent $l$, in the sense that for all $a \otimes 1 \in A_{\mathcal{O}}$ with $a \in A$, we have

$$
l(a \otimes 1)=\text { the image of }\{a, \tilde{l}\} \text { under } W \rightarrow \mathcal{O}
$$

Since $l\left(\Lambda^{\vee} \otimes \mathcal{O}\right) \subseteq l(P)=0$, we know that $\{v, \tilde{l}\} \in p^{i+1} W$ for all $v \in \Lambda^{\vee}$. Since $\Lambda^{\vee} \subseteq C=N_{0}^{\Phi}$, applying (2.2.0.4) we see that $\{\tilde{l}, v\} \in p^{i+1} W$ for all $v \in \Lambda^{\vee}$. Therefore

$$
p^{-i-1} \tilde{l} \in\left(\Lambda_{W}^{\vee}\right)^{\vee}=\Lambda_{W}
$$

and thus $\tilde{l} \in p^{i+1} \Lambda_{W}$, which is contained in $p \Lambda_{W}^{\vee}$ by hypothesis. Since $\Lambda$ is $\Phi$-invariant, we also have $\Phi(\tilde{l}) \in p \Lambda_{W}^{\vee}$. But $\Lambda_{W}^{\vee} \subseteq A^{\vee}$ by Corollary 3.1.4, so $\Phi(\tilde{l}) \in p A^{\vee}$. It follows that for all $a \in A$, we have $\{\Phi(\tilde{l}), a\} \in p W$, and therefore

$$
\{a, \tilde{l}\} \stackrel{(2.2 .0 .4)}{=} \sigma^{-1}(\{\Phi(\tilde{l}), a\}) \in p W
$$

contradicting the second condition in (3.2.1.1).
Corollary 3.2.2. Let $\Lambda$ be a vertex lattice of type $t$ and let $x \in \mathcal{N}_{\Lambda}(k)$. Then the tangent space $\mathcal{T}_{x} \mathcal{N}_{\Lambda, k}$ to $\mathcal{N}_{\Lambda, k}$ at $x$, where $\mathcal{N}_{\Lambda, k}$ is the special fiber (i.e., base change to $k$ ) of $\mathcal{N}_{\Lambda}$, is of $k$-dimension $(t-1) / 2$.
Proof. This can be deduced from Theorem 3.1.3 elementarily, in the same way as in [Li and Zhu 2017, §4.2]. Here we provide a shorter proof. Firstly we make an easy observation. Denote by $\mathscr{C}_{k}$ the full subcategory of $\mathscr{C}$ consisting of characteristic $p$ objects. Let $\mathcal{W}_{1}$ and $\mathcal{W}_{2}$ be two formal schemes over $k$. For $i=1,2$ fix $y_{i} \in \mathcal{W}_{i}(k)$ and define the set-valued functor $\mathcal{F}_{i}$ on $\mathscr{C}_{k}$ sending $\mathcal{O}$ to the set of $\mathcal{O}$-points of $\mathcal{W}_{i}$ which induce $y_{i}$ under the structure map $\mathcal{O} \rightarrow k$. Assume $\mathcal{F}_{1} \cong \mathcal{F}_{2}$. Then the tangent spaces $\mathcal{T}_{x_{i}} \mathcal{W}_{i}$ are isomorphic. In fact, this observation is a direct consequence of the definition of the vector space structure on the tangent spaces from the point of view of functor of points, as recalled in the proof of [ Li and Zhu 2017, Lemma 4.2.6] for instance.

Denote by $B$ the $k$-subspace of $A_{k}$ spanned by the image of $\Lambda^{\vee}$ in $A_{k}$. Consider the Grassmannian $\operatorname{Gr}_{n-1}\left(A_{k}\right)$ parametrizing hyperplanes in the $n$-dimensional $k$ vector space $A_{k}$. Let $\mathcal{W}_{1}$ be the subvariety of $\mathrm{Gr}_{n-1}\left(A_{k}\right)$ defined by the condition that the hyperplane should contain $B$, and let $y_{1} \in \mathcal{W}_{1}(k)$ corresponding to $\operatorname{Fil}^{1}\left(A_{k}\right) \subseteq A_{k}$. Let $\mathcal{W}_{2}:=\mathcal{N}_{\Lambda, k}$ and $y_{2}:=x$. By Theorem 3.1.3, the assumption on $\left(\mathcal{W}_{i}, y_{i}\right), i=1,2$ in the previous paragraph is satisfied. Hence it suffices to compute the dimension of $\mathcal{T}_{y_{1}} \mathcal{W}_{1}$. Note that $\mathcal{W}_{1}$ is itself a Grassmannian, parametrizing hyperplanes in $A_{k} / B$. The proof is finished once we know that $A_{k} / B$ has $k$-dimension $(t+1) / 2$. But this is true by the ( $\sigma$-linear) duality between the $k$-vector spaces $A_{k} / B=A / \Lambda_{W}^{\vee}$ and $\Lambda_{W} / A^{\vee}$ under the $\sigma$-sesquilinear form on $\Omega(\Lambda)$ obtained by extension of scalars from the $\sigma$-Hermitian form $(\cdot, \cdot)$ on $\Omega_{0}(\Lambda)$ (see Section 2.4) and the fact that $A / A^{\vee}$ is a 1 -dimensional $k$-vector space (see Definition 2.8.1).

