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Dedicated to the memory of Jon Rogawski


#### Abstract

This article presents a new proof of a theorem of Karl Rubin relating values of the Katz $p$-adic $L$-function of an imaginary quadratic field at certain points outside its range of classical interpolation to the formal group logarithms of rational points on CM elliptic curves. The approach presented here is based on the $\boldsymbol{p}$-adic Gross-Zagier type formula proved by the three authors in previous work. As opposed to the original proof which relied on a comparison between Heegner points and elliptic units, it only makes use of Heegner points, and leads to a mild strengthening of Rubin's original result. A generalization to the case of modular abelian varieties with complex multiplication is also included.


## 1. Introduction

The aim of this article is to present a new proof of a theorem of Karl Rubin (see [Rubin 1992] and Theorem 1 below) relating values of the Katz $p$-adic $L$ function of an imaginary quadratic field at certain points outside its range of classical interpolation to the formal group logarithms of rational points on CM elliptic curves. This theorem has been seminal in providing a motivation for PerrinRiou's formulation 1993; 2000 of the $p$-adic Beilinson conjectures. The new proof described in this work is based on the $p$-adic Gross-Zagier type formula of [Bertolini et al. 2012b], and only makes use of Heegner points as opposed to the original proof, which relied on a comparison between Heegner points and elliptic units. Hence, it should be adaptable to more general situations, for example to the setting of general CM fields.

Let $A$ be an elliptic curve over $\mathbb{Q}$ with complex multiplication by the ring of integers of a quadratic imaginary field $K$. A classical result of Deuring identifies the Hasse-Weil $L$-series $L(A, s)$ of $A$ with the $L$-series $L\left(v_{A}, s\right)$ attached to a

[^0]Hecke character $v_{A}$ of $K$ of infinity type $(1,0)$. When $p$ is a prime which splits in $K$ and does not divide the conductor of $A$, the Hecke $L$-function $L\left(v_{A}, s\right)$ has a $p$-adic analog, namely the Katz two-variable $p$-adic $L$-function attached to $K$. It is a $p$-adic analytic function, denoted by $v \mapsto \mathscr{L}_{p}(\nu)$, on the space of Hecke characters equipped with its natural $p$-adic analytic structure. Section 3A recalls the definition of this $L$-function: the values $\mathscr{L}_{p}(\nu)$ at Hecke characters of infinity type ( $1+j_{1},-j_{2}$ ) with $j_{1}, j_{2} \geq 0$ are defined by interpolation of the classical $L$-values $L\left(v^{-1}, 0\right)$. Letting $v^{*}:=v \circ c$, where $c$ denotes complex conjugation on the ideals of $K$, it is readily seen by comparing Euler factors that $L(v, s)=L\left(\nu^{*}, s\right)$. A similar equality need not hold in the $p$-adic setting, because the involution $v \mapsto \nu^{*}$ corresponds to the map $\left(j_{1}, j_{2}\right) \mapsto\left(j_{2}, j_{1}\right)$ on weight space and therefore does not preserve the lower right quadrant of weights of Hecke characters that lie in the range of classical interpolation. Since $\nu_{A}$ lies in the domain of classical interpolation, the $p$-adic $L$-value $\mathscr{L}_{p}\left(v_{A}\right)$ is a simple multiple of $L\left(v_{A}^{-1}, 0\right)=L(A, 1)$. Suppose that it vanishes. (The Birch and Swinnerton-Dyer conjecture predicts then that $A(\mathbb{Q})$ is infinite; this is known to be true when the order of vanishing is exactly one.) The value $\mathscr{L}_{p}\left(v_{A}^{*}\right)$ is a second, a priori more mysterious $p$-adic avatar of the leading term of $L(A, s)$ at $s=1$. Rubin's theorem gives a formula for this quantity in the analytic rank-one case:

Theorem 1 [Rubin 1992]. Let $v_{A}$ be a Hecke character of type (1, 0) attached to an elliptic curve $A / \mathbb{Q}$ with complex multiplication, and suppose that $L(A, s)$ vanishes to order one at the central point $s=1$. Then there exists a global point $P \in A(\mathbb{Q})$ of infinite order such that

$$
\begin{equation*}
\mathscr{L}_{p}\left(v_{A}^{*}\right)=\Omega_{p}(A)^{-1} \log _{\omega_{A}}(P)^{2}\left(\bmod K^{\times}\right), \tag{1-1}
\end{equation*}
$$

where

- $\Omega_{p}(A)$ is the p-adic period attached to $A$ as in Section 2C;
- $\omega_{A} \in \Omega^{1}(A / \mathbb{Q})$ is a regular differential on $A$ over $\mathbb{Q}$, and $\log _{\omega_{A}}: A\left(\mathbb{Q}_{p}\right) \rightarrow \mathbb{Q}_{p}$ denotes the p-adic formal group logarithm with respect to $\omega_{A}$.
(For a more precise statement without the $K^{\times}$ambiguity, see [Rubin 1992].) Formula (1-1) is peculiar to the $p$-adic world and suggests that $p$-adic $L$-functions encode arithmetic information that is not readily apparent in their complex counterparts.

Rubin's proof of Theorem 1 breaks up naturally into three parts:
(1) He exploits the Euler system of elliptic units to construct a global cohomology class $\kappa_{A}$ belonging to a pro- $p$ Selmer group $\operatorname{Sel}_{p}(A / \mathbb{Q})$ attached to $A$. The close connection between elliptic units and the Katz $p$-adic $L$-function is then parlayed into the explicit evaluation of two natural $p$-adic invariants attached to $\kappa_{A}$ : the $p$-adic
formal group logarithm $\log _{A, p}\left(\kappa_{A}\right)$ and the cyclotomic $p$-adic height $\left\langle\kappa_{A}, \kappa_{A}\right\rangle$ :

$$
\begin{align*}
\log _{A, p}\left(\kappa_{A}\right) & =\left(1-\beta_{p}^{-1}\right)^{-1} \mathscr{L}_{p}\left(v_{A}^{*}\right) \Omega_{p}(A),  \tag{1-2}\\
\left\langle\kappa_{A}, \kappa_{A}\right\rangle & =\left(1-\alpha_{p}^{-1}\right)^{-2} \mathscr{L}_{p}^{\prime}\left(v_{A}\right) \mathscr{L}_{p}\left(v_{A}^{*}\right), \tag{1-3}
\end{align*}
$$

where

- $\alpha_{p}$ and $\beta_{p}$ denote the roots of the Hasse polynomial $x^{2}-a_{p}(A) x+p$, ordered in such a way that $\operatorname{ord}_{p}\left(\alpha_{p}\right)=0$ and $\operatorname{ord}_{p}\left(\beta_{p}\right)=1$;
- the quantity $\mathscr{L}_{p}^{\prime}\left(\nu_{A}\right)$ denotes the derivative of $\mathscr{L}_{p}$ at $v_{A}$ in the direction of the cyclotomic character.
(2) Independently of the construction of $\kappa_{A}$, the theory of Heegner points can be used to construct a canonical point $P \in A(\mathbb{Q})$, which is of infinite order when $L\left(v_{A}, s\right)=$ $L(A, s)$ vanishes to exact order one at $s=1$. The $\operatorname{Selmer} \operatorname{group} \operatorname{Sel}_{p}(A / \mathbb{Q}) \otimes \mathbb{Q}$ is of rank one (by the results of Kolyvagin) and the image $\kappa_{P} \in \operatorname{Sel}_{p}(A / \mathbb{Q})$ of $P$ under the connecting homomorphism of Kummer theory supplies us with a generator for $\operatorname{Sel}_{p}(A / \mathbb{Q}) \otimes \mathbb{Q}$. Furthermore, the $p$-adic analog of the Gross-Zagier formula proved by Perrin-Riou [1987] shows that

$$
\begin{equation*}
\left\langle\kappa_{P}, \kappa_{P}\right\rangle=\mathscr{L}_{p}^{\prime}\left(v_{A}\right) \Omega_{p}(A)^{-1}\left(\bmod K^{\times}\right) \tag{1-4}
\end{equation*}
$$

Finally, a theorem of Bertrand shows that the $p$-adic height pairing is nondegenerate in the above situation, that is, $\left\langle\kappa_{P}, \kappa_{P}\right\rangle \neq 0$. In particular, one concludes from (1-4) that $\mathscr{L}_{p}^{\prime}\left(\nu_{A}\right) \neq 0$.
(3) Using that $\mathscr{L}_{p}^{\prime}\left(\nu_{A}\right)$ is nonzero, Rubin shows that $\kappa_{A}$ is nonzero in $\operatorname{Sel}_{p}(A / \mathbb{Q}) \otimes \mathbb{Q}$ and therefore is a second generator of this one-dimensional $\mathbb{Q}_{p}$-vector space. (See Theorem 8.1 and Corollary 8.3 of [Rubin 1992].) Equations (1-2) and (1-3) then show that $\mathscr{L}_{p}\left(v_{A}^{*}\right) \neq 0$, and further, for any generator $\kappa$ of the $\mathbb{Q}_{p}$-vector space $\operatorname{Sel}_{p}(A / \mathbb{Q}) \otimes \mathbb{Q}$, one has

$$
\begin{equation*}
\frac{\log _{A, p}^{2}(\kappa)}{\langle\kappa, \kappa\rangle}=\frac{\left(1-\beta_{p}^{-1}\right)^{-2} \mathscr{L}_{p}\left(v_{A}^{*}\right) \Omega_{p}(A)^{2}}{\left(1-\alpha_{p}^{-1}\right)^{-2} \mathscr{L}_{p}^{\prime}\left(v_{A}\right)}, \tag{1-5}
\end{equation*}
$$

since the quantity on the left-hand side does not depend on the choice of $\kappa$. Rubin obtains Theorem 1 by setting $\kappa=\kappa_{P}$ in (1-5) and using (1-4) to eliminate the quantities $\left\langle\kappa_{P}, \kappa_{P}\right\rangle$ and $\mathscr{L}_{p}^{\prime}\left(\nu_{A}\right)$.

The reader will note the key role that is played in Rubin's proof by both the Euler systems of elliptic units and of Heegner points. The new approach to Theorem 1 described in this paper relies solely on Heegner points, and requires neither elliptic units nor Perrin-Riou's $p$-adic height calculations. Instead, the key ingredient in this approach is the $p$-adic variant of the Gross-Zagier formula, arising from the results
of [Bertolini et al. 2012b], which is stated in Theorem 3.12. This formula expresses $p$-adic logarithms of Heegner points in terms of the special values of a $p$-adic Rankin $L$-function attached to a cusp form $f$ and an imaginary quadratic field $K$, and may be of some independent interest insofar as it exhibits a strong analogy with Rubin's formula but applies to arbitrary - not necessarily CM - elliptic curves over $\mathbb{Q}$. When $f$ is the theta series attached to a Hecke character of $K$, Theorem 1 follows from the factorization of the associated $p$-adic Rankin $L$-function into a product of two Katz $L$-functions, a factorization which is a simple manifestation of the Artin formalism for these $p$-adic $L$-series.

One might expect that the statement of Theorem 1 should generalize to the setting where $v_{A}$ is replaced by an algebraic Hecke character $v$ of infinity type $(1,0)$ of a quadratic imaginary field $K$ (of arbitrary class number) satisfying

$$
\begin{equation*}
\left.\nu\right|_{A_{\mathbb{Q}}}=\varepsilon_{K} \cdot N, \tag{1-6}
\end{equation*}
$$

where $\varepsilon_{K}$ is the quadratic Dirichlet character associated to $K / \mathbb{Q}$ and $N: \mathbb{A}_{\mathbb{Q}}^{\times} \rightarrow \mathbb{R}^{\times}$ is the adelic norm character. Chapter 3 treats this more general setting, which (although probably amenable as well to the original methods of [Rubin 1992]) is not yet covered in the literature. Assumption (1-6) implies that the classical functional equation relates $L(v, s)$ to $L(v, 2-s)$. Assume further that the sign $w_{v}$ in this functional equation satisfies

$$
\begin{equation*}
w_{v}=-1, \tag{1-7}
\end{equation*}
$$

so that $L(\nu, s)$ vanishes to odd order at $s=1$. For technical reasons, it will also be convenient to make two further assumptions. Firstly, we assume that

$$
\begin{equation*}
\text { the discriminant }-D \text { of } K \text { is odd. } \tag{1-8}
\end{equation*}
$$

Secondly, we note that assumption (1-6) implies that $\mathfrak{d}_{K}:=\sqrt{-D}$ necessarily divides the conductor of $\nu$, and we further restrict the setting by imposing the assumption that the conductor of $v$ is exactly divisible by $\mathfrak{d}_{K}$.

The statement of Theorem 2 below requires some further notions, which we now introduce. Let $E_{\nu}$ be the subfield of $\mathbb{C}$ generated by the values of the Hecke character $\nu$, and let $T_{\nu}$ be its ring of integers. A general construction which is recalled in Sections 2B and 3F attaches to $v$ an abelian variety $B_{v}$ over $K$ of dimension $\left[E_{v}: K\right]$, equipped with inclusions

$$
T_{v} \subset \operatorname{End}_{K}\left(B_{v}\right), \quad E_{v} \subset \operatorname{End}_{K}\left(B_{v}\right) \otimes \mathbb{Q} .
$$

Given $\lambda \in T_{\nu}$, denote by $[\lambda]$ the corresponding endomorphism of $B_{\nu}$, and set

$$
\begin{align*}
\Omega^{1}\left(B_{v} / E_{\nu}\right)^{T_{v}} & :=\left\{\omega \in \Omega^{1}\left(B_{v} / E_{v}\right) \mid[\lambda]^{*} \omega=\lambda \omega, \text { for all } \lambda \in T_{\nu}\right\},  \tag{1-10}\\
\left(B_{v}(K) \otimes E_{v}\right)^{T_{v}} & :=\left\{P \in B_{v}(K) \otimes_{\mathbb{Z}} E_{v} \mid[\lambda] P=\lambda P, \text { for all } \lambda \in T_{v}\right\} . \tag{1-11}
\end{align*}
$$

The vector space $\Omega^{1}\left(B_{v} / E_{v}\right)^{T_{v}}$ is one-dimensional over $E_{v}$.
After fixing a $p$-adic embedding $K \subset \mathbb{Q}_{p}$, the formal group logarithm on $B_{v}$ gives rise to a bilinear pairing

$$
\left.\begin{array}{rl}
\langle,\rangle: \Omega^{1}\left(B_{v} / K\right) \times & B_{v}(K)
\end{array}\right) \mathbb{Q}_{p}, \quad, \quad(\omega, P) \mapsto \log _{\omega} P,
$$

satisfying $\left\langle[\lambda]^{*} \omega, P\right\rangle=\langle\omega,[\lambda] P\rangle$ for all $\lambda \in T_{\nu}$. This pairing can be extended by $E_{\nu}$-linearity to an $E_{\nu} \otimes \mathbb{Q}_{p}$-valued pairing between $\Omega^{1}\left(B_{\nu} / E_{\nu}\right)$ and $B_{\nu}(K) \otimes E_{\nu}$. When $\omega$ and $P$ belong to these $E_{\nu}$-vector spaces, we will continue to write $\log _{\omega}(P)$ for $\langle\omega, P\rangle$.