In the following corollary we reprove [RTZ, Theorems 9.4 and 10.1].
Corollary 3.2.3. Let $\Lambda$ be a vertex lattice. Then $\mathcal{N}_{\Lambda}=\mathcal{N}_{\Lambda, k}=\mathcal{N}_{\Lambda}^{\text {red }}$ and it is regular.
Proof. Let $t$ be the type of $\Lambda$. Recall from Section 2.6 that $\mathcal{N}_{\Lambda}^{\text {red }}$ is a smooth $k$-scheme of dimension $(t-1) / 2$. By Corollary 3.2.2, all the tangent spaces of $\mathcal{N}_{\Lambda, k}$ have $k$-dimension $(t-1) / 2$, and so $\mathcal{N}_{\Lambda, k}$ is regular. In particular $\mathcal{N}_{\Lambda, k}$ is reduced, namely $\mathcal{N}_{\Lambda, k}=\mathcal{N}_{\Lambda}^{\text {red }}$. Knowing that $\mathcal{N}_{\Lambda, k}$ is regular, and that $\mathcal{N}_{\Lambda}$ has no ( $W / p^{2}$ )-points (Corollary 3.2.1), it follows that $\mathcal{N}_{\Lambda}=\mathcal{N}_{\Lambda, k}$ by the general criterion [RTZ, Lemma 10.3].

## 4. The intersection length formula

4.1. The arithmetic intersection as a fixed point scheme. Fix a regular semisimple and minuscule element $g \in J\left(\mathbb{Q}_{p}\right)$. Let $L:=L(\boldsymbol{v}(g))$ and $\Lambda:=L^{\vee}$. They are both $\mathcal{O}_{E}$-lattices in $C$. Recall from the end of Section 2.10 that $\Lambda$ is a vertex lattice and $\mathcal{Z}(L)=\mathcal{N}_{\Lambda}$. From now on we assume $\mathcal{N}^{g}(k) \neq \varnothing$. As shown in [RTZ, §5], this assumption implies that both $L$ and $\Lambda$ are $g$-cyclic and stable under $g$. In particular, the natural action of $g$ on $\mathcal{N}$ stabilizes $\mathcal{N}_{\Lambda}$.

Let $\Omega_{0}=\Omega_{0}(\Lambda)$ and $\Omega=\Omega(\Lambda)$. Let $t=t_{\Lambda}$ and $d=(t-1) / 2$ as in Section 2.5. Let $\bar{g} \in \operatorname{GL}\left(\Omega_{0}\right)\left(\mathbb{F}_{p^{2}}\right)$ be the induced action of $g$ on $\Omega_{0}$. Then $\bar{g}$ preserves the Hermitian form $(\cdot, \cdot)$ on $\Omega_{0}$ and hence acts on $\mathcal{V}(\Lambda)$. It is clear from the definition of the isomorphism (2.6.0.1) given in [Vollaard and Wedhorn 2011, §4] that it is equivariant for the actions of $g$ and $\bar{g}$ on the two sides.
Remark 4.1.1. Since $\Lambda$ and $\Lambda^{\vee}$ are $g$-cyclic, the linear operator $\bar{g} \in \operatorname{GL}\left(\Omega_{0}\right)\left(\mathbb{F}_{p^{2}}\right)$ has equal minimal polynomial and characteristic polynomial. Equivalently, in the Jordan normal form of $\bar{g}$ (over $k$ ) there is a unique Jordan block associated to any eigenvalue.

Proposition 4.1.2. $\delta(\mathcal{M}) \cap \mathcal{N}^{g}$ is a scheme of characteristic $p$ (i.e., a $k$-scheme) isomorphic to $\mathcal{V}(\Lambda)^{\bar{g}}$.
Proof. Recall from Section 2.11 that $\delta(\mathcal{M})=\mathcal{Z}(u)$. Since the $\mathcal{O}_{E}$-module $L$ is generated by $u, g u, \cdots, g^{n-1} u$ and stable under $g$, we have $\delta(\mathcal{M}) \cap \mathcal{N}^{g}=$ $\mathcal{Z}(L)^{g}=\mathcal{N}_{\Lambda}^{g}$. By Corollary 3.2.3, we know that $\mathcal{N}_{\Lambda}^{g}=\left(\mathcal{N}_{\Lambda}^{\text {red }}\right)^{g}$. But the latter is isomorphic to the characteristic $p$ scheme $\mathcal{V}(\Lambda)^{\bar{g}}$ under (2.6.0.1).
4.2. Study of $\mathcal{V}(\boldsymbol{\Lambda})^{\bar{g}}$. We start with an alternative moduli interpretation of $\mathcal{V}\left(\Omega_{0}\right)$. The idea is to rewrite (in Lemma 4.2.2) the procedure of taking the complement $U \mapsto U^{\perp}$ with respect to the Hermitian form, in terms of taking the complement with respect to some quadratic form and taking a Frobenius. The alternative moduli interpretation is given in Corollary 4.2.3 below.

Let $\Theta_{0}$ be a $t$-dimensional nondegenerate quadratic space over $\mathbb{F}_{p}$. Let $[\cdot, \cdot]$ : $\Theta_{0} \times \Theta_{0} \rightarrow \mathbb{F}_{p}$ be the associated bilinear form. Since there is a unique isomorphism class of nondegenerate $\sigma$-Hermitian spaces over $\mathbb{F}_{p^{2}}$, we may assume that $\Omega_{0}=$ $\Theta_{0} \otimes_{\mathbb{F}_{p}} \mathbb{F}_{p^{2}}$ and that the $\sigma$-Hermitian form $(\cdot, \cdot)$ (see Section 2.5) is obtained by extension of scalars (linearly in the first variable and $\sigma$-linearly in the second variable) from $[\cdot, \cdot]$.
Definition 4.2.1. Let $R$ be an $\mathbb{F}_{p}$-algebra. We define $[\cdot, \cdot]_{R}$ to be the $R$-bilinear form on $\Theta_{0} \otimes_{\mathbb{F}_{p}} R$ obtained from $[\cdot, \cdot]$ by extension of scalars. For any $R$-submodule $\mathcal{L} \subset \Theta_{0} \otimes_{\mathbb{F}_{p}} R$, define

$$
\mathcal{L}^{\operatorname{lin} \perp}:=\left\{v \in \Theta_{0} \otimes_{\mathbb{F}_{p}} R:[v, l]_{R}=0, \forall l \in \mathcal{L}\right\} .
$$

Define $\sigma_{*}(\mathcal{L})$ to be the $R$-module generated by the image of $\mathcal{L}$ under the map

$$
\sigma: \Theta_{0} \otimes_{\mathbb{F}_{p}} R \rightarrow \Theta_{0} \otimes_{\mathbb{F}_{p}} R, \quad v \otimes r \mapsto v \otimes r^{p}
$$

Let $R$ be an $\mathbb{F}_{p^{2}}$-algebra. Let $U$ be an $R$-submodule of $\Omega_{0} \otimes_{\mathbb{F}_{p^{2}}} R$. Since $\Omega_{0} \otimes_{\mathbb{F}_{p^{2}}} R=\Theta_{0} \otimes_{\mathbb{F}_{p}} R$, we may view $U$ as an $R$-submodule of the latter and define $\sigma_{*}(U)$ as in Definition 4.2.1.
Lemma 4.2.2. We have $U^{\perp}=\left(\sigma_{*}(U)\right)^{\text {lin } \perp}$.
Proof. Consider two arbitrary elements

$$
x=\sum_{j} u_{j} \otimes r_{j} \quad \text { and } \quad y=\sum_{k} v_{k} \otimes s_{k}
$$

of $\Theta_{0} \otimes_{\mathbb{F}_{p}} R$. We have

$$
(y, x)_{R}=\sum_{j, k} s_{k} r_{j}^{p} \cdot\left[v_{k}, u_{j}\right]_{R}=\sum_{j, k} s_{k} r_{j}^{p} \cdot\left[u_{j}, v_{k}\right]_{R}=[\sigma(x), y]_{R}
$$

Hence for $y \in \Theta_{0} \otimes_{\mathbb{F}_{p}} R$, we have $y \in U^{\perp}$ if and only if $(y, x)_{R}=0$ for all $x \in U$, if and only if $[\sigma(x), y]_{R}=0$ for all $x \in U$, if and only if $y \in\left(\sigma_{*}(U)\right)^{\operatorname{lin} \perp}$.

Corollary 4.2.3. For any $\mathbb{F}_{p^{2}}$-algebra $R$, the set $\mathcal{V}\left(\Omega_{0}\right)(R)$ is equal to the set of $R$-submodules $U$ of

$$
\Omega_{0} \otimes_{\mathbb{F}_{p^{2}}} R=\Theta_{0} \otimes_{\mathbb{F}_{p}} R,
$$

such that $U$ is an $R$-module local direct summand of rank $d+1$, satisfying

$$
\left(\sigma_{*}(U)\right)^{\operatorname{lin} \perp} \subseteq U
$$

Proof. This is a direct consequence of (2.5.0.1) and Lemma 4.2.2.
In the following we denote $\mathcal{V}(\Lambda)$ by $\mathcal{V}$ for simplicity, where $\Lambda$ is always fixed as in the beginning of Section 4.1. Denote $\Theta:=\Theta_{0} \otimes_{\mathbb{F}_{p}} k$. Fix a point $x_{0} \in \mathcal{V}^{\bar{g}}(k)$. Let $U_{0}$ correspond to $x_{0}$ under (2.5.0.1) or Corollary 4.2.3. Define

$$
\mathcal{L}_{d+1}:=U_{0} \quad \text { and } \quad \mathcal{L}_{d}:=\left(\sigma_{*}\left(U_{0}\right)\right)^{\text {lin } \perp} \stackrel{\text { Lemma }}{=} \stackrel{4.2 .2}{ } U_{0}^{\perp}
$$

They are subspaces of $\Theta$ stable under $\bar{g}$, of $k$-dimensions $d+1$ and $d$ respectively.
Definition 4.2.4. Define $\mathcal{I}:=\mathbb{P}\left(\Theta / \mathcal{L}_{d}\right)$, a projective space of dimension $d$ over $k$.
Then $\mathcal{L}_{d+1}$ defines an element in $\mathcal{I}(k)$, which we still denote by $x_{0}$ by abuse of notation. We have a natural action of $\bar{g}$ on $\mathcal{I}$ that fixes $x_{0}$. Let $\mathcal{R}_{p}$ and $\mathcal{S}_{p}$ be the quotient of the local ring of $\mathcal{I}^{\bar{g}}$ and of $\mathcal{V}^{\bar{g}}$ at $x_{0}$ divided by the $p$-th power of its maximal ideal, respectively.
Lemma 4.2.5. There is a $k$-algebra isomorphism $\mathcal{R}_{p} \cong \mathcal{S}_{p}$.