Theorem 2. Let v be an algebraic Hecke character of infinity type $(1,0)$ satisfying (1-6), (1-7), (1-8) and (1-9) above. Then there exists $P_{v} \in B_{v}(K) \otimes \mathbb{Q}$ such that

$$
\mathscr{L}_{p}\left(\nu^{*}\right)=\Omega_{p}\left(v^{*}\right)^{-1} \log _{\omega_{v}}\left(P_{\nu}\right)^{2}\left(\bmod E_{v}^{\times}\right),
$$

where $\Omega_{p}\left(v^{*}\right) \in \mathbb{C}_{p}$ is the $p$-adic period attached to $v$ in Definition 2.13, and $\omega_{\nu}$ is a nonzero element of $\Omega^{1}\left(B_{\nu} / E_{\nu}\right)^{T_{v}}$. The point $P_{\nu}$ is nonzero if and only if $L^{\prime}(\nu, 1) \neq 0$.

Remark 3. Assumptions (1-8) and (1-9) could certainly be relaxed with more work. For instance, (1-8) is needed since the main theorem of [Bertolini et al. 2012b] is only proved for imaginary quadratic fields of odd discriminant. Likewise, removing (1-9) would require generalizing the cited result to the case of Shimura curves over $\mathbb{Q}$.

Remark 4. In [Bertolini et al. 2012c], we give a conjectural construction of rational points on CM elliptic curves (called Chow-Heegner points) using cycles on higher-dimensional varieties. While this construction of points is contingent on a certain case of the Tate conjecture, the corresponding construction at the level of cohomology classes can be made unconditionally. The results of this paper (especially Theorem 2 above) combined with those of [Bertolini et al. 2012b] are used in [Bertolini et al. 2012c] to establish that these cohomology classes indeed correspond to global points via the Kummer map.

The careful reader will notice that the hypothesis in Theorem 2 on the order of vanishing of $L(v, s)$ is weaker than that in [Rubin 1992], since $L(v, s)$ is only assumed to vanish to odd order rather than to exact order one. In the case that
the order of vanishing is at least 3, the point $P_{\nu}$ (which comes from a Heegner construction) is torsion, so Theorem 2 just says that $\mathscr{L}_{p}\left(\nu^{*}\right)=0$.

Corollary 5. Let v be an algebraic Hecke character of infinity type $(1,0)$ satisfying (1-6), (1-7), (1-8) and (1-9) above. Suppose that $\mathscr{L}_{p}\left(\nu^{*}\right) \neq 0$. Then
(1) $L(v, s)$ vanishes to exact order one at the center $s=1$.
(2) $\left(B_{v}(K) \otimes E_{\nu}\right)^{T_{v}}$ is one-dimensional over $E_{\nu}$.
(3) The Shafarevich-Tate group $\amalg\left(B_{v}\right)$ is finite.

Indeed, the nonvanishing of $\mathscr{L}_{p}\left(v^{*}\right)$ implies that a Heegner point on $A_{v}$ is nontorsion, and the conclusion then follows from results of Gross-Zagier and Kolyvagin (see [Kolyvagin 1990; Kolyvagin and Logachëv 1989]). Corollary 5 appears to be new; it would be interesting to see if it can also be obtained via the more indirect methods of [Rubin 1992].

Remark 6. The methods used in the proof of Theorem 2 also give information about the special values $\mathscr{L}_{p}\left(v^{*}\right)$ for Hecke characters $v$ of type $(1+j,-j)$ satisfying (1-6) with $j \geq 0$. A discussion of this point will be taken up in future work. (See [Bertolini et al. 2012a].)

The results of this article concern $p$-adic $L$-functions for the unitary group $\mathrm{U}(1)$; its proofs rely on $p$-adic $L$-functions for the unitary group $\mathrm{U}(2)$, as well as the theorem of Waldspurger relating periods of automorphic forms on $\mathrm{U}(2)$ along an embedded $\mathrm{U}(1)$ to central values of Rankin-Selberg $L$-functions, the latter being the main ingredient in the proof of the main result of [Bertolini et al. 2012b]. One should expect fruitful generalizations of the present work to the setting of higherdimensional unitary groups. The authors are therefore pleased to dedicate this article to the memory of Jon Rogawski, whose deep ideas on automorphic forms, periods and $L$-functions for unitary groups are destined to play a key role in such eventual generalizations.

## 2. Hecke characters and periods

Throughout this article, all number fields that arise are viewed as being embedded in a fixed algebraic closure $\overline{\mathbb{Q}}$ of $\mathbb{Q}$. A complex embedding $\overline{\mathbb{Q}} \rightarrow \mathbb{C}$ and $p$-adic embeddings $\overline{\mathbb{Q}} \rightarrow \mathbb{C}_{p}$ for each rational prime $p$ are also fixed from the outset, so that any finite extension of $\mathbb{Q}$ is simultaneously realized as a subfield of $\mathbb{C}$ and of $\mathbb{C}_{p}$.

2A. Algebraic Hecke characters. We will briefly recall some key definitions regarding algebraic Hecke characters, mainly to fix notation. The reader is referred to [Schappacher 1988, Chapter 0] for more details. Let $K$ and $E$ be number fields.

Given a $\mathbb{Z}$-linear combination

$$
\phi=\sum_{\sigma} n_{\sigma} \sigma \in \mathbb{Z}[\operatorname{Hom}(K, \overline{\mathbb{Q}})]
$$

of embeddings of $K$ into $\overline{\mathbb{Q}}$, we define

$$
\alpha^{\phi}:=\prod_{\sigma}(\sigma \alpha)^{n_{\sigma}},
$$

for all $\alpha \in K^{\times}$. Let $I_{\mathrm{f}}$ denote the group of fractional ideals of $K$ which are prime to a given integral ideal $\mathfrak{f}$ of $K$, and let

$$
J_{\mathfrak{f}}:=\{(\alpha) \mid \alpha \gg 0 \text { and } \alpha-1 \in \mathfrak{f}\} \subseteq I_{\mathfrak{f}} .
$$

Definition 2.1. An E-valued algebraic Hecke character (or simply Hecke character) of $K$ of infinity type $\phi$ and conductor dividing $\mathfrak{f}$ is a homomorphism

$$
\chi: I_{\mathrm{f}} \rightarrow E^{\times}
$$

such that

$$
\begin{equation*}
\chi((\alpha))=\alpha^{\phi}, \quad \text { for all }(\alpha) \in J_{\mathrm{f}} . \tag{2-1}
\end{equation*}
$$

The smallest integral ideal $\mathfrak{g}$ such that $\chi$ can be extended to a Hecke character of conductor dividing $\mathfrak{g}$ is called the conductor of $\chi$, and is denoted $\mathfrak{f}_{\chi}$.

The most basic examples of algebraic Hecke characters are the norm characters on $\mathbb{Q}$ and on $K$ respectively, which are given by

$$
\boldsymbol{N}((a))=|a|, \quad N_{K}:=N \circ \mathrm{~N}_{\mathbb{Q}}^{K} .
$$

Note that the infinity type $\phi$ of a Hecke character $\chi$ must be trivial on all totally positive units congruent to $1 \bmod \mathfrak{f}$. Hence, the existence of such a $\chi$ implies there is an integer $w(\chi)$ (called the weight of $\chi$ or of $\phi$ ) such that for any choice of embedding of $\overline{\mathbb{Q}}$ into $\mathbb{C}$,

$$
n_{\sigma}+n_{\bar{\sigma}}=w(\chi), \quad \text { for all } \sigma \in \operatorname{Hom}(K, \overline{\mathbb{Q}})
$$

Let $U_{\mathfrak{f}} \subset U_{\mathfrak{f}}^{\prime} \subset \mathbb{A}_{K}^{\times}$be the subgroups defined by

$$
U_{\mathfrak{f}}^{\prime}:=\left\{\begin{array}{l|l}
\left(x_{v}\right) \in \mathbb{A}_{K}^{\times} & \begin{array}{l}
x_{v} \equiv 1(\bmod \mathfrak{f}) \\
x_{v}>0
\end{array} \\
x_{v} \text { for all } v \mid \mathfrak{f}, \\
\text { for all real } v
\end{array}\right\}
$$

and

$$
U_{\mathfrak{f}}:=\left\{\left(x_{v}\right) \in U_{\mathfrak{f}}^{\prime} \mid x_{v} \in \mathbb{O}_{K_{v}}^{\times}, \text {for all nonarchimedean } v\right\} .
$$

A Hecke character $\chi$ of conductor dividing $\mathfrak{f}$ may also be viewed as a character on $\mathrm{A}_{K}^{\times} / U_{\mathrm{f}}$ (denoted by the same symbol by a common abuse of notation),

$$
\begin{equation*}
\chi: \mathbb{A}_{K}^{\times} / U_{\mathfrak{f}} \rightarrow E^{\times}, \quad \text { satisfying }\left.\chi\right|_{K^{\times}}=\phi . \tag{2-2}
\end{equation*}
$$

To wit, given $x \in \mathbb{A}_{K}^{\times}$, we define $\chi(x)$ by choosing $\alpha \in K^{\times}$such that $\alpha x$ belongs to $U_{f}^{\prime}$, and setting

$$
\begin{equation*}
\chi(x)=\chi(i(\alpha x)) \phi(\alpha)^{-1} \tag{2-3}
\end{equation*}
$$

where the symbol $i(x)$ denotes the fractional ideal of $K$ associated to $x$. This definition is independent of the choice of $\alpha$ by (2-1). In the opposite direction, given a character $\chi$ as in (2-2), we can set

$$
\chi(\mathfrak{a})=\chi(x) \quad \text { for any } x \in U_{\mathfrak{f}}^{\prime} \text { such that } i(x)=\mathfrak{a} .
$$

The subfield of $E$ generated by the values of $\chi$ on $I_{\mathfrak{f}}$ is easily seen to be independent of the choice of $\mathfrak{f}$ and will be denoted $E_{\chi}$.

Definition 2.2. The central character $\varepsilon_{\eta}$ of a Hecke character $\eta$ of $K$ is the finiteorder character of $\mathbb{Q}$ given by

$$
\left.\eta\right|_{\mathbb{A}_{\mathbb{Q}}^{\times}}=\varepsilon_{\eta} \cdot \boldsymbol{N}^{w(\eta)}
$$

The infinity type $\phi$ defines a homomorphism $\operatorname{Res}_{K / \mathbb{Q}}\left(\boldsymbol{G}_{m}\right) \rightarrow \operatorname{Res}_{E / \mathbb{Q}}\left(\boldsymbol{G}_{m}\right)$ of algebraic groups, and therefore induces a homomorphism

$$
\phi_{\mathbb{A}}: \mathbb{A}_{K}^{\times} \rightarrow \mathbb{A}_{E}^{\times}
$$

on adelic points. Given a Hecke character $\chi$ with values in $E$ and a place $\lambda$ of $E$ (either finite or infinite), we may use $\phi_{\mathbb{A}}$ to define an idèle class character

$$
\chi_{\lambda}: \mathbb{A}_{K}^{\times} / K^{\times} \rightarrow E_{\lambda}^{\times},
$$

by setting

$$
\chi_{\lambda}(x)=\chi(x) / \phi_{A}(x)_{\lambda} .
$$

If $\lambda$ is an infinite place, the character $\chi_{\lambda}$ is a Grossencharacter of $K$ of type $A_{0}$. If $\lambda$ is a finite place, then $\chi_{\lambda}$ factors through $G_{K}^{\text {ab }}$ and gives a Galois character (denoted by $\rho_{\chi, \lambda}$ ) valued in $E_{\lambda}^{\times}$, satisfying

$$
\rho_{\chi, \lambda}\left(\operatorname{Frob}_{\mathfrak{p}}\right)=\chi(\mathfrak{p})
$$

for any prime ideal $\mathfrak{p}$ of $K$ not dividing $\mathfrak{f} \lambda$.
Let $\mathfrak{g}$ be any integral ideal of $K$. The $L$-function and $L$-function with modulus $\mathfrak{g}$ attached to $\chi$ are defined by

$$
L(\chi, s)=\prod_{\mathfrak{p}}\left(1-\frac{\chi(\mathfrak{p})}{\mathrm{Np}^{s}}\right)^{-1}, \quad L_{\mathfrak{g}}(\chi, s)=\prod_{\mathfrak{p} \not \mathfrak{g}}\left(1-\frac{\chi(\mathfrak{p})}{\mathrm{Np}^{s}}\right)^{-1}
$$

Note that $L(\chi, s)=L_{f_{\chi}}(\chi, s)$.
The following definition will only be used in Section 3F.

Definition 2.3. Let $E=\prod_{i} E_{i}$ be a product of number fields. An $E$-valued algebraic Hecke character of conductor dividing $\mathfrak{f}$ is a character

$$
\chi: I_{\mathfrak{f}} \rightarrow E^{\times}
$$

whose projection to each component $E_{i}$ is an algebraic Hecke character in the sense defined above.

2B. Abelian varieties associated to characters of type $(\mathbf{1}, \mathbf{0})$. In this section, we limit the discussion to the case where $K$ is an imaginary quadratic field. Let $\tau: K \mapsto \mathbb{C}$ be the given complex embedding of $K$. A Hecke character of infinity type $\phi=n_{\tau} \tau+n_{\bar{\tau}} \bar{\tau}$ will also be said to be of infinity type ( $n_{\tau}, n_{\bar{\tau}}$ ).

Let $v$ be a Hecke character of $K$ of infinity type $(1,0)$ and conductor $\mathfrak{f}_{v}$, let $E_{\nu} \supset K$ denote the subfield of $\overline{\mathbb{Q}}$ generated by its values, and let $T_{\nu}$ be the ring of integers of $E_{v}$. The Hecke character $v$ gives rise to a compatible system of one-dimensional $\ell$-adic representations of $G_{K}$ with values in $\left(E_{\nu} \otimes \mathbb{Q}_{\ell}\right)^{\times}$, denoted $\rho_{\nu, \ell}$, satisfying

$$
\rho_{v, \ell}\left(\sigma_{\mathfrak{a}}\right)=v(\mathfrak{a}), \quad \text { for all } \mathfrak{a} \in I_{f_{v} \ell},
$$

where $\sigma_{\mathfrak{a}} \in \operatorname{Gal}(\bar{K} / K)$ denotes the Frobenius conjugacy class attached to $\mathfrak{a}$. The theory of complex multiplication realizes the representations $\rho_{\nu, \ell}$ on the division points of CM abelian varieties:

Definition 2.4. A $C M$ abelian variety over $K$ is a pair $(B, E)$ where
(1) $B$ is an abelian variety over $K$;
(2) $E$ is a product of CM fields equipped with the structure of a $K$-algebra and an inclusion

$$
i: E \rightarrow \operatorname{End}_{K}(B) \otimes \mathbb{Q}
$$

satisfying $\operatorname{dim}_{K}(E)=\operatorname{dim} B$;
(3) for all $\lambda \in K \subset E$, the endomorphism $i(\lambda)$ acts on the cotangent space $\Omega^{1}(B / K)$ as multiplication by $\lambda$.