Proof. The proof is based on exactly the same idea as [Li and Zhu 2017, Lemma 5.2.9]. Let $\tilde{\mathcal{R}}_{p}$ and $\tilde{\mathcal{S}}_{p}$ be the quotient of the local ring of $\mathcal{I}$ and of $\mathcal{V}$ at $x_{0}$ divided by the $p$-th power of its maximal ideal, respectively. Let $R$ be an arbitrary local $k$-algebra with residue field $k$ such that the $p$-th power of its maximal ideal is zero. Then by Lemma 4.2.2, the $R$-points of $\mathcal{V}$ lifting $x_{0}$ classify $R$-module local direct summands $U$ of $\Theta \otimes_{k} R$ of rank $d+1$ that lift $\mathcal{L}_{d+1}$, and such that

$$
U \supseteq\left(\sigma_{*}(U)\right)^{\operatorname{lin} \perp}
$$

But by the assumption that the $p$-th power of the maximal ideal of $R$ is zero, we have

$$
\sigma_{R, *}(U)=\left(\sigma_{k, *}\left(\mathcal{L}_{d+1}\right)\right) \otimes_{k} R,
$$

where we have written $\sigma_{R}$ and $\sigma_{k}$ to distinguish between the Frobenius on $R$ and on $k$. Therefore

$$
\left(\sigma_{R, *}(U)\right)^{\operatorname{lin} \perp}=\left(\left(\sigma_{k, *}\left(\mathcal{L}_{d+1}\right)\right) \otimes_{k} R\right)^{\operatorname{lin} \perp}=\left(\sigma_{k, *}\left(\mathcal{L}_{d+1}\right)\right)^{\operatorname{lin} \perp} \otimes_{k} R=\mathcal{L}_{d} \otimes_{k} R
$$

Thus we see that the set of $R$-points of $\mathcal{V}$ lifting $x_{0}$ is in canonical bijection with the set of $R$-points of $\mathcal{I}$ lifting $x_{0}$. We thus obtain a canonical $\tilde{\mathcal{R}}_{p}$-point of $\mathcal{V}$ lifting $x_{0} \in \mathcal{V}(k)$, and a canonical $\tilde{\mathcal{S}}_{p}$-point of $\mathcal{I}$ lifting $x_{0} \in \mathcal{I}(k)$. These two points induce maps $\tilde{\mathcal{S}}_{p} \rightarrow \tilde{\mathcal{R}}_{p}$ and $\tilde{\mathcal{R}}_{p} \rightarrow \tilde{\mathcal{S}}_{p}$ respectively. From the moduli interpretation of these two maps we see that they are $k$-algebra homomorphisms inverse to each other and equivariant with respect to the actions of $\bar{g}$ on both sides. Note that $\mathcal{S}_{p}$ and $\mathcal{R}_{p}$ are the quotients of $\tilde{\mathcal{S}}_{p}$ and $\tilde{\mathcal{R}}_{p}$ by the augmentation ideal for the $\bar{g}$-action, respectively. It follows that $\mathcal{R}_{p} \cong \mathcal{S}_{p}$.

### 4.3. Study of $\mathcal{I}^{\bar{g}}$.

Definition 4.3.1. Let $\lambda$ be the eigenvalue of $\bar{g}$ on the 1 -dimensional $k$-vector space $\mathcal{L}_{d+1} / \mathcal{L}_{d}=U_{0} / U_{0}^{\perp}$, and let $c$ be the size of the unique (see Remark 4.1.1) Jordan block of $\left.\bar{g}\right|_{\mathcal{L}_{d+1}}$ associated to $\lambda$. Notice our $c$ is denoted by $c+1$ in [RTZ, $\left.\S 9\right]$.
Remark 4.3.2. By the discussion before [Rapoport et al. 2013, Proposition 9.1], $c$ is the size of the unique Jordan block associated to $\lambda$ of $\bar{g}$ on $\Theta / \mathcal{L}_{d}=\Omega / U_{0}^{\perp}$, and is also equal to the quantity $\frac{1}{2}(m(Q(T))+1)$ introduced in the introduction.
Proposition 4.3.3. The local ring $\mathcal{O}_{\mathcal{I}^{\bar{\xi}}, x_{0}}$ of $\mathcal{I}^{\bar{g}}$ at $x_{0}$ is isomorphic to $k[X] / X^{c}$ as a k-algebra.

Proof. By Remark 4.3.2 and Definition 4.2.4, the proposition is a consequence of the following general lemma applied to $\mathcal{L}=\Theta / \mathcal{L}_{d}$ and $h=\bar{g}$.

Lemma 4.3.4. Let $\mathcal{L}$ be a $k$-vector space of dimension $d+1$. Let $\mathbb{P}(\mathcal{L})=\mathbb{P}^{d}$ be the associated projective space. Let $x_{0} \in \mathbb{P}(\mathcal{L})(k)$, represented by a vector $\ell \in \mathcal{L}$. Let $h \in \operatorname{GL}(\mathcal{L})(k)=\mathrm{GL}_{d+1}(k)$. Assume that:
(1) The natural action of $h$ on $\mathbb{P}(\mathcal{L})$ fixes $x_{0}$. Denote the eigenvalue of $h$ on $\ell$ by $\lambda$.
(2) There is a unique Jordan block of $h$ associated to the eigenvalue $\lambda$. Denote its size by $c$.
Let $R:=\mathcal{O}_{\mathbb{P}(\mathcal{L})^{h}, x_{0}}$ be the local ring of the fixed point scheme $\mathbb{P}(\mathcal{L})^{h}$ at $x_{0}$. Then

$$
R \cong k[X] / X^{c}
$$

Proof. Extend $\ell$ to a basis $\left\{\ell_{0}=\ell, \ell_{1}, \ldots, \ell_{d}\right\}$ of $\mathcal{L}$ such that the matrix $\left(h_{i j}\right)_{0 \leq i, j \leq d}$ of $h$ under this basis is in the Jordan normal form. Under this basis, the point $x_{0}$ has projective coordinates $\left[X_{0}: \cdots: X_{d}\right]=[1: 0: \cdots: 0] \in \mathbb{P}^{d}$. Let $Z_{i}=X_{i} / X_{0}$ $(1 \leq i \leq d)$ and let $\mathbb{A}^{d}$ be the affine space with coordinates $\left(Z_{1}, \ldots, Z_{d}\right)$. Then we can identify the local ring of $\mathbb{P}^{d}$ at $x_{0}$ with the local ring of $\mathbb{A}^{d}$ at the origin. Since $h$ fixes $x_{0}$, we know that $h$ acts on the local ring of $\mathbb{A}^{d}$ at the origin (although $h$ does not stabilize $\mathbb{A}^{d}$ in general). Since $\left(h_{i j}\right)$ is in the Jordan normal form, we know that the action of $h$ on the latter is given explicitly by

$$
h Z_{i}=\frac{h_{i, i} X_{i}+h_{i, i+1} X_{i+1}}{h_{0,0} X_{0}+h_{0,1} X_{1}}=\frac{h_{i, i} Z_{i}+h_{i, i+1} Z_{i+1}}{h_{0,0}+h_{0,1} Z_{1}}, \quad 1 \leq i \leq d
$$

where $h_{i, i+1} Z_{i+1}$ is understood as 0 when $i=d$. Hence the local equations at the origin of $\mathbb{A}^{d}$ which cut out the $h$-fixed point scheme are given by

$$
\left(h_{0,0}-h_{i, i}\right) Z_{i}+h_{0,1} Z_{1} Z_{i}=h_{i, i+1} Z_{i+1}, \quad 1 \leq i \leq d
$$

By hypothesis (2), we have $h_{0,0}-h_{i, i} \neq 0$ if and only if $i \geq c$. Thus when $i \geq c$, we know that $\left(h_{0,0}-h_{i, i}\right)+h_{0,1} Z_{1}$ is a unit in the local ring of $\mathbb{A}^{d}$ at the origin, and so $Z_{i}$ can be solved as a multiple of $h_{i, i+1} Z_{i+1}$ when $i \geq c$. It follows that

$$
Z_{i}=0, \quad i \geq c
$$

If $c=1$, then $Z_{1}=\cdots=Z_{d}=0$ and the local ring $R$ in question is isomorphic to $k$ as desired. If $c>1$, then $h_{0,1}=1$ and we find the equations for $i=1, \cdots, c-1$ simplify to

$$
Z_{1} Z_{1}=Z_{2}, Z_{1} Z_{2}=Z_{3}, \cdots, Z_{1} Z_{c-2}=Z_{c-1}, Z_{1} Z_{c-1}=0
$$

Hence the local ring $R$ in question is isomorphic to (the localization at the ideal $\left(Z_{1}, Z_{2}, \cdots, Z_{c-1}\right)$ or ( $Z_{1}$ ) of)

$$
k\left[Z_{1}, Z_{2}, \ldots, Z_{c-1}\right] /\left(Z_{1}^{2}-Z_{2}, Z_{1}^{3}-Z_{3}, \cdots, Z_{1}^{c-1}-Z_{c-1}, Z_{1}^{c}\right) \cong k\left[Z_{1}\right] / Z_{1}^{c}
$$

as desired.
Theorem 4.3.5. Let $g \in J\left(\mathbb{Q}_{p}\right)$ be regular semisimple and minuscule. Let $x_{0}$ be a point in $\left(\delta(\mathcal{M}) \cap \mathcal{N}^{g}\right)(k)$. Also denote by $x_{0}$ the image of $x_{0}$ in $\mathcal{V}(\Lambda)(k)$ as in