The abelian varieties ( $B, E$ ) over $K$ with complex multiplication by a fixed $E$ form a category denoted $\mathscr{C} \mathcal{M}_{K, E}$ in which a morphism from $B_{1}$ to $B_{2}$ is a morphism $j: B_{1} \rightarrow B_{2}$ of abelian varieties over $K$ for which the diagrams

commute, for all $e \in E$ which belong to both $\operatorname{End}_{K}\left(B_{1}\right)$ and $\operatorname{End}_{K}\left(B_{2}\right)$. An isogeny in $\mathscr{C} \mathcal{M}_{K, E}$ is simply a morphism in this category arising from an isogeny on the underlying abelian varieties.

If $(B, E)$ is a CM abelian variety, its endomorphism ring over $K$ contains a finite-index subring $T^{0}$ of the integral closure $T$ of $\mathbb{Z}$ in $E$. After replacing $B$ by the $K$-isogenous abelian variety $\operatorname{Hom}_{T_{0}}(T, B)$, we can assume that $\operatorname{End}_{K}(B)$ contains $T$. This assumption, which is occasionally convenient, will consistently be made from now on.

Let $(B, E)$ be a CM-abelian variety with $E$ a field, and let $E^{\prime} \supset E$ be a finite extension of $E$ with ring of integers $T^{\prime}$. The abelian variety $B \otimes_{T} T^{\prime}$ is defined to be the variety whose $L$-rational points, for any $L \supset K$, are given by

$$
\left(B \otimes_{T} T^{\prime}\right)(L)=\left(B(\overline{\mathbb{Q}}) \otimes_{T} T^{\prime}\right)^{\operatorname{Gal}(\overline{\mathbb{Q}} / L)} .
$$

This abelian variety is equipped with an action of $T^{\prime}$ by $K$-rational endomorphisms, described by multiplication on the right, and therefore $\left(B \otimes_{T} T^{\prime}, E^{\prime}\right)$ is an object of $\mathscr{C} \mathcal{M}_{K, E^{\prime}}$. Note that $B \otimes_{T} T^{\prime}$ is isogenous to $t:=\operatorname{dim}_{E}\left(E^{\prime}\right)$ copies of $B$, and that the action of $T$ on $B \otimes_{T} T^{\prime}$ agrees with the "diagonal" action of $T$ on $B^{t}$.

Let $\ell$ be a rational prime. For each CM abelian variety $(B, E)$, let

$$
T_{\ell}(B):=\lim _{\leftarrow, n} B\left[\ell^{n}\right](\bar{K}), \quad V_{\ell}(B):=T_{\ell}(B) \otimes_{\mathbb{Z}_{\ell}} \mathbb{Q}_{\ell}
$$

be the $\ell$-adic Tate module and $\ell$-adic representation of $G_{K}$ attached to $B$. The $\mathbb{Q}_{\ell}$-vector space $V_{\ell}(B)$ is a free $E \otimes \mathbb{Q}_{\ell}$-module of rank one via the action of $E$ by endomorphisms. The natural action of $G_{K}:=\operatorname{Gal}(\bar{K} / K)$ on $V_{\ell}(B)$ commutes with this $E \otimes \mathbb{Q}_{\ell}$-action, and the collection $\left\{V_{\ell}(B)\right\}$ thus gives rise to a compatible system of one-dimensional $\ell$-adic representations of $G_{K}$ with values in $\left(E \otimes \mathbb{Q}_{\ell}\right)^{\times}$, denoted $\rho_{B, \ell}$. We note in passing that for any extension $E^{\prime} \supset E$ where $T^{\prime}$ is the integral closure of $T$ in $E^{\prime}$, we have

$$
T_{\ell}\left(B \otimes_{T} T^{\prime}\right)=T_{\ell}(B) \otimes_{T} T^{\prime}, \quad V_{\ell}\left(B \otimes_{T} T^{\prime}\right)=V_{\ell}(B) \otimes_{E} E^{\prime}
$$

The following result is due to Casselman (cf. [Shimura 1971, Theorem 6]).
Theorem 2.5. Let $v$ be a Hecke character of $K$ of type $(1,0)$ as above, and let $\rho_{v, \ell}$ be the associated one-dimensional $\ell$-adic representation with values in $\left(E_{v} \otimes \mathbb{Q}_{\ell}\right)^{\times}$. Then:
(1) There exists a CM abelian variety $\left(B_{v}, E_{v}\right)$ satisfying

$$
\rho_{B_{v}, \ell} \simeq \rho_{v, \ell}
$$

(2) The CM abelian variety $B_{v}$ is unique up to isogeny over $K$. More generally, if $(B, E)$ is any $C M$ abelian variety with $E \supset E_{v}$ satisfying $\rho_{B, \ell} \simeq \rho_{\nu, \ell} \otimes_{E_{v}} E$ as $\left(E \otimes \mathbb{Q}_{\ell}\right)\left[G_{K}\right]$-modules, then there is an isogeny in $\mathscr{C} \mathcal{M}_{K, E}$ from $B$ to $B_{v} \otimes_{T_{v}} T$.

Let $\psi$ be a Hecke character of infinity type ( 1,0 ), and let $\chi$ be a finite-order Hecke character of $K$, so that $\psi \chi^{-1}$ also has infinity type ( 1,0 ). In comparing the abelian varieties $B_{\psi}$ and $B_{\psi \chi^{-1}}$, it is useful to introduce a CM abelian variety $B_{\psi, \chi}$ over $K$, which we now describe.

Let $E_{\chi}$ denote the field generated by $K$ and the values of $\chi$. We denote by $E_{\psi, \chi}$ the compositum of $E_{\psi}$ and $E_{\chi}$, and by $T_{\psi, \chi} \subset E_{\psi, \chi}$ its ring of integers. We also write $H_{\chi}$ for the abelian extension of $K$ which is cut out by $\chi$ viewed as a Galois character of $G_{K}$. Consider first the abelian variety over $K$ with endomorphisms by $T_{\psi, \chi}:$

$$
B_{\psi, \chi}^{0}:=B_{\psi} \otimes_{T_{\psi}} T_{\psi, \chi} .
$$

The natural inclusion $i_{\psi}: T_{\psi} \rightarrow T_{\psi, \chi}$ induces a morphism

$$
\begin{equation*}
i: B_{\psi} \rightarrow B_{\psi, \chi}^{0} \tag{2-4}
\end{equation*}
$$

with finite kernel, which is compatible with the $T_{\psi}$-actions on both sides and is given by

$$
i(P)=P \otimes 1 .
$$

Lemma 2.6. Let $F$ be any number field containing $E_{\psi, \chi}$. With notations as in Equation (1-10) of the Introduction, the restriction map $i^{*}$ induces an isomorphism

$$
\begin{equation*}
i^{*}: \Omega^{1}\left(B_{\psi, \chi}^{0} / F\right)^{T_{\psi, \chi}} \rightarrow \Omega^{1}\left(B_{\psi} / F\right)^{T_{\psi}} \tag{2-5}
\end{equation*}
$$

of one-dimensional $F$-vector spaces.
Proof. The fact that $B_{\psi}$ and $B_{\psi, x}^{0}$ are CM abelian varieties over $F$ implies that the spaces $\Omega^{1}\left(B_{\psi} / F\right)$ and $\Omega^{1}\left(B_{\psi, \chi}^{0} / F\right)$ of regular differentials over $F$ are free of rank one over $T_{\psi} \otimes_{\varrho_{K}} F$ and $T_{\psi, \chi} \otimes_{\varrho_{K}} F$ respectively. In particular, the source and target in (2-5) are both one-dimensional over $F$. The space $\Omega^{1}\left(B_{\psi, \chi}^{0} / F\right)$ is canonically identified with $\operatorname{Hom}_{T_{\psi}}\left(T_{\psi, \chi}, \Omega^{1}\left(B_{\psi} / F\right)\right)$, and under this identification, the pullback

$$
i^{*}: \Omega^{1}\left(B_{\psi, \chi}^{0} / F\right) \rightarrow \Omega^{1}\left(B_{\psi} / F\right)
$$

corresponds to the natural restriction

$$
\operatorname{Hom}_{T_{\psi}}\left(T_{\psi, \chi}, \Omega^{1}\left(B_{\psi} / F\right)\right) \rightarrow \Omega^{1}\left(B_{\psi} / F\right)
$$

sending the function $\varphi$ to $\varphi(1)$. (To see this, consider the map $i_{*}$ on tangent spaces and dualize.) It follows directly from this description that

$$
\operatorname{ker}\left(i^{*}\right) \cap \Omega^{1}\left(B_{\psi, \chi}^{0} / F\right)^{T_{\psi, \chi}}=0 ;
$$

hence the restriction of $i^{*}$ to the one-dimensional $F$-vector space $\Omega^{1}\left(B_{\psi, \chi}^{0} / F\right)^{T \psi, \chi}$ is injective.

Fix $\omega_{\psi} \in \Omega^{1}\left(B_{\psi} / E_{\psi}\right)^{T_{\psi}}$ and define $\omega_{\psi, \chi}^{0} \in \Omega^{1}\left(B_{\psi, \chi}^{0} / \overline{\mathbb{Q}}\right)^{T_{\psi, \chi}}$ by

$$
\begin{equation*}
i^{*}\left(\omega_{\psi, \chi}^{0}\right)=\omega_{\psi} \tag{2-6}
\end{equation*}
$$

It follows from Lemma 2.6 that such an $\omega_{\psi, \chi}^{0}$ exists and is unique (once $\omega_{\psi}$ has been fixed), and further that $\omega_{\psi, \chi}^{0}$ belongs to $\Omega^{1}\left(B_{\psi, \chi}^{0} / E_{\psi, \chi}\right)$.

The character $\chi^{-1}: \operatorname{Gal}\left(H_{\chi} / K\right) \rightarrow T_{\chi}^{\times}$can be viewed as a one-cocycle in

$$
H^{1}\left(\operatorname{Gal}\left(H_{\chi} / K\right), T_{\psi, \chi}^{\times}\right) \subset H^{1}\left(\operatorname{Gal}\left(H_{\chi} / K\right), \operatorname{Aut}_{K}\left(B_{\psi, \chi}^{0}\right)\right) .
$$

Let

$$
\begin{equation*}
B_{\psi, \chi}:=\left(B_{\psi, \chi}^{0}\right)^{x^{-1}} \tag{2-7}
\end{equation*}
$$

denote the twist of $B_{\psi, \chi}^{0}$ by this cocycle. There is a natural identification $B_{\psi, \chi}^{0}(\bar{K})=$ $B_{\psi, \chi}(\bar{K})$ of sets, arising from an isomorphism of varieties over $H_{\chi}$, where $H_{\chi}$ is the extension of $K$ cut out by $\chi$. The actions of $G_{K}$ on $B_{\psi, \chi}^{0}(\bar{K})$ and $B_{\psi, \chi}(\bar{K})$, denoted $*_{0}$ and $*$ respectively, are related by

$$
\begin{equation*}
\sigma * P=\left(\sigma *_{0} P\right) \otimes \chi^{-1}(\sigma), \quad \text { for all } \sigma \in G_{K} . \tag{2-8}
\end{equation*}
$$

In particular, for any $L \supset K$, we have

$$
\begin{equation*}
B_{\psi, \chi}(L)=\left\{P \in B_{\psi}(\overline{\mathbb{Q}}) \otimes_{T_{\psi}} T_{\psi, \chi} \mid \sigma P=P \cdot \chi(\sigma), \text { for all } \sigma \in \operatorname{Gal}(\overline{\mathbb{Q}} / L)\right\} . \tag{2-9}
\end{equation*}
$$

Likewise, the natural actions of $G_{K}$ on $\Omega^{1}\left(B_{\psi, \chi}^{0} / \bar{K}\right)$ and on $\Omega^{1}\left(B_{\psi, \chi} / \bar{K}\right)$ are related by

$$
\begin{equation*}
\sigma * \omega=\left[\chi^{-1}(\sigma)\right]^{*}\left(\sigma *_{0} \omega\right) \quad \text { for all } \sigma \in G_{K} . \tag{2-10}
\end{equation*}
$$

The isomorphism of $B_{\psi, \chi}^{0}$ and $B_{\psi, \chi}$ as CM abelian varieties over $H_{\chi}$ gives natural identifications

$$
\Omega^{1}\left(B_{\psi, \chi}^{0} / H_{\chi}\right)=\Omega^{1}\left(B_{\psi, \chi} / H_{\chi}\right), \quad \Omega^{1}\left(B_{\psi, \chi}^{0} / E_{\psi, \chi}^{\prime}\right)^{T_{\psi, \chi}}=\Omega^{1}\left(B_{\psi, \chi} / E_{\psi, \chi}^{\prime}\right)^{T_{\psi, \chi}},
$$

where $E_{\psi, \chi}^{\prime}$ denotes the subfield of $\overline{\mathbb{Q}}$ generated by $H_{\chi}$ and $E_{\psi, \chi}$.
Let $\omega_{\psi, \chi}$ be any generator of $\Omega^{1}\left(B_{\psi, \chi} / E_{\psi, \chi}\right)^{T_{\psi, x}}$ as an $E_{\psi, \chi}$-vector space. Since $\omega_{\psi, \chi}^{0}$ (defined in (2-6)) and $\omega_{\psi, \chi}$ both generate $\Omega^{1}\left(B_{\psi, \chi} / E_{\psi, \chi}^{\prime}\right)^{T_{\psi, \chi}}$ as an $E_{\psi, \chi^{-}}^{\prime}$ vector space, they necessarily differ by a nonzero scalar in $E_{\psi, \chi}^{\prime}$. To spell out the relation between $\omega_{\psi, \chi}^{0}$ and $\omega_{\psi, \chi}$ more precisely, it will be useful to introduce the notion of a generalized Gauss sum attached to any finite-order character $\chi$ of $G_{K}$. Given such a character, let

$$
E\{\chi\}:=\left\{\lambda \in E_{\chi} H_{\chi} \mid \lambda^{\sigma}=\chi(\sigma) \lambda \text { for all } \sigma \in \operatorname{Gal}\left(E_{\chi} H_{\chi} / E_{\chi}\right)\right\} .
$$

This set is a one-dimensional $E_{\chi}$-vector space in a natural way. It is not closed under multiplication, but

$$
\begin{equation*}
E\left\{\chi_{1}\right\} \cdot E\left\{\chi_{2}\right\}=E\left\{\chi_{1} \chi_{2}\right\}\left(\bmod \left(E_{\chi_{1}} E_{\chi_{2}}\right)^{\times}\right) . \tag{2-11}
\end{equation*}
$$

Definition 2.7. An $E_{\chi}$-vector space generator of $E\{\chi\}$ is called a Gauss sum attached to the character $\chi$, and is denoted $\mathfrak{g}(\chi)$.

By definition, the Gauss sum $\mathfrak{g}(\chi)$ belongs to $E\{\chi\} \cap\left(E_{\chi} H_{\chi}\right)^{\times}$, but is only well-defined up to multiplication by $E_{\chi}^{\times}$. It follows from (2-11) that (2-12)

$$
\mathfrak{g}\left(\chi_{1} \chi_{2}\right)=\mathfrak{g}\left(\chi_{1}\right) \mathfrak{g}\left(\chi_{2}\right)\left(\bmod \left(E_{\chi_{1}} E_{\chi_{2}}\right)^{\times}\right), \quad \mathfrak{g}\left(\chi^{-1}\right)=\mathfrak{g}(\chi)^{-1}\left(\bmod E_{\chi}^{\times}\right) .
$$

The following lemma pins down the relationship between the differentials $\omega_{\psi, \chi}^{0}$ and $\omega_{\psi, \chi}$.

Lemma 2.8. For all Hecke characters $\psi$ and $\chi$ as above,

$$
\omega_{\psi, \chi}=\mathfrak{g}(\chi) \omega_{\psi, \chi}^{0}\left(\bmod E_{\psi, \chi}^{\times}\right) .
$$

Proof. Let $\lambda \in\left(H_{\chi} E_{\psi, \chi}\right)^{\times}$be the scalar satisfying

$$
\begin{equation*}
\omega_{\psi, \chi}=\lambda \omega_{\psi, \chi}^{0} . \tag{2-13}
\end{equation*}
$$

Since $\omega_{\psi, \chi}$ is an $E_{\psi, \chi}$-rational differential on $B_{\psi, \chi}$, for all $\sigma \in \operatorname{Gal}\left(\bar{K} / E_{\psi, \chi}\right)$ we have

$$
\begin{equation*}
\omega_{\psi, \chi}=\sigma * \omega_{\psi, \chi}=\left[\chi^{-1}(\sigma)\right]^{*} \sigma *_{0} \omega_{\psi, \chi}=\chi^{-1}(\sigma) \lambda^{\sigma} \omega_{\psi, \chi}^{0}, \tag{2-14}
\end{equation*}
$$

where the second equality follows from (2-10) and the last from the fact that the differential $\omega_{\psi, \chi}^{0}$ belongs to $\Omega^{1}\left(B_{\psi, \chi}^{0} / E_{\psi, \chi}\right)^{T_{\psi, \chi}}$. Comparing (2-13) and (2-14) gives $\lambda^{\sigma}=\chi(\sigma) \lambda$, and hence $\lambda=\mathfrak{g}(\chi)\left(\bmod E_{\psi, \chi}^{\times}\right)$.

The following lemma relates the abelian varieties $B_{\psi, \chi}$ and $B_{v}$, where $\nu=\psi \chi^{-1}$.
Lemma 2.9. There is an isogeny defined over $K$ :

$$
i_{v}: B_{\psi, \chi} \rightarrow B_{v} \otimes_{T_{v}} T_{\psi, \chi},
$$

which is compatible with the action of $T_{\psi, \chi}$ by endomorphisms on both sides.
Proof. The pair ( $B_{\psi, \chi}^{0}, E_{\psi, \chi}$ ) is a CM abelian variety having $\psi$ (viewed as taking values in $E_{\psi, \chi}$ ) as its associated Hecke character. The Hecke character attached to the Galois twist $B_{\psi, \chi}$ is therefore $\psi \chi^{-1}=v$. The second part of Theorem 2.5 implies that $B_{\psi, \chi}$ and $B_{\nu} \otimes_{T_{v}} T_{\psi, \chi}$ are isogenous over $K$ as CM abelian varieties.

2C. Complex periods and special values of L-functions. This section recalls certain periods attached to the quadratic imaginary field $K$ and to Hecke characters of this field. We begin by fixing:
(1) An elliptic curve $A$ with complex multiplication by $\mathbb{O}_{K}$, defined over a finite extension $F$ of $K$. (Note that $F$ necessarily contains the Hilbert class field of $K$.)
(2) A regular differential $\omega_{A} \in \Omega^{1}(A / F)$.
(3) A nonzero element $\gamma$ of $H_{1}(A(\mathbb{C}), \mathbb{Q})$.

The complex period attached to this data is defined by

$$
\begin{equation*}
\Omega(A):=\frac{1}{2 \pi i} \int_{\gamma} \omega_{A} \tag{2-15}
\end{equation*}
$$

Note that $\Omega(A)$ depends on the pair $\left(\omega_{A}, \gamma\right)$. A different choice of $\omega_{A}$ or $\gamma$ has the effect of multiplying $\Omega(A)$ by a scalar in $F^{\times}$, and therefore $\Omega(A)$ can be viewed as a well-defined element of $\mathbb{C}^{\times} / F^{\times}$.

For any Hecke character $\psi$ of $K$, recall that $\psi^{*}$ is the Hecke character defined as in the Introduction by $\psi^{*}(x)=\psi(\bar{x})$. Suppose that $\psi$ is of infinity type $(1,0)$, and as before let $E_{\psi} \subset \overline{\mathbb{Q}} \subset \mathbb{C}$ denote the field generated by the values of $\psi$ (or, equivalently, $\psi^{*}$ ). Choose (arbitrary) nonzero elements

$$
\omega_{\psi} \in \Omega^{1}\left(B_{\psi} / E_{\psi}\right)^{T_{\psi}}, \quad \gamma \in H_{1}\left(B_{\psi}(\mathbb{C}), \mathbb{Q}\right)
$$

with $B_{\psi}$ the CM abelian variety attached to $\psi$ by Theorem 2.5 , and $\Omega^{1}\left(B_{\psi} / E_{\psi}\right)^{T_{\psi}}$ defined in (1-10). The period $\Omega\left(\psi^{*}\right)$ attached to $\psi^{*}$ is defined by setting

$$
\Omega\left(\psi^{*}\right)=\frac{1}{2 \pi i} \int_{\gamma} \omega_{\psi} \quad\left(\bmod E_{\psi}^{\times}\right)
$$

Note that the complex number $\Omega\left(\psi^{*}\right)$ does not depend, up to multiplication by $E_{\psi}^{\times}$, on the choices of $\omega_{\psi}$ and $\gamma$ that were made to define it.

Lemma 2.10. If $\psi$ is a Hecke character of infinity type (1, 0), and $\chi$ is a finite-order character, then

$$
\begin{equation*}
\Omega\left(\psi^{*} \chi\right)=\Omega\left(\psi^{*}\right) \mathfrak{g}\left(\chi^{*}\right)^{-1}\left(\bmod E_{\psi, \chi}^{\times}\right) \tag{2-16}
\end{equation*}
$$

Proof. Choose a nonzero generator $\gamma$ of $H_{1}\left(B_{\psi, \chi}^{0}(\mathbb{C}), \mathbb{Q}\right)=H_{1}\left(B_{\psi, \chi}(\mathbb{C}), \mathbb{Q}\right)$ (viewed as a one-dimensional $E_{\psi, \chi}$ vector space via the endomorphism action). By definition,

$$
\Omega\left(\left(\psi \chi^{-1}\right)^{*}\right)=\int_{\gamma} \omega_{\psi, \chi}=\mathfrak{g}(\chi) \int_{\gamma} \omega_{\psi, \chi}^{0}=\mathfrak{g}(\chi) \Omega\left(\psi^{*}\right)\left(\bmod E_{\psi, \chi}^{\times}\right)
$$

where the second equality follows from Lemma 2.8. The result now follows after substituting $\chi^{*-1}$ for $\chi$.

As in [Schappacher 1988, §1.8], one can also attach a period $\Omega(\psi)$ to an arbitrary Hecke character $\psi$ of $K$; these satisfy the following relations:

Proposition 2.11. Let $\psi$ be a Hecke character of infinity type ( $k, j$ ). Then
(1) The ratio

$$
\frac{\Omega\left(\psi^{*}\right)}{(2 \pi i)^{j} \Omega(A)^{k-j}}
$$

is algebraic.
(2) For all $\psi$ and $\psi^{\prime}$,

$$
\Omega\left(\psi \psi^{\prime}\right)=\Omega(\psi) \Omega\left(\psi^{\prime}\right)\left(\bmod E_{\psi, \psi^{\prime}}^{\times}\right),
$$

where $E_{\psi, \psi^{\prime}}$ is the subfield of $\overline{\mathbb{Q}}$ generated by $E_{\psi}$ and $E_{\psi^{\prime}}$.
The following theorem is due to Goldstein and Schappacher [1981] in certain cases and Blasius [1986] in the general case (even CM fields).
Theorem 2.12. Suppose that $\psi$ has infinity type $(k, j)$ with $k>j$, and that $m$ is a critical integer for $L\left(\psi^{-1}, s\right)$. Then

$$
\frac{L\left(\psi^{-1}, m\right)}{(2 \pi i)^{m} \Omega\left(\psi^{*}\right)} \text { belongs to } E_{\psi},
$$

and for all $\tau \in \operatorname{Gal}\left(E_{\psi} / K\right)$,

$$
\left(\frac{L\left(\psi^{-1}, m\right)}{(2 \pi i)^{m} \Omega\left(\psi^{*}\right)}\right)^{\tau}=\frac{L\left(\left(\psi^{-1}\right)^{\tau}, m\right)}{(2 \pi i)^{m} \Omega\left(\left(\psi^{*}\right)^{\tau}\right)} .
$$

Here the action of $\tau \in \operatorname{Gal}\left(E_{\psi} / K\right)$ on $\psi^{-1}$ and $\psi^{*}$ is defined by viewing them as $E_{\chi}$-valued characters on ideals and applying $\tau$ to their values.

2D. p-adic periods. Fix a prime $p$ that splits in $K$. We will need $p$-adic analogs of the periods $\Omega(A)$ and $\Omega\left(\nu^{*}\right)$. The $p$-adic analog $\Omega_{p}(A)$ of $\Omega(A)$ is obtained by considering the base change $A_{\mathbb{C}_{p}}$ of $A$ to $\mathbb{C}_{p}$ (via our fixed embedding of $F$ into $\mathbb{C}_{p}$ ). Assume that $A$ has good reduction at the maximal ideal of $\mathbb{O}_{C_{p}}$, that is, that $A_{\mathbb{C}_{p}}$ extends to a smooth proper model $A_{\mathscr{C}_{p}}$ over $\mathbb{O}_{\mathbb{C}_{p}}$. The $p$-adic completion $\hat{A}_{\mathbb{C}_{C_{p}}}$ of $A$ along its special fiber is isomorphic to $\hat{\boldsymbol{G}}_{m}$. Following [de Shalit 1987, II, §4.4], choose an isomorphism $\iota_{p}: \hat{A} \rightarrow \hat{\boldsymbol{G}}_{m}$ over $\mathbb{O}_{\mathbb{C}_{p}}$, and define $\Omega_{p}(A) \in \mathbb{C}_{p}^{\times}$ by the rule

$$
\begin{equation*}
\omega_{A}=\Omega_{p}(A) \cdot \iota_{p}^{*}(d u / u), \tag{2-17}
\end{equation*}
$$

where $u$ is the standard coordinate on $\hat{\boldsymbol{G}}_{m}$. The invariant $\Omega_{p}(A) \in \mathbb{C}_{p}^{\times}$thus defined depends on the choices of $\omega_{A}$ and $\iota_{p}$, but only up to multiplication by a scalar in $F^{\times}$. Observe also that $\Omega(A)$ and $\Omega_{p}(A)$ each depend linearly in the same way on the choice of the global differential $\omega_{A}$.

The $p$-adic period $\Omega_{p}(A)$ can be used to define $p$-adic analogs of the complex periods that appear in the statement of Theorem 2.12.
Definition 2.13. Let $v$ be a Hecke character of $K$ of type ( 1,0 ). The $p$-adic period $\Omega_{p}\left(v^{*}\right)$ is defined by

$$
\Omega_{p}\left(v^{*}\right):=\Omega_{p}(A) \cdot \frac{\Omega\left(v^{*}\right)}{\Omega(A)} .
$$

More generally, for any character $v$ of infinity type $(k, j)$, we define

$$
\Omega_{p}\left(\nu^{*}\right):=\Omega_{p}(A)^{k-j} \cdot \frac{\Omega\left(\nu^{*}\right)}{(2 \pi i)^{j} \Omega(A)^{k-j}} .
$$

It can be seen from this definition that the period $\Omega_{p}\left(\nu^{*}\right)$, like its complex counterpart $\Omega\left(\nu^{*}\right)$, is well-defined up to multiplication by a scalar in $E_{v}^{\times}$. The following $p$-adic analog of Lemma 2.10 is a direct consequence of this lemma combined with the definition of $\Omega_{p}(\psi)$ :

Lemma 2.14. If $\psi$ is a Hecke character of infinity type $(1,0)$, and $\chi$ is a finite order character, then

$$
\begin{equation*}
\Omega_{p}\left(\psi^{*} \chi\right)=\Omega_{p}\left(\psi^{*}\right) \mathfrak{g}\left(\chi^{*}\right)^{-1}\left(\bmod E_{\psi, \chi}^{\times}\right) . \tag{2-18}
\end{equation*}
$$

Likewise, Proposition 2.11 implies:
Proposition 2.15. Let $\psi$ be a Hecke character of infinity type ( $k, j$ ). Then:
(1) The ratio

$$
\frac{\Omega_{p}\left(\psi^{*}\right)}{(2 \pi i)^{j} \Omega_{p}(A)^{k-j}}
$$

is algebraic.
(2) For all $\psi$ and $\psi^{\prime}$,

$$
\begin{equation*}
\Omega_{p}\left(\psi \psi^{\prime}\right)=\Omega_{p}(\psi) \Omega_{p}\left(\psi^{\prime}\right)\left(\bmod E_{\psi, \psi^{\prime}}^{\times}\right) . \tag{2-19}
\end{equation*}
$$

## 3. $\boldsymbol{p}$-adic $\boldsymbol{L}$-functions and Rubin's formula

3A. The Katz p-adic L-function. Throughout this chapter, we will fix a prime $p$ that is split in $K$. Let $\mathfrak{c}$ be an integral ideal of $K$ which is prime to $p$, and let $\Sigma(\mathfrak{c})$ denote the set of all Hecke characters of $K$ of conductor dividing $\mathfrak{c}$. Denote by $\mathfrak{p}$ the prime above $p$ corresponding to the chosen embedding $K \hookrightarrow \overline{\mathbb{Q}}_{p}$.