Proposition 4.1.2 and define $\lambda$ and $c$ as in Definition 4.3.1. Assume $p>c$. Then the complete local ring of $\delta(\mathcal{M}) \cap \mathcal{N}^{g}$ at $x_{0}$ is isomorphic to $k[X] / X^{c}$.
Proof. Let $\hat{\mathcal{S}}$ be the complete local ring of $\delta(\mathcal{M}) \cap \mathcal{N}^{g}$ at $x_{0}$. By Proposition 4.1.2 and by the fact that $\mathcal{V}(\Lambda)$ is smooth of dimension $d$ (Section 2.5), we know that $\hat{\mathcal{S}}$ is a quotient of the power series ring $k \llbracket X_{1}, \cdots, X_{d} \rrbracket$. By Proposition 4.1.2, Lemma 4.2.5 and Proposition 4.3.3, we know that $\hat{\mathcal{S}} / \mathfrak{m}_{\hat{\mathcal{S}}}^{p}$ is isomorphic to $k[X] / X^{c}$ as a $k$-algebra. In such a situation, it follows from the next abstract lemma that $\hat{\mathcal{S}} \cong k[X] / X^{c}$.
Lemma 4.3.6. Let I be a proper ideal of $k \llbracket X_{1}, \cdots, X_{d} \rrbracket$ and let

$$
\hat{\mathcal{S}}=k \llbracket X_{1}, \cdots, X_{d} \rrbracket / I .
$$

Let $\mathfrak{m}$ be the maximal ideal of $k \llbracket X_{1}, \cdots, X_{d} \rrbracket$ and let $\mathfrak{m}_{\hat{\mathcal{S}}}$ be the maximal ideal of $\hat{\mathcal{S}}$. Assume there is a $k$-algebra isomorphism $\beta: \hat{\mathcal{S}} / \mathfrak{m}_{\hat{\mathcal{S}}}^{p} \xrightarrow{\sim} k[X] / X^{c}$ for some integer $1 \leq c<p$. Then $\hat{\mathcal{S}}$ is isomorphic to $k[X] / X^{c}$ as a $k$-algebra.
Proof. We first notice that if $R_{1}$ is any quotient ring of $k \llbracket X_{1}, \cdots, X_{d} \rrbracket$ with its maximal ideal $\mathfrak{m}_{1}$ satisfying $\mathfrak{m}_{1}=\mathfrak{m}_{1}^{2}$ (i.e., $R_{1}$ has zero cotangent space), then $R_{1}=k$. In fact, $R_{1}$ is noetherian and we have $\mathfrak{m}_{1}^{l}=\mathfrak{m}_{1}$ for all $l \in \mathbb{Z}_{\geq 1}$, so by Krull's intersection theorem we conclude that $\mathfrak{m}_{1}=0$ and $R_{1}=k$.

Suppose $c=1$. Then $\hat{\mathcal{S}} / \mathfrak{m}_{\hat{\mathcal{S}}}^{p} \cong k$, so $\hat{\mathcal{S}}$ has zero cotangent space and thus $\hat{\mathcal{S}}=k$ as desired. Next we treat the case $c \geq 2$. Let $\alpha$ be the composite

$$
\alpha: k \llbracket X_{1}, \cdots, X_{d} \rrbracket \rightarrow \hat{\mathcal{S}} / \mathfrak{m}_{\hat{\mathcal{S}}}^{p} \xrightarrow{\beta} k[X] / X^{c} .
$$

Let $J=\operatorname{ker} \alpha$. Since $\alpha$ is surjective, we reduce to prove that $I=J$. Note that because $\beta$ is an isomorphism we have

$$
\begin{equation*}
I+\mathfrak{m}^{p}=J \tag{4.3.6.1}
\end{equation*}
$$

In the following we prove $\mathfrak{m}^{p} \subset I$, which will imply $I=J$ and hence the lemma. The argument is a variant of [RTZ, Lemma 11.1].

Let $Y \in k \llbracket X_{1}, \cdots, X_{d} \rrbracket$ be such that $\alpha(Y)=X$. Since $X$ generates the maximal ideal in $k[X] / X^{c}$, we have

$$
\begin{equation*}
\mathfrak{m}=J+(Y) \tag{4.3.6.2}
\end{equation*}
$$

Then by (4.3.6.1) and (4.3.6.2) we have $\mathfrak{m}=I+(Y)+\mathfrak{m}^{p}$, and so the local ring $k \llbracket X_{1}, \cdots, X_{d} \rrbracket /(I+(Y))$ has zero cotangent space. We have observed that the cotangent space being zero implies that the ring has to be $k$, or equivalently

$$
\begin{equation*}
\mathfrak{m}=I+(Y) \tag{4.3.6.3}
\end{equation*}
$$

Now we start to show $\mathfrak{m}^{p} \subset I$. By (4.3.6.3) we have $\mathfrak{m}^{p} \subset I+\left(Y^{p}\right)$, so we only need to prove $Y^{p} \in I$. We will show the stronger statement that $Y^{c} \in I$. By Krull's
intersection theorem, it suffices to show that $Y^{c} \in I+\mathfrak{m}^{p l}$ for all $l \geq 1$. In the following we show this by induction on $l$.

Assume $l=1$. Note that $\alpha\left(Y^{c}\right)=0$, so by (4.3.6.1) we have

$$
Y^{c} \in J=I+\mathfrak{m}^{p} .
$$

Suppose $Y^{c} \in I+\mathfrak{m}^{p l}$ for an integer $l \geq 1$. Write

$$
\begin{equation*}
Y^{c}=i+m, \quad i \in I, m \in \mathfrak{m}^{p l} . \tag{4.3.6.4}
\end{equation*}
$$

By (4.3.6.2) we know

$$
\mathfrak{m}^{p l} \subset(J+(Y))^{p l} \subset \sum_{s=0}^{p l} J^{s}(Y)^{p l-s} .
$$