A character $v \in \Sigma(\mathfrak{c})$ is called a critical character if $L\left(v^{-1}, 0\right)$ is a critical value in the sense of Deligne, that is, if the $\Gamma$-factors that arise in the functional equation for $L\left(\nu^{-1}, s\right)$ are nonvanishing and have no poles at $s=0$. The set $\Sigma_{\text {crit }}(\mathfrak{c})$ of critical characters can be expressed as the disjoint union

$$
\Sigma_{\text {crit }}(\mathfrak{c})=\Sigma_{\text {crit }}^{(1)}(\mathfrak{c}) \cup \Sigma_{\text {crit }}^{(2)}(\mathfrak{c}),
$$



Figure 1. Critical infinity types for the Katz $p$-adic $L$-function.
where

$$
\begin{aligned}
& \Sigma_{\text {crit }}^{(1)}(\mathfrak{c})=\left\{v \in \Sigma(\mathfrak{c}) \text { of type }\left(\ell_{1}, \ell_{2}\right) \text { with } \ell_{1} \leq 0, \ell_{2} \geq 1\right\}, \\
& \Sigma_{\text {crit }}^{(2)}(\mathfrak{c})=\left\{v \in \Sigma(\mathfrak{c}) \text { of type }\left(\ell_{1}, \ell_{2}\right) \text { with } \ell_{1} \geq 1, \ell_{2} \leq 0\right\} .
\end{aligned}
$$

The possible infinity types of Hecke characters in $\Sigma_{\text {crit }}^{(2)}(\mathfrak{c})$ lie in the shaded region in Figure 1 whose reflection about the principal diagonal corresponds likewise to $\Sigma_{\text {crit }}^{(1)}(\mathfrak{c})$. Note in particular that when $\mathfrak{c}=\overline{\mathfrak{c}}$, the regions $\Sigma_{\text {crit }}^{(1)}(\mathfrak{c})$ and $\Sigma_{\text {crit }}^{(2)}(\mathfrak{c})$ are interchanged by the involution $v \mapsto \nu^{*}$. The set $\Sigma_{\text {crit }}(\mathfrak{c})$ is endowed with a natural $p$-adic topology as described in Section 5.2 of [Bertolini et al. 2012b]. The subsets $\Sigma_{\text {crit }}^{(1)}(\mathfrak{c})$ and $\Sigma_{\text {crit }}^{(2)}(\mathfrak{c})$ are each dense in the completion $\hat{\Sigma}_{\text {crit }}(\mathfrak{c})$ relative to this topology.

Recall that $\mathfrak{p}$ is the prime above $p$ induced by our chosen embedding of $K$ into $\mathbb{C}_{p}$. The following theorem on the existence of the $p$-adic $L$-function is due to Katz. The statement below is a restatement of [de Shalit 1987, II, Theorem 4.14] with a minor correction, and restricted to characters unramified at $p$. We remark that since our characters are unramified at $p$, the Gauss sum in the interpolation formula in [loc. cit.] is equal to 1 .

Theorem 3.1. There exists a p-adic analytic function $v \mapsto \mathscr{L}_{p, \mathfrak{c}}(v)$ (valued in $\mathbb{C}_{p}$ ) on $\hat{\Sigma}_{\text {crit }}(\mathfrak{c})$ which is determined by the interpolation property

$$
\begin{equation*}
\frac{\mathscr{L}_{p, \mathfrak{c}}(v)}{\Omega_{p}(A)^{\ell_{1}-\ell_{2}}}=\left(\frac{\sqrt{D}}{2 \pi}\right)^{\ell_{2}}\left(\ell_{1}-1\right)!(1-v(\mathfrak{p}) / p)\left(1-v^{-1}(\overline{\mathfrak{p}})\right) \frac{L_{\mathfrak{c}}\left(v^{-1}, 0\right)}{\Omega(A)^{\ell_{1}-\ell_{2}}}, \tag{3-1}
\end{equation*}
$$

for all critical characters $v \in \Sigma_{\text {crit }}^{(2)}(\mathfrak{c})$ of infinity type $\left(\ell_{1}, \ell_{2}\right)$.
The right-hand side of (3-1) belongs to $\overline{\mathbb{Q}}$, by part (1) of Proposition 2.11 and Theorem 2.12 with $m=0$. Equation (3-1) should be interpreted to mean that the left-hand side also belongs to $\overline{\mathbb{Q}}$, viewed as a subfield of $\mathbb{C}_{p}$ under the chosen embeddings, and agrees with the right-hand side. Note that although both sides of (3-1) depend on the choice of the differential $\omega_{A}$ that was made in the definition of the periods $\Omega(A)$ and $\Omega_{p}(A)$, the quantity $\mathscr{L}_{p, \mathfrak{c}}(\nu)$, just like its complex counterpart $L_{\mathfrak{c}}\left(v^{-1}, 0\right)$, does not depend on this choice.

Remark 3.2. Once a choice of $\mathfrak{c}$ is fixed, we shall often drop the subscript $\mathfrak{c}$ and simply write $\mathscr{L}_{p}$ for the $p$-adic $L$-function.

The following corollary is the $p$-adic counterpart of Theorem 2.12.
Corollary 3.3. Suppose that $v \in \Sigma_{\text {crit }}^{(2)}(\mathfrak{c})$. Then

$$
\frac{\mathscr{L}_{p, \mathrm{c}}(\nu)}{\Omega_{p}\left(\nu^{*}\right)} \text { belongs to } E_{\nu} .
$$

Proof. Suppose that $v$ has infinity type $\left(\ell_{1}, \ell_{2}\right)$. By the definition of $\Omega_{p}\left(v^{*}\right)$ and the interpolation property of the Katz $p$-adic $L$-function in Theorem 3.1, we have

$$
\begin{aligned}
\frac{\mathscr{L}_{p, \mathfrak{c}}(v)}{\Omega_{p}\left(v^{*}\right)} & =\frac{\mathscr{L}_{p, \mathfrak{c}}(v)}{\Omega_{p}(A)^{\ell_{1}-\ell_{2}}} \times \frac{(2 \pi i)^{\ell_{2}} \Omega(A)^{\ell_{1}-\ell_{2}}}{\Omega\left(v^{*}\right)} \\
& =\sqrt{-D}^{\ell_{2}} \cdot\left(\ell_{1}-1\right)!(1-v(\mathfrak{p}))\left(1-v^{-1}(\overline{\mathfrak{p}})\right) \frac{L_{\mathfrak{c}}\left(v^{-1}, 0\right)}{\Omega\left(v^{*}\right)} .
\end{aligned}
$$

The result is now a direct consequence of Theorem 2.12 with $m=0$.
Corollary 3.3 expresses $\mathscr{L}_{p, \mathfrak{c}}(\nu)$ as an $E_{\nu}$-multiple of a $p$-adic period $\Omega_{p}\left(\nu^{*}\right)$, when $v$ lies in the range $\Sigma_{\text {crit }}^{(2)}(\mathfrak{c})$ of classical interpolation for the Katz $p$-adic $L$ function. On the other hand, the characters in $\Sigma_{\text {crit }}^{(1)}(\mathfrak{c})$ are outside the range of interpolation, and so Corollary 3.3 does not directly say anything about these values, and indeed the main goal of this paper is to obtain analogous results for certain characters in $\Sigma_{\text {crit }}^{(1)}(\mathfrak{c})$. It turns out that the methods of this paper only allow us to study $\mathscr{L}_{p, \mathfrak{c}}(v)$ for characters $v$ in $\Sigma_{\text {crit }}^{(2)}(\mathfrak{c})$ satisfying the following auxiliary (but not unnatural) condition:

$$
\begin{equation*}
v \text { is a self-dual Hecke character with } \varepsilon_{v}=\varepsilon_{K} \text {. } \tag{3-2}
\end{equation*}
$$

For the benefit of the reader, we now recall this key definition.

Definition 3.4. A Hecke character $v \in \Sigma_{\text {crit }}(\mathfrak{c})$ is said to be self-dual or anticyclotomic if

$$
\nu v^{*}=N_{K} .
$$

The reason for the terminology in Definition 3.4 is that the functional equation for the $L$-series $L\left(v^{-1}, s\right)$ relates $L\left(v^{-1}, s\right)$ to $L\left(v^{-1},-s\right)$, and therefore $s=0$ is the central critical point for this complex $L$-series. Note that a self-dual character is necessarily of infinity type $(1+j,-j)$ for some $j \in \mathbb{Z}$. Also the conductor of a self-dual character is clearly invariant under complex conjugation. If $\mathfrak{c}$ is an integral ideal such that $\mathfrak{c}=\overline{\mathfrak{c}}$, we denote by $\Sigma_{\text {sd }}(\mathfrak{c})$ the set of self-dual Hecke characters of conductor exactly $\mathfrak{c}$, and write

$$
\Sigma_{\text {sd }}^{(1)}(\mathfrak{c})=\Sigma_{\text {crit }}^{(1)}(\mathfrak{c}) \cap \Sigma_{\text {sd }}(\mathfrak{c}), \quad \Sigma_{\text {sd }}^{(2)}(\mathfrak{c})=\Sigma_{\text {crit }}^{(2)}(\mathfrak{c}) \cap \Sigma_{\text {sd }}(\mathfrak{c}) .
$$

In particular, the possible infinity types of characters in $\Sigma_{\text {sd }}^{(2)}(\mathfrak{c})$ correspond to the black dots in Figure 1.

For convenience, we restate Theorem 3.1 for self-dual characters.
Proposition 3.5. For all characters $v \in \Sigma_{\text {sd }}^{(2)}(\mathfrak{c})$ of infinity type $(1+j,-j)$ with $j \geq 0$,

$$
\begin{equation*}
\frac{\mathscr{L}_{p, \mathfrak{c}}(v)}{\Omega_{p}(A)^{1+2 j}}=\left(1-v^{-1}(\overline{\mathfrak{p}})\right)^{2} \times \frac{j!(2 \pi)^{j} L_{\mathfrak{c}}\left(v^{-1}, 0\right)}{\sqrt{D}^{j} \Omega(A)^{1+2 j}} . \tag{3-3}
\end{equation*}
$$

Remark 3.6. In the proposition above, we could equally write $L\left(v^{-1}, 0\right)$ instead of $L_{\mathfrak{c}}\left(v^{-1}, 0\right)$, since $v$ has conductor exactly equal to $c$.

Remark 3.7. The central character of such a $v$ is very restricted. Indeed, for any Hecke character $\nu$, it is clear that $\varepsilon_{\bar{v}}=\overline{\varepsilon_{v}}$, while $\varepsilon_{\nu^{*}}=\varepsilon_{v}$. Further, if $\nu$ is a self-dual character, it follows that for any $x \in \mathbb{A}_{K}^{\times}$,

$$
\nu\left(\mathrm{N}_{\mathbb{Q}}^{K}(x)\right)=\nu(x \bar{x})=\left(\nu \nu^{*}\right)(x)=\boldsymbol{N}_{K}(x)=\boldsymbol{N}\left(\mathrm{N}_{\mathbb{Q}}^{K}(x)\right) .
$$

Hence

$$
\left.\nu\right|_{N_{\mathbb{Q}}^{K} A_{K}^{X}}=N \quad \text { and }\left.\quad \varepsilon_{v}\right|_{N_{\mathbb{Q}}^{K} A_{K}^{x}}=1 .
$$

This implies that the central character $\varepsilon_{\nu}$ of a self-dual character $v$ is either 1 or $\varepsilon_{K}$, where $\varepsilon_{K}$ denotes the quadratic Dirichlet character corresponding to the extension $K / \mathbb{Q}$. Conversely, it is easy to see that if $v$ is a Hecke character with $w(\nu)=1$ and $\varepsilon_{v}=1$ or $\varepsilon_{K}$, then $v$ is a self-dual character.

We define:

$$
\begin{equation*}
\Sigma_{\mathrm{sd}}(\mathfrak{c})^{+}:=\left\{v \in \Sigma_{\mathrm{sd}}(\mathfrak{c}) ; \quad \varepsilon_{v}=1\right\}, \quad \Sigma_{\mathrm{sd}}(\mathfrak{c})^{-}:=\left\{v \in \Sigma_{\mathrm{sd}}(\mathfrak{c}) ; \quad \varepsilon_{v}=\varepsilon_{K}\right\} . \tag{3-4}
\end{equation*}
$$

The sets $\Sigma_{\text {sd }}^{(1)}(\mathfrak{c})^{ \pm}$and $\Sigma_{\text {sd }}^{(2)}(\mathfrak{c})^{ \pm}$are defined similarly.

Our approach to studying $\mathscr{L}_{p, \mathrm{c}}(v)$ for characters $v$ satisfying (3-2), that is, those $\nu$ lying in $\Sigma_{\text {sd }}^{(1)}(\mathfrak{c})^{-}$for some $\mathfrak{c}$, relies on a different kind of $p$-adic $L$-function. This latter $p$-adic $L$-function is attached to Rankin-Selberg $L$-series and is recalled in the following section.

3B. p-adic Rankin L-series. In this section, we consider $p$-adic $L$-functions obtained by interpolating special values of Rankin-Selberg $L$-series associated to modular forms and Hecke characters of a quadratic imaginary field $K$ of odd discriminant. We briefly recall the definition of this $p$-adic $L$-function that is given in Section 5 of [Bertolini et al. 2012b], referring the reader to that work for a more detailed description.

Let $S_{k}\left(\Gamma_{0}(N), \varepsilon\right)$ denote the space of cusp forms of weight $k \geq 2$ and character $\varepsilon$ on $\Gamma_{0}(N)$. Let $f \in S_{k}\left(\Gamma_{0}(N), \varepsilon\right)$ be a normalized newform and let $E_{f}$ denote the subfield of $\mathbb{C}$ generated by its Fourier coefficients.

Definition 3.8. The pair $(f, K)$ is said to satisfy the Heegner hypothesis if $0_{K}$ contains a cyclic ideal of norm $N$, that is, an integral ideal $\mathfrak{N}$ of $\mathbb{O}_{K}$ with $\mathcal{O}_{K} / \mathfrak{N}=$ $\mathbb{Z} / N \mathbb{Z}$.

Assume from now on that $(f, K)$ satisfies the Heegner hypothesis, and let $\mathfrak{N}$ be a cyclic $\widehat{O}_{K}$-ideal of norm $N$. We write $\mathfrak{N}_{\varepsilon}$ for the unique ideal dividing $\mathfrak{N}$ of norm $N_{\varepsilon}$, where $N_{\varepsilon}$ is the conductor of $\varepsilon$.
Definition 3.9. A Hecke character $\chi$ of $K$ of infinity type $\left(\ell_{1}, \ell_{2}\right)$ is said to be central critical for $f$ if

$$
\ell_{1}+\ell_{2}=k \quad \text { and } \quad \varepsilon_{\chi}=\varepsilon .
$$

The reason for the terminology of Definition 3.9 is that when $\chi$ satisfies these hypotheses, the complex Rankin $L$-series $L\left(f, \chi^{-1}, s\right)$ is self-dual and $s=0$ is its central (critical) point.

Definition 3.10. Let $c$ be a rational integer prime to $p N$. Then $\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$ is defined to be the set of Hecke characters $\chi$ of $K$ such that
(1) $\chi$ is central critical for $f$;
(2) $\mathfrak{f}_{\chi}=c \cdot \mathfrak{N}_{\varepsilon}$;
(3) the local $\operatorname{sign} \varepsilon_{q}\left(f, \chi^{-1}\right)=+1$ for all finite primes $q$.

It is easily checked that this agrees with the definition of $\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$ given in [Bertolini et al. 2012b, §4.1], where this is just denoted $\Sigma_{\text {cc }}(\mathfrak{N})$. Further, as in [loc. cit.], given conditions (1) and (2) above, condition (3) is automatic except possibly for primes $q$ lying in the set $S_{f}$ defined by

$$
S_{f}:=\left\{q: q \mid(N, D), q \nmid N_{\varepsilon}\right\} .
$$



Figure 2. Critical infinity types for the $p$-adic Rankin $L$-function.

The set $\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$ can be expressed as a disjoint union

$$
\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)=\Sigma_{\mathrm{cc}}^{(1)}(c, \mathfrak{N}, \varepsilon) \cup \Sigma_{\mathrm{cc}}^{(2)}(c, \mathfrak{N}, \varepsilon)
$$

where $\Sigma_{\mathrm{cc}}^{(1)}(c, \mathfrak{N}, \varepsilon)$ and $\Sigma_{\mathrm{cc}}^{(2)}(c, \mathfrak{N}, \varepsilon)$ denote the subsets consisting of characters of infinity type $(k+j,-j)$, with $1-k \leq j \leq-1$ and $j \geq 0$ respectively. We shall also denote by $\hat{\Sigma}_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$ the completion of $\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$ relative to the $p$-adic compact open topology on $\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$, which is defined in Section 5.2 of [Bertolini et al. 2012b]. The infinity types of Hecke characters in $\Sigma_{c c}^{(1)}(c, \mathfrak{N}, \varepsilon)$ and $\Sigma_{c c}^{(2)}(c, \mathfrak{N}, \varepsilon)$ correspond respectively to the white and black dots in the shaded regions in Figure 2. We note that the set $\Sigma_{\text {cc }}^{(2)}(c, \mathfrak{N}, \varepsilon)$ of classical central critical characters "of type 2 " is dense in $\hat{\Sigma}_{\text {cc }}(c, \mathfrak{N}, \varepsilon)$.

For all $\chi \in \Sigma_{\text {cc }}^{(2)}(c, \mathfrak{N}, \varepsilon)$ of infinity type $(k+j,-j)$ with $j \geq 0$, let $E_{f, \chi}$ denote the subfield of $\mathbb{C}$ generated by $E_{f}$ and the values of $\chi$, and let $E_{f, \chi, \varepsilon}$ be the field generated by $E_{f, \chi}$ and by the abelian extension of $\mathbb{Q}$ cut out by $\varepsilon$. The algebraic part of $L\left(f, \chi^{-1}, 0\right)$ is defined by the rule

$$
\begin{equation*}
L_{\mathrm{alg}}\left(f, \chi^{-1}, 0\right):=w(f, \chi)^{-1} C(f, \chi, c) \cdot \frac{L\left(f, \chi^{-1}, 0\right)}{\Omega(A)^{2(k+2 j)}} \tag{3-5}
\end{equation*}
$$

where $w(f, \chi)^{-1} \in E_{f, \chi, \varepsilon}$ and $C(f, \chi, c)$ are respectively the scalar (of complex norm 1) and the explicit real constant defined in [Bertolini et al. 2012b, (5.1.11) and Theorem 4.6]; we have

$$
\begin{equation*}
C(f, \chi, c)=\frac{2^{k+2 j-2} \pi^{k+2 j-1} j!(k+j-1)!w_{K}}{\sqrt{D}} \prod_{q \mid c} \frac{q-\varepsilon_{K}(q)}{q-1}, \tag{3-6}
\end{equation*}
$$

where $w_{K}=\# \mathbb{O}_{K}^{\times}$is the number of roots of unity in $K$. Theorems 5.5 and 5.10 of [loc. cit.] show respectively that the values $L_{\text {alg }}\left(f, \chi^{-1}, 0\right)$ belong to $\overline{\mathbb{Q}}$, and that they interpolate $p$-adically:
Proposition 3.11. Let $\chi \mapsto L_{p}(f, \chi)$ be the function on $\Sigma_{\mathrm{cc}}^{(2)}(c, \mathfrak{N}, \varepsilon)$ defined by (3-7)
$L_{p}(f, \chi):=\Omega_{p}(A)^{2(k+2 j)}\left(1-\chi^{-1}(\overline{\mathfrak{p}}) a_{p}(f)+\chi^{-2}(\overline{\mathfrak{p}}) \varepsilon(p) p^{k-1}\right)^{2} L_{\text {alg }}\left(f, \chi^{-1}, 0\right)$,
for $\chi$ of infinity type $(k+j,-j)$ with $j \geq 0$. This function extends (uniquely) to a p-adically continuous function on $\hat{\Sigma}_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$.

The function $\chi \mapsto L_{p}(f, \chi)$ on $\hat{\Sigma}_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$ will be referred to as the p-adic Rankin L-function attached to the cusp form $f$.

3C. A p-adic Gross-Zagier formula. In this section, we specialize to the case where the newform $f$ is of weight $k=2$, and assume that $\chi$ is a finite-order Hecke character of $K$ satisfying

$$
\chi N_{K} \text { belongs to } \Sigma_{\mathrm{cc}}^{(1)}(c, \mathfrak{N}, \varepsilon)
$$

In particular, the character $\chi \boldsymbol{N}_{K}$ lies outside the domain $\Sigma_{\mathrm{cc}}^{(2)}(c, \mathfrak{N}, \varepsilon)$ of classical interpolation defining $L_{p}(f,-)$. The $p$-adic Gross-Zagier formula alluded to in the title of this section relates the special value $L_{p}\left(f, \chi N_{K}\right)$ to the formal group logarithm of a Heegner point on the modular abelian variety attached to $f$.

The Eichler-Shimura construction associates to $f$ an abelian variety $B_{f}$ with endomorphism by an order in the ring of integers $T_{f} \subset E_{f}$, and a surjective morphism

$$
\Phi_{f}: J_{1}(N) \rightarrow B_{f}
$$

of abelian varieties over $\mathbb{Q}$, called the modular parametrization, which is welldefined up to a rational isogeny. Let

$$
\omega_{f}=2 \pi i f(\tau) d \tau \in \Omega^{1}\left(X_{1}(N) / E_{f}\right)
$$

be the differential form on $X_{1}(N)$ attached to $f$; we use the same symbol $\omega_{f}$ to denote the associated one-form on $J_{1}(N)$. Let $\omega_{B_{f}} \in \Omega^{1}\left(B_{f} / E_{f}\right)^{T_{f}}$ be the unique one-form satisfying

$$
\begin{equation*}
\Phi_{f}^{*}\left(\omega_{B_{f}}\right)=\omega_{f} . \tag{3-8}
\end{equation*}
$$

Let $A^{\prime}$ be an elliptic curve with endomorphism ring isomorphic to the order $\mathfrak{O}_{c}=\mathbb{Z}+c 0_{K}$ of conductor $c$, defined over the ring class field $H_{c}$ of conductor $c$. The pair $\left(A^{\prime}, A^{\prime}[\mathfrak{N}]\right)$ corresponds to a point on $X_{0}(N)\left(H_{c}\right)$. Let $t$ be any generator of $A^{\prime}[\mathfrak{N}]$. Then the triple ( $A^{\prime}, A^{\prime}[\mathfrak{N}], t$ ) corresponds to a point in $X_{1}(N)$, whose field of definition $H_{c, \mathfrak{N}}$ is an abelian extension of $K$, independent of the choice of $t$, and the finite-order Hecke character $\chi$ can be viewed as a character

$$
\chi: \operatorname{Gal}\left(H_{c, \mathfrak{N}} / K\right) \rightarrow E_{\chi}
$$

Fix a cusp $\infty$ of $X_{1}(N)$ which is defined over $\mathbb{Q}$, and let

$$
\begin{equation*}
\Delta=\left[A^{\prime}, A^{\prime}[\mathfrak{N}], t\right]-(\infty) \quad \in J_{1}(N)\left(H_{c, \mathfrak{N}}\right) . \tag{3-9}
\end{equation*}
$$

To the pair $(f, \chi)$ we associate a Heegner point by letting $G=\operatorname{Gal}\left(H_{c, \mathfrak{N}} / K\right)$ and setting

$$
\begin{equation*}
P_{f}(\chi):=\sum_{\sigma \in G} \chi^{-1}(\sigma) \Phi_{f}\left(\Delta^{\sigma}\right) \quad \in B_{f}\left(H_{c, \mathfrak{N}}\right) \otimes_{T_{f}} E_{f, \chi} \tag{3-10}
\end{equation*}
$$

Note that, since $P_{f}(\chi)^{\sigma}=P_{f}(\chi)$ for any $\sigma \in \operatorname{Gal}\left(H_{c, \mathfrak{N}} / H_{\chi}\right)$, the point $P_{f}(\chi)$ lies in the subspace $B_{f}\left(H_{\chi}\right) \otimes_{T_{f}} E_{f, \chi}$. The embedding of $\overline{\mathbb{Q}}$ into $\mathbb{C}_{p}$ that was fixed from the outset allows us to consider the formal group logarithm

$$
\log _{\omega_{B_{f}}}: B_{f}\left(H_{c, \mathfrak{N}}\right) \rightarrow \mathbb{C}_{p} .
$$

We extend this function to $B_{f}\left(H_{c, \mathfrak{N}}\right) \otimes_{T_{f}} E_{f, \chi}$ by $E_{f, \chi}$-linearity.
Theorem 3.12. With notations and assumptions as above,

$$
L_{p}\left(f, \chi N_{K}\right)=\left(1-\chi^{-1}(\overline{\mathfrak{p}}) p^{-1} a_{p}(f)+\chi^{-2}(\overline{\mathfrak{p}}) \varepsilon(p) p^{-1}\right)^{2} \log _{\omega_{B_{f}}}^{2}\left(P_{f}(\chi)\right) .
$$

Proof. Let

$$
\mathscr{E}(f, \chi):=\left(1-\chi^{-1}(\overline{\mathfrak{p}}) p^{-1} a_{p}(f)+\chi^{-2}(\overline{\mathfrak{p}}) \varepsilon(p) p^{-1}\right)^{2} \quad \in E_{f, \chi}^{\times}
$$

be the Euler factor appearing in the statement of Theorem 3.12. Let $F^{\prime}$ denote the $p$-adic completion of $H_{c, \mathfrak{N}}$. Theorem 5.13 of [Bertolini et al. 2012b] in the case $k=2$ and $r=j=0$, with $\chi$ replaced by $\chi \mathrm{N}_{K}$, gives

$$
\begin{equation*}
L_{p}\left(f, \chi N_{K}\right)=\mathscr{E}(f, \chi) \times\left(\sum_{\sigma \in G} \chi^{-1}(\sigma) \cdot \operatorname{AJ}_{F^{\prime}}\left(\Delta^{\sigma}\right)\left(\omega_{f}\right)\right)^{2} \tag{3-11}
\end{equation*}
$$

Note that in this context, the $p$-adic Abel-Jacobi map $\mathrm{AJ}_{F^{\prime}}$ that appears in (3-11) is related to the formal group logarithm by

$$
\operatorname{AJ}_{F^{\prime}}(\Delta)\left(\omega_{f}\right)=\log _{\omega_{f}}(\Delta)
$$

Therefore,

$$
\begin{equation*}
L_{p}\left(f, \chi N_{K}\right)=\mathscr{E}(f, \chi)\left(\sum_{\sigma \in G} \chi^{-1}(\sigma) \log _{\omega_{f}}\left(\Delta^{\sigma}\right)\right)^{2} \tag{3-12}
\end{equation*}
$$

Theorem 3.12 follows from this formula and the fact that, by (3-8),

$$
\log _{\omega_{f}}(\Delta)=\log _{\Phi_{f}^{*}\left(\omega_{B_{f}}\right)}(\Delta)=\log _{\omega_{B_{f}}}\left(\Phi_{f}(\Delta)\right) .
$$

In the special case where $f$ has rational Fourier coefficients and $\chi=1$ is the trivial character, the abelian variety $B_{f}$ is an elliptic curve quotient of $J_{0}(N)$ and the Heegner point $P_{f}:=P_{f}(1)$ belongs to $B_{f}(K)$. Theorem 3.12 implies in this case that

$$
\begin{equation*}
L_{p}\left(f, \boldsymbol{N}_{K}\right)=\left(\frac{p+1-a_{p}(f)}{p}\right)^{2} \log ^{2}\left(P_{f}\right) \tag{3-13}
\end{equation*}
$$

where $\log : B_{f}\left(K_{\mathfrak{p}}\right) \rightarrow K_{\mathfrak{p}}$ is the formal group logarithm attached to a rational differential on $B_{f} / \mathbb{Q}$. Equation (3-13) exhibits a strong analogy with Theorem 1 of the Introduction, although it applies to arbitrary (modular) elliptic curves and not just elliptic curves with complex multiplication.

The remainder of Chapter 3 explains how Theorem 3.12 can in fact be used to prove Theorems 1 and 2 of the Introduction. The key to this proof is a relation between the Katz $p$-adic $L$-function of Section 3A and the $p$-adic Rankin $L$-function $L_{p}(f, \chi)$ of Section 3B in the special case where $f$ is a theta series attached to a Hecke character of the imaginary quadratic field $K$. This explicit relation is described in the following section.

3D. A factorization of the p-adic Rankin L-series. This section focuses on the Rankin $L$-function $L_{p}(f, \chi)$ of $f$ and $K$ in the special case where $f$ is a theta series associated to a Hecke character of the same imaginary quadratic field $K$.

More precisely, let $\psi$ be a fixed Hecke character of $K$ of infinity type ( $k-1,0$ ) with $k=r+2 \geq 2$. Consider the associated theta series

$$
\theta_{\psi}:=\sum_{\mathfrak{a}} \psi(\mathfrak{a}) q^{\mathrm{Na}}=\sum_{n=1}^{\infty} a_{n}\left(\theta_{\psi}\right) q^{n}
$$

where the first sum is taken over integral ideals of $K$. The Fourier coefficients of $\theta_{\psi}$ generate a number field $E_{\theta_{\psi}}$ which is clearly contained in $E_{\psi}$.