Thus we can decompose $m \in \mathfrak{m}^{p l}$ into a sum

$$
\begin{equation*}
m=\sum_{s=0}^{p l} j_{s} Y^{p l-s}, \quad j_{s} \in J^{s} \tag{4.3.6.5}
\end{equation*}
$$

By (4.3.6.4) and (4.3.6.5), we have

$$
Y^{c}=i+\sum_{s=0}^{p l} j_{s} Y^{p l-s}
$$

and so

$$
\begin{equation*}
Y^{c}-\sum_{s=0}^{p l-c} j_{s} Y^{p l-s}=i+\sum_{s=p l-c+1}^{p l} j_{s} Y^{p l-s} . \tag{4.3.6.6}
\end{equation*}
$$

Denote

$$
A:=\sum_{s=0}^{p l-c} j_{s} Y^{p l-s-c} .
$$

Then the left hand side of (4.3.6.6) is equal to $(1-A) Y^{c}$. Hence we have

$$
\begin{aligned}
(1-A) Y^{c}=i+ & \sum_{s=}^{p l} j_{s} Y^{p l-c+1} \subset I+J^{p l-c+1} \\
& \xrightarrow{(4.3 .6 .1)} I+\left(I+\mathfrak{m}^{p}\right)^{p l-c+1}=I+\mathfrak{m}^{p(p l-c+1)} \subset I+\mathfrak{m}^{p(l+1)},
\end{aligned}
$$

where for the last inclusion we have used $c<p$. Since $1-A$ is a unit in $k \llbracket X_{1}, \cdots, X_{d} \rrbracket$ (because $c<p$ ), we have $Y^{c} \in I+\mathfrak{m}^{p(l+1)}$. By induction, $Y^{c} \in I+\mathfrak{m}^{p l}$ for all $l \in \mathbb{Z}_{\geq 1}$, as desired.

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[^0]:    MSC2010: primary 11G18; secondary 11R39, 14C17, 14C25, 14G35.
    Keywords: Supersingular locus, Special fiber of Shimura varieties, Deligne-Lusztig varieties, Tate conjecture.

[^1]:    ${ }^{1}$ Strictly speaking, the moduli space $\mathcal{S h}(G)_{K}$ is $\# \operatorname{ker}^{1}(\mathbb{Q}, G)$-copies of the classical Shimura variety whose $\mathbb{C}$-points are given by the double coset space $G(\mathbb{Q}) \backslash G(\mathbb{A}) / K_{\infty} K$, where $K_{\infty} \subseteq G(\mathbb{R})$ is the maximal compact subgroup modulo center. See [Kottwitz 1992b, page 400] for details.

[^2]:    ${ }^{2}$ By a $(1,1)$-elementary matrix, we mean an $n \times n$-matrix whose $(1,1)$-entry is 1 and whose other entries are zero.

[^3]:    ${ }^{3}$ As explained in the proof of [Harris and Taylor 2001, Lemma I.7.1], when $n$ is odd, such $\beta_{a_{\bullet}}$ always exists, and when $n$ is even, existence of $\beta_{a_{\bullet}}$ depends on the parity of $a_{1}+\cdots+a_{f}$. See also the proof of Lemma 2.9.

[^4]:    ${ }^{4}$ Although one can descend $\mathcal{S} h_{a_{\bullet}}$ to the subring $\mathcal{O}_{E_{h, \wp}}$ of $\mathbb{Z}_{p}$, we ignore this minor improvement.

[^5]:    ${ }^{5}$ This automatically implies that $\pi_{\infty}$ has the same central and infinitesimal characters as the contragradient of $\xi_{\mathbb{C}}$.
    ${ }^{6}$ Rigorously speaking, Kottwitz's theorem describes the direct sum of the $\pi$-component of all cohomological degrees. Since our $\pi_{p}$ is tempered, so $\pi$ appears only in the middle degree for purity reasons because $\mathrm{Sh}_{a_{\bullet}}$ is compact.

[^6]:    ${ }^{7}$ Conjecturally, the Frobenius action on the étale $\ell$-adic cohomology groups of a projective smooth variety over a finite field is always semisimple.
    ${ }^{8}$ For example, if $r=1$ and $\alpha_{1}=\alpha_{2}$, the eigenvalues $\alpha_{1} \cdot \alpha_{1} \alpha_{3} \alpha_{4} \cdots \alpha_{n}$ is equal to $\alpha_{1} \cdots \alpha_{n}$ and hence is $p^{n(n-1)}$ times a root of unity. So to be in the generic case, we will need to require that $\alpha_{i} / \alpha_{j}$ for $i \neq j$ is not a root of unity if $r=1$. For another example, if $r=2$, "generic" will mean that $\alpha_{i} / \alpha_{j}$ for $i \neq j$ and $\alpha_{i} \alpha_{i^{\prime}} / \alpha_{j} \alpha_{j^{\prime}}$ for $\left\{i, i^{\prime}\right\} \neq\left\{j, j^{\prime}\right\}$ are not roots of unity.

[^7]:    ${ }^{9}$ This isomorphism depends on the choice of the isomorphism $\gamma_{a_{\bullet}, b_{\bullet}}$ made earlier.
    ${ }^{10}$ This assumption is satisfied when $\pi$ is the finite part of an automorphic cuspidal representation of $G_{a_{\bullet}}(\mathbb{A})$ which admits a base change to a cuspidal automorphic representation of $\mathrm{GL}_{n}\left(\mathbb{A}_{E}\right) \times \mathbb{A}_{E}^{\times}$. Indeed, in this case, White [2012, Theorem E] proved that $m_{a_{\bullet}}(\pi)=m_{b_{\bullet}^{(j)}}(\pi)=1$.
    ${ }^{11}$ The $\pi$-isotypic component is the same as the $\pi^{p}$-isotypic component according to Lemma 4.17.

[^8]:    ${ }^{12}$ The Siegel varieties are Shimura varieties associated to $\mathrm{GSp}_{2 g}(\mathbb{Q})$. The Langlands dual group is isogenous to $\operatorname{Spin}(2 g+1)$ and the associated representation $r_{\mu}$ is the spin representation, which is minuscule and hence does not contain trivial weight subspace.

[^9]:    ${ }^{13}$ As pointed out above, we have to work with the Hilbert-Siegel setup as opposed to the usual Siegel setup because $r_{\text {spin }}$ is a minuscule representation.

[^10]:    ${ }^{14}$ Here and after, by a subbundle of a locally free coherent sheaf, we mean a locally free coherent sheaf that is Zariski locally a direct factor.
    ${ }^{15}$ The notation $F$ for Frobenius was also used to denote the real quadratic field. But we think the chance for confusion is minimal.

[^11]:    ${ }^{16}$ Note that the two skew-Hermitian forms $(H,\langle-,-\rangle) \tilde{\lambda}$ and $\left(V_{0, n},\langle-,-\rangle_{0, n}\right)$ are not necessarily isomorphic over $\mathbb{Q}$. However, they differ at most only by a scalar in $F$, hence define the same similitude unitary group. See [Kottwitz 1992b, p. 400] for details.

[^12]:    ${ }^{17}$ We point out that, for (4.14.1), $F$ is an isomorphism if and only if $V$ is an isomorphism, because this is equivalent to requiring the source and the target to have the same dimension.

[^13]:    ${ }^{18}$ This is in fact a corollary of (2) and (4).

[^14]:    ${ }^{19}$ We may also view elements of $\mathscr{H}_{K}$ as $\mathbb{Z}$-valued locally constant and compactly supported functions on $\mathrm{GL}_{n}(F)$ which are bi-invariant under $K$, and define the product of $f, g \in \mathscr{H}_{K}$ as $(f * g)(x)=\int_{\mathrm{GL}_{n}(F)} f(y) g\left(y^{-1} x\right) d y$, where $d y$ means the unique bi-invariant Haar measure on $\mathrm{GL}_{n}(F)$ with $\int_{K} d y=1$. For the equivalence between these two definitions, see [Gross 1998, p. 4].

[^15]:    MSC2010: primary 11G18; secondary 11G35, 11E57, 20 G 30.
    Keywords: Shimura varieties, conjugation, models.

[^16]:    MSC2010: primary 11J97; secondary 11J87, 14G05.
    Keywords: Schmidt's subspace theorem, Roth's theorem, Diophantine approximation, Vojta's conjecture.

[^17]:    MSC2010: primary 11R23; secondary 11F33.
    Keywords: Iwasawa theory, Hida theory, Selmer groups, Heegner points, special values of
    L-functions.

[^18]:    I am grateful to the anonymous referee for helpful comments that have improved the exposition.
    MSC2010: primary 11F80; secondary 11R37.
    Keywords: Galois representations, Kuga-Satake construction.

[^19]:    ${ }^{1}$ For some, admittedly limited, examples, see [Patrikis 2016a; 2016b].

[^20]:    ${ }^{2}$ To see the relevance of bounding the coefficients, the reader may contrast the case of elliptic curves (over $\mathbb{Q}$, say) unramified outside $S$ with that of all weight-two modular forms unramified outside $S$ : of the former there are finitely many isogeny classes, since the conductor is bounded, whereas the latter can have level divisible by arbitrarily high powers of the primes in $S$.

[^21]:    ${ }^{3}$ To be precise, this depends on $\iota_{\infty}$ and $\iota_{\lambda}$; but $\iota_{\infty}$ is fixed throughout the paper.

[^22]:    ${ }^{4}$ Note that $\phi$ is valued in $C[m]$, not $H_{1} \cap C$, so $\delta(\phi)$ need not be a coboundary in $Z^{2}\left(\Gamma_{F_{v}}, H_{1} \cap C\right)$.

[^23]:    ${ }^{5}$ In concrete terms, this says that if an element of $F^{\times}$is everywhere locally a ( $2 m$ )-th power, then it is globally an $m$-th power.
    ${ }^{6}$ This is easy to make explicit, using local duality, in terms of $\mu_{\infty}\left(F_{v}\right)$.

[^24]:    ${ }^{7}$ Of course, we only need to consider the cokernel of the map to the direct sum over $v \in S_{\lambda} \backslash\left(S_{\lambda} \cap(S \cup T)\right)$; to lighten the notation we will work with all $v \in S_{\lambda}$, taking some arbitrary (e.g., trivial) choice of $\chi_{\lambda, v}$ at any places in $(S \cup T) \cap S_{\lambda}$.

[^25]:    MSC2010: primary 11G18; secondary 14G17, 22E55.
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    Rapoport-Zink spaces, special cycles.

[^26]:    ${ }^{1}$ Such a pair $\left(N_{0}, \Phi\right)$ is sometimes called a relative isocrystal.