The following classical proposition is due to Hecke and Schoenberg. (See [Ogg 1969] or Section 3.2 of [Zagier 2008]).

Proposition 3.13. The theta series $\theta_{\psi}$ belongs to $S_{k}\left(\Gamma_{0}(N), \varepsilon\right)$, where
(1) the level $N$ is equal to $D M$, with $M=\mathrm{N}_{\mathbb{Q}}^{K} \mathfrak{f}_{\psi}$;
(2) the Nebentypus character $\varepsilon$ is equal to $\varepsilon_{K} \varepsilon_{\psi}$.

Lemma 3.14. If the conductor $\mathfrak{f}_{\psi}$ of $\psi$ is a cyclic ideal $\mathfrak{m}$ of norm $M$ prime to $D$, then $f:=\theta_{\psi}$ satisfies the Heegner hypothesis relative to $K$.
Proof. In this case, the modular form $\theta_{\psi}$ is of level $N=D M$, by Proposition 3.13. But then the ideal

$$
\begin{equation*}
\mathfrak{N}:=\mathfrak{d}_{K} \mathfrak{m} \tag{3-14}
\end{equation*}
$$

with $\mathfrak{d}_{K}:=\left(\sqrt{-D_{K}}\right)$, is a cyclic ideal of $K$ of norm $N$.
We will assume from now on that the condition in Lemma 3.14 is satisfied. Furthermore, we will always take $\mathfrak{N}$ to be the ideal in (3-14).

The goal of this section is to factor the $p$-adic Rankin $L$-function $L_{p}\left(\theta_{\psi}, \cdot\right)$ as a product of two Katz $p$-adic $L$-functions. As a preparation to stating the main result, we record the following two lemmas:
Lemma 3.15. For $f:=\theta_{\psi}$, the $L$-function $L\left(f, \chi^{-1}, s\right)$ factors as

$$
L\left(f, \chi^{-1}, s\right)=L\left(\psi \chi^{-1}, s\right) \cdot L\left(\psi^{*} \chi^{-1}, s\right) .
$$

Proof. Let ( $\rho_{f, \ell}$ ) denote the compatible system of $\ell$-adic Galois representations associated to $f$. The factorization above then follows from the fact that $L\left(f, \chi^{-1}, s\right)$ is the $L$-function of the compatible system of Galois representations

$$
\left.\rho_{f, \ell}\right|_{\mathrm{Gal}(\bar{K} / K)} \otimes \chi_{\ell}^{-1}=\left(\psi_{\ell} \oplus \psi_{\ell}^{*}\right) \otimes \chi_{\ell}^{-1}=\psi_{\ell} \chi_{\ell}^{-1} \oplus \psi_{\ell}^{*} \chi_{\ell}^{-1} .
$$

Lemma 3.16. Let $c$ be an integer prime to $p N$ and let $\chi$ be any character in $\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$.
(1) If $\chi$ belongs to $\Sigma_{\mathrm{cc}}^{(2)}(c, \mathfrak{N}, \varepsilon)$, then $\psi^{-1} \chi$ belongs to $\Sigma_{\mathrm{sd}}^{(2)}\left(c \mathfrak{d}_{K}\right)^{-}$and $\psi^{*-1} \chi$ belongs to $\Sigma_{\mathrm{sd}}^{(2)}\left(c \mathfrak{d}_{K} M\right)^{-}$.
(2) If $\chi$ belongs to $\Sigma_{\mathrm{cc}}^{(1)}(c, \mathfrak{N}, \varepsilon)$, then $\psi^{-1} \chi$ belongs to $\Sigma_{\mathrm{sd}}^{(1)}\left(c \mathfrak{d}_{K}\right)^{-}$and $\psi^{*-1} \chi$ belongs to $\Sigma_{\mathrm{sd}}^{(2)}\left(c \mathfrak{d}_{K} M\right)^{-}$.
Proof. We first note that when $\chi$ is of type $(k+j,-j)$, then $\psi^{-1} \chi$ is of infinity type $(1+j,-j)$ and $\psi^{*-1} \chi$ is of infinity type $(k+j, 1-(k+j))$. Since $\chi \in \Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$, we have

$$
\varepsilon_{\chi}=\varepsilon=\varepsilon_{\psi} \cdot \varepsilon_{K}
$$

Thus $\varepsilon_{\psi^{-1} \chi}$ equals $\varepsilon_{K}$ and the same holds for $\varepsilon_{\psi^{*-1} \chi}$ since $\varepsilon_{\psi^{*}}=\varepsilon_{\psi}$. It follows then from Remark 3.7 that $\psi^{-1} \chi$ and $\psi^{*-1} \chi$ are self-dual characters.

Let $q$ be a rational prime dividing $M$. Since $\mathfrak{m}$ is a cyclic $0_{K}$-ideal, it follows that $q=\mathfrak{q} \overline{\mathfrak{q}}$ must be split in $K$, and exactly one of $\mathfrak{q}, \overline{\mathfrak{q}}$ divides $\mathfrak{m}$. From this, it is easy to see that $\varepsilon_{\psi}$ has conductor exactly $M$, and hence $\varepsilon$ has conductor exactly $N$ and $\mathfrak{N}_{\varepsilon}=\mathfrak{N}$. Thus $\mathfrak{f}_{\chi}=c \mathfrak{N}=c \mathfrak{d}_{K} \mathfrak{m}$ and $\mathfrak{f}_{\psi^{*-1} \chi}=c \mathfrak{d}_{K} \mathfrak{m} \overline{\mathfrak{m}}=c \mathfrak{d}_{K} M$. On the other hand, since $\varepsilon_{\chi}=\varepsilon_{\psi} \varepsilon_{K}$, it follows that $\mathfrak{f}_{\psi^{-1} \chi}=c \mathfrak{d}_{K}$.

The preceding remarks imply that if $\chi$ is in $\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$, then $\psi^{-1} \chi$ lies in $\Sigma_{\mathrm{sd}}\left(c \mathfrak{d}_{K}\right)^{-}$and $\psi^{*-1} \chi$ lies in $\Sigma_{\mathrm{sd}}\left(c \mathfrak{d}_{K} M\right)^{-}$. To finish, we note that if $j \geq 0$, then both $\psi^{-1} \chi$ and $\psi^{*-1} \chi$ lie in $\Sigma_{\text {sd }}^{(2)}$, while if $-(k-1) \leq j \leq-1$, then $\psi^{*-1} \chi$ is in $\Sigma_{\text {sd }}^{(2)}$ while $\psi^{-1} \chi$ lies in $\Sigma_{\text {sd }}^{(1)}$.
Theorem 3.17. For all $\chi \in \Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$,

$$
\begin{align*}
& L_{p}\left(\theta_{\psi}, \chi\right)  \tag{3-15}\\
& =\frac{w\left(\theta_{\psi}, \chi\right)^{-1} w_{K}}{2 c^{k+2 j-1}} \prod_{q \mid c} \frac{q-\varepsilon_{K}(q)}{q-1} \times \mathscr{L}_{p, c \mathrm{o}_{K}}\left(\psi^{-1} \chi\right) \times \mathscr{L}_{p, c \mathrm{o}_{K} M}\left(\psi^{*-1} \chi\right) .
\end{align*}
$$

Proof. Since $\Sigma_{\text {cc }}^{(2)}(c, \mathfrak{N}, \varepsilon)$ is dense in $\hat{\Sigma}_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$, it suffices to prove the formula for the characters $\chi$ in this range, where it follows directly from the interpolation properties defining the respective $p$-adic $L$-functions. More precisely, by (3-7),

$$
\begin{equation*}
\frac{L_{p}\left(\theta_{\psi}, \chi\right)}{\Omega_{p}(A)^{2(k+2 j)}}=\left(\left(1-\psi \chi^{-1}(\overline{\mathfrak{p}})\right)\left(1-\psi^{*} \chi^{-1}(\overline{\mathfrak{p}})\right)\right)^{2} L_{\mathrm{alg}}\left(\theta_{\psi}, \chi^{-1}, 0\right) \tag{3-16}
\end{equation*}
$$

Let

$$
\delta_{c}:=\prod_{q \mid c} \frac{q-\varepsilon_{K}(q)}{q-1} .
$$

By Lemma 3.15 and the definition of $L_{\text {alg }}\left(\theta_{\psi}, \chi^{-1}, 0\right)$ given in (3-5) and (3-6),

$$
\begin{align*}
L_{\mathrm{alg}}\left(\theta_{\psi}, \chi^{-1}, 0\right)= & w\left(\theta_{\psi}, \chi\right)^{-1} C\left(\theta_{\psi}, \chi, c\right) \frac{L\left(\theta_{\psi}, \chi^{-1}, 0\right)}{\Omega(A)^{2(k+2 j)}}  \tag{3-17}\\
= & w\left(\theta_{\psi}, \chi\right)^{-1} w_{K} \delta_{c} \frac{2^{r+2 j} \pi^{k+2 j-1} j!(k+j-1)!}{\sqrt{D}^{k+2 j-1} c^{k+2 j-1}} \\
& \times \frac{L\left(\psi \chi^{-1}, 0\right) L\left(\psi^{*} \chi^{-1}, 0\right)}{\Omega(A)^{2(k+2 j)}} \\
= & \frac{w\left(\theta_{\psi}, \chi\right)^{-1} w_{K} \delta_{c}}{2 c^{k+2 j-1}}\left(\frac{j!(2 \pi)^{j} L\left(\psi \chi^{-1}, 0\right)}{\sqrt{D}^{j} \Omega(A)^{1+2 j}}\right) \\
& \times\left(\frac{(k+j-1)!(2 \pi)^{k+j-1} L\left(\psi^{*} \chi^{-1}, 0\right)}{\sqrt{D}^{k+j-1} \Omega(A)^{1+2(k+j-1)}}\right) .
\end{align*}
$$

Combining (3-16) and (3-17) with the interpolation property of the Katz p-adic $L$-function given in Proposition 3.5, we obtain

$$
\begin{equation*}
\frac{L_{p}\left(\theta_{\psi}, \chi\right)}{\Omega_{p}(A)^{2(k+2 j)}}=\frac{w\left(\theta_{\psi}, \chi\right)^{-1} w_{K} \delta_{c}}{2 c^{k+2 j-1}} \times \frac{\mathscr{L}_{p, c \mathrm{c}_{K}}\left(\psi^{-1} \chi\right)}{\Omega_{p}(A)^{1+2 j}} \times \frac{\mathscr{L}_{p, c \mathrm{c}_{K} M}\left(\psi^{*-1} \chi\right)}{\Omega_{p}(A)^{1+2(k+j-1)}} . \tag{3-18}
\end{equation*}
$$

Clearing the powers of $\Omega_{p}(A)$ on both sides gives the desired result.

The Nebentypus character $\varepsilon$ can be viewed as a finite-order Galois character of $G_{\mathbb{Q}}$. Recall that $E_{\psi, \chi, \varepsilon}$ denotes the smallest extension of $E_{\psi, \chi}$ containing the field through which this character factors.

Corollary 3.18. For all $\chi \in \Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$,

$$
L_{p}\left(\theta_{\psi}, \chi\right)=\mathscr{L}_{p, c \mathfrak{o}_{K}}\left(\psi^{-1} \chi\right) \times \mathscr{L}_{p, c \mathfrak{o}_{K} M}\left(\psi^{*-1} \chi\right)\left(\bmod E_{\psi, \chi, \varepsilon}^{\times}\right) .
$$

Proof. This follows from Theorem 3.17 in light of the fact that the constant that appears on the right-hand side of (3-15) belongs to $E_{\psi, \chi, \varepsilon}^{\times}$.

3E. Proof of Rubin's theorem. The goal of this section is to prove Theorem 2 of the Introduction. Let $\mathfrak{c}=\overline{\mathfrak{c}}$ be an integral ideal in $\mathbb{O}_{K}$ invariant under complex conjugation and let $v \in \Sigma_{\mathrm{sd}}(\mathfrak{c})^{-}$be a Hecke character of $K$ of infinity type $(1,0)$. Since $\varepsilon_{v}=\varepsilon_{K}$, it follows that $\mathfrak{d}_{K} \mid \mathfrak{c}$. We will also assume that $v$ satisfies the following additional conditions:
(i) The sign $w_{\nu}$ of the functional equation of the $L$-function $L(\nu, s)$ is -1 .
(ii) $\mathfrak{d}_{K} \| \mathfrak{c}$. Thus $\mathfrak{c}=(c) \mathfrak{d}_{K}$ for a unique positive rational integer $c$ that is prime to $D$.

Let $p$ be a rational prime split in $K$ that is prime to $c$.
Definition 3.19. A pair $(\psi, \chi)$ of Hecke characters is said to be good for $v$ if it satisfies the following conditions.
(1) The character $\psi$ is of type $(1,0)$ and has conductor $\mathfrak{m}$, where $\mathfrak{m}$ is a cyclic ${O_{K}}^{-}$ ideal prime to $p D$. Thus $\theta_{\psi}$ is a newform in $S_{2}\left(\Gamma_{0}(N), \varepsilon\right)$, where $N=M D$ and $\varepsilon=\varepsilon_{\psi} \varepsilon_{K}$ is a Dirichlet character of conductor exactly $N$. Let $\mathfrak{N}:=\mathfrak{m d}_{K}$.
(2) The character $\chi$ is of finite order, and $\chi \boldsymbol{N}_{K}$ belongs to $\Sigma_{\mathrm{cc}}^{(1)}(c, \mathfrak{N}, \varepsilon)$. This implies (on account of Lemma 3.16 applied to $\chi \boldsymbol{N}_{K}$ ) that $\psi^{-1} \chi \boldsymbol{N}_{K}$ lies in $\Sigma_{\text {sd }}^{(1)}(\mathfrak{c})$ and $\psi^{*-1} \chi N_{K}$ lies in $\Sigma_{\text {sd }}^{(2)}(\mathfrak{c} M)$.
(3) The character $\psi \chi^{-1}$ is equal to $v$, that is, $\psi^{-1} \chi N_{K}=v^{*}$.
(4) The classical $L$-value $L\left(\psi^{*} \chi^{-1} N_{K}^{-1}, 0\right)=L\left(\psi^{*} \chi^{-1}, 1\right)$ is nonzero, and hence $\mathscr{L}_{p, \mathrm{c} M}\left(\psi^{*-1} \chi N_{K}\right) \neq 0$.
Remark 3.20. Suppose that a pair $(\psi, \chi)$ satisfies (1) and (3) above with $\mathfrak{m}$ prime to $\mathfrak{c}$. Then such a pair automatically satisfies (2) also. Indeed, the character $\chi \boldsymbol{N}_{K}=\psi \nu^{*}$ is of type $(1,1)$ and its central character is equal to

$$
\varepsilon_{\chi}=\varepsilon_{\psi} \varepsilon_{\nu^{*}}=\varepsilon_{\psi} \varepsilon_{K}=\varepsilon
$$

where $\varepsilon$ is the Nebentypus character attached to $\theta_{\psi}$. Further, $\mathfrak{f}_{\chi}=\mathfrak{f}_{\psi} \mathfrak{f}_{\nu^{*}}=\mathfrak{m} \cdot \boldsymbol{c d _ { K }}$. It follows that the character $\chi \boldsymbol{N}_{K}$ belongs to $\Sigma_{\mathrm{cc}}(c, \mathfrak{N}, \varepsilon)$, with $\mathfrak{N}=\mathfrak{d}_{K} \mathfrak{m}$. (The set $S_{\theta_{\psi}}$ in the discussion below Definition 3.10 is empty, since $D \mid N_{\varepsilon}$.)

Remark 3.21. Suppose that a pair ( $\psi, \chi$ ) satisfies conditions (1), (2) and (3) above. Since $\chi \boldsymbol{N}_{K}$ lies in $\Sigma_{\mathrm{cc}}^{(1)}(c, \mathfrak{N}, \varepsilon)$, the sign in the functional equation of $L\left(\theta_{\psi}, \chi^{-1}, s\right)$ is -1 . As seen previously, this $L$-function factors as

$$
\begin{aligned}
L\left(\theta_{\psi},\left(\chi \boldsymbol{N}_{K}\right)^{-1}, s\right) & =L\left(\psi \chi^{-1} \boldsymbol{N}_{K}^{-1}, s\right) L\left(\psi^{*} \chi^{-1} \boldsymbol{N}_{K}^{-1}, s\right) \\
& =L\left(\nu \boldsymbol{N}_{K}^{-1}, s\right) L\left(\psi^{*} \chi^{-1} \boldsymbol{N}_{K}^{-1}, s\right) .
\end{aligned}
$$

The normalization here is such that the central point is $s=0$. Since the sign of $L(v, s)$ is -1 , it follows that the sign of $L\left(\psi^{*} \chi^{-1} N_{K}^{-1}, s\right)$ is +1 . Hence condition (4) would be expected to hold generically.

The modular abelian variety $B_{\theta_{\psi}}$ attached to $\psi$ is a CM abelian variety in the sense of Definition 2.4. Hence it is $K$-isogenous to the CM abelian variety $B_{\psi}$ constructed in Section 2B. In particular, the modular parametrization $\Phi_{\psi}:=\Phi_{\theta_{\psi}}$ can be viewed as a surjective morphism of abelian varieties over $K$ :

$$
\begin{equation*}
\Phi_{\psi}: J_{1}(N) \rightarrow B_{\psi} . \tag{3-19}
\end{equation*}
$$

Given a good pair $(\psi, \chi)$, recall the Heegner divisor $\Delta \in J_{1}(N)\left(H_{c, \mathfrak{N}}\right)$ that was constructed in Section 3C, and the Heegner point

$$
\begin{equation*}
P_{\psi}(\chi):=P_{\theta_{\psi}}(\chi)=\sum_{\sigma \in G} \Phi_{\psi}\left(\Delta^{\sigma}\right) \otimes \chi^{-1}(\sigma) \in B_{\psi}\left(H_{\chi}\right) \otimes_{T_{\psi}} E_{\psi, \chi} \tag{3-20}
\end{equation*}
$$

that was defined in Equation (3-10) of that section with $f=\theta_{\psi}$. Recall also that $\omega_{\psi}$ is an $E_{\psi}$-vector space generator of $\Omega^{1}\left(B_{\psi} / E_{\psi}\right)^{T_{\psi}}$. Viewing the point $P_{\psi}(\chi)$ as a formal linear combination of elements of $B_{\psi}\left(H_{\chi}\right)$ with coefficients in $E_{\psi, \chi}$, we define the expression $\log _{\omega_{\psi}}\left(P_{\psi}(\chi)\right)$ by $E_{\psi, \chi \text {-linearity. }}$

In the rest of this section, we will denote by $E_{\psi, \chi}^{\prime}$ the subfield of $\overline{\mathbb{Q}}$ generated by $E_{\psi}, E_{\chi}$, and the abelian extension $H_{\chi}^{\prime}$ of $K$ cut out by the finite-order characters $\chi$ and $\chi^{*}$. The motivation for singling out good pairs for a special definition lies in the following proposition.
Proposition 3.22. For any pair $(\psi, \chi)$ which is good for $\nu$,

$$
\begin{equation*}
\mathscr{L}_{p, \mathfrak{c}}\left(\nu^{*}\right)=\Omega_{p}\left(\nu^{*}\right)^{-1} \log _{\omega_{\psi}}^{2}\left(P_{\psi}(\chi)\right)\left(\bmod \left(E_{\psi, \chi}^{\prime}\right)^{\times}\right), \tag{3-21}
\end{equation*}
$$

where $\Omega_{p}\left(v^{*}\right)$ is the p-adic period from Definition 2.13.
Proof. By Theorem 3.12 applied to $f=\theta_{\psi}$,

$$
\begin{equation*}
L_{p}\left(\theta_{\psi}, \chi \boldsymbol{N}_{K}\right)=\log _{\omega_{\psi}}^{2}\left(P_{\psi}(\chi)\right)\left(\bmod E_{\psi, \chi}^{\times}\right) \tag{3-22}
\end{equation*}
$$

On the other hand, since $E_{\psi, \chi}^{\prime}$ contains $E_{\psi, \chi, \varepsilon}$, Corollary 3.18 implies that

$$
\begin{align*}
L_{p}\left(\theta_{\psi}, \chi \boldsymbol{N}_{K}\right) & =\mathscr{L}_{p, \mathrm{c}}\left(\psi^{-1} \chi \boldsymbol{N}_{K}\right) \mathscr{L}_{p, \mathrm{c} M}\left(\psi^{*-1} \chi \boldsymbol{N}_{K}\right)\left(\bmod \left(E_{\psi, \chi}^{\prime}\right)^{\times}\right)  \tag{3-23}\\
& =\mathscr{L}_{p, \mathfrak{c}}\left(v^{*}\right) \mathscr{L}_{p, \mathrm{c} M}\left(\psi^{*-1} \chi \boldsymbol{N}_{K}\right)\left(\bmod \left(E_{\psi, \chi}^{\prime}\right)^{\times}\right),
\end{align*}
$$

where the second equality follows from condition (3) in the definition of a good pair. The value $\mathscr{L}_{p, c M}\left(\psi^{*-1} \chi N_{K}\right)$ is nonzero by condition (4) in the definition of a good pair. Therefore, by Corollary 3.3,

$$
\begin{equation*}
\mathscr{L}_{p, c M}\left(\psi^{*-1} \chi \boldsymbol{N}_{K}\right)=\Omega_{p}\left(\psi^{-1} \chi^{*} \boldsymbol{N}_{K}\right)\left(\bmod E_{\psi, \chi}^{\times}\right) . \tag{3-24}
\end{equation*}
$$

Since $\psi \chi^{-1}=v$, we have

$$
\begin{align*}
\Omega_{p}\left(\psi^{-1} \chi^{*} \boldsymbol{N}_{K}\right) & =\Omega_{p}\left(v^{-1} \chi^{-1} \chi^{*} \boldsymbol{N}_{K}\right)  \tag{3-25}\\
& =\Omega_{p}\left(v^{*} \cdot \chi^{*} / \chi\right)=\Omega_{p}\left(v^{*}\right)\left(\bmod \left(E_{\psi, \chi}^{\prime}\right)^{\times}\right),
\end{align*}
$$

where the last equality follows from Lemma 2.14. The proposition now follows by combining Equations (3-22) through (3-25).

To go further, we will analyze the expression $\log _{\omega_{\psi}}\left(P_{\psi}(\chi)\right)$ and relate it to quantities depending solely on $v$ and not on the good pair $(\psi, \chi)$. It will be useful to view the point $P_{\psi}(\chi)$ appearing in (3-21) as an element of $B_{\psi, \chi}^{0}\left(H_{c, \mathfrak{N}}\right)$ or as a $K$-rational point on the abelian variety $B_{\psi, \chi}$ that was introduced in Section 2B. More precisely, after setting

$$
\begin{equation*}
\boldsymbol{P}_{\psi}(\chi):=\sum_{\sigma \in G} \Phi_{\psi}\left(\Delta^{\sigma}\right) \otimes \chi^{-1}(\sigma) \in B_{\psi}(\bar{K}) \otimes_{T_{\psi}} T_{\psi, \chi}=B_{\psi, \chi}^{0}(\bar{K}), \tag{3-26}
\end{equation*}
$$

we observe that, for all $\tau \in \operatorname{Gal}(\bar{K} / K)$,

$$
\tau *_{0} \boldsymbol{P}_{\psi}(\chi)=\sum_{\sigma \in G} \Phi_{\psi}\left(\Delta^{\tau \sigma}\right) \otimes \chi^{-1}(\sigma)=\sum_{\sigma \in G} \Phi_{\psi}\left(\Delta^{\sigma}\right) \otimes \chi^{-1}\left(\sigma \tau^{-1}\right)=\boldsymbol{P}_{\psi}(\chi) \chi(\tau) .
$$

The point $\boldsymbol{P}_{\psi}(\chi)$ therefore belongs to $B_{\psi, \chi}(K)$ by (2-9). For the following lemmas, recall the differentials $\omega_{\psi, \chi}^{0} \in \Omega^{1}\left(B_{\psi, \chi}^{0} / E_{\psi, \chi}\right)^{T_{\psi, x}}$ and $\omega_{\psi, \chi} \in \Omega^{1}\left(B_{\psi, \chi} / E_{\psi, \chi}\right)^{T_{\psi, \chi}}$.
Lemma 3.23. For all good pairs $(\psi, \chi)$ attached to $v=\psi \chi^{-1}$,

$$
\log _{\omega_{\psi}}\left(P_{\psi}(\chi)\right)=\log _{\omega_{\psi, \chi}^{0}}\left(\boldsymbol{P}_{\psi}(\chi)\right)
$$

Proof. Let $G=\operatorname{Gal}\left(H_{c, \mathfrak{N}} / K\right)$ and let $P=\Phi_{\psi}(\Delta)$. Also, let $i$ be the map defined in (2-4). Then

$$
\begin{aligned}
\log _{\omega_{\psi}}\left(P_{\psi}(\chi)\right) & =\sum_{\sigma \in G} \chi(\sigma)^{-1} \log _{\omega_{\psi}}\left(P^{\sigma}\right)=\sum_{\sigma \in G} \chi(\sigma)^{-1} \log _{i^{*} * \omega_{\psi, \chi}^{0}}\left(P^{\sigma}\right) \\
& =\sum_{\sigma \in G} \chi(\sigma)^{-1} \log _{\omega_{\psi, \chi}^{0}}\left(P^{\sigma} \otimes 1\right)=\sum_{\sigma \in G} \log _{\chi(\sigma)^{-1} \omega_{\psi, \chi}^{0}}\left(P^{\sigma} \otimes 1\right) \\
& =\sum_{\sigma \in G} \log _{\omega_{\psi, \chi}^{0}}\left(P^{\sigma} \otimes \chi(\sigma)^{-1}\right) \\
& =\log _{\omega_{\psi, \chi}^{0}}\left(\sum_{\sigma \in G} P^{\sigma} \otimes \chi(\sigma)^{-1}\right)=\log _{\omega_{\psi, \chi}^{0}}\left(\boldsymbol{P}_{\psi, \chi}\right) .
\end{aligned}
$$

Lemma 3.24. $\quad \log _{\omega_{\psi, \chi}^{0}}\left(\boldsymbol{P}_{\psi}(\chi)\right)=\log _{\omega_{\psi, \chi}}\left(\boldsymbol{P}_{\psi}(\chi)\right)\left(\bmod \left(E_{\psi, \chi}^{\prime}\right)^{\times}\right)$.
Proof. This follows from Lemma 2.8 since the Gauss sum $\mathfrak{g}(\chi)$ lies in $\left(E_{\psi, \chi}^{\prime}\right)^{\times}$. $\square$
Lemma 3.25. There exist $P_{v} \in B_{v}(K) \otimes \mathbb{Q}$ and $\omega_{\nu} \in \Omega^{1}\left(B_{v} / E_{\nu}\right)^{T_{v}}$ such that

$$
\log _{\omega_{\psi, \chi}}\left(\boldsymbol{P}_{\psi}(\chi)\right)=\log _{\omega_{\nu}}\left(P_{\nu}\right)\left(\bmod \left(E_{\psi, \chi}^{\prime}\right)^{\times}\right)
$$

The point $P_{\nu}$ is nonzero if and only if $L(v, s)$ vanishes to exact order one at $s=1$.
Proof. Recall from Lemma 2.9 that there is a $K$-rational isogeny

$$
B_{v} \otimes_{T_{v}} T_{\psi, \chi} \rightarrow B_{\psi, \chi} .
$$

Composing it with the natural morphism $B_{v} \rightarrow B_{v} \otimes_{T_{v}} T_{\psi, \chi}$, we obtain a $T_{\nu^{-}}$ equivariant morphism $j: B_{v} \rightarrow B_{\psi, \chi}$ defined over $K$ with finite kernel. By the Gross-Zagier theorem (see [Gross and Zagier 1986; Yuan et al. 2011] and the remark below), the point $\boldsymbol{P}_{\psi}(\chi)$ is nonzero if and only if $L^{\prime}\left(\theta_{\psi}, \chi^{-1}, 1\right) \neq 0$. By Remark 3.21, we have

$$
L^{\prime}\left(\theta_{\psi}, \chi^{-1}, 1\right)=L^{\prime}(\nu, 1) \cdot L\left(\psi^{*} \chi^{-1}, 1\right) .
$$

Since $L\left(\psi^{*} \chi^{-1}, 1\right) \neq 0$, we see that $\boldsymbol{P}_{\psi}(\chi)$ is nonzero if and only if $L^{\prime}(\nu, 1) \neq 0$. Thus, if $L(\nu, s)$ vanishes to order strictly greater than one (hence order $\geq 3$ ), the lemma holds with $P_{\nu}:=0$.

We may assume therefore that $L(v, s)$ has a simple zero at $s=1$. This implies (by [Gross and Zagier 1986] and [Kolyvagin 1990]; see also [Kolyvagin and Logachëv 1989]) that $B_{v}(K) \otimes \mathbb{Q}$ is one-dimensional over $E_{v}$, and therefore that $B_{\psi, \chi}(K) \otimes \mathbb{Q}$ is one-dimensional over $E_{\psi, \chi}$. In particular, if $P_{\nu}$ is any generator of $B_{\nu}(K) \otimes \mathbb{Q}$, we may write

$$
\boldsymbol{P}_{\psi}(\chi)=\lambda j\left(P_{v}\right)
$$

for some nonzero scalar $\lambda \in E_{\psi, \chi}^{\times}$. But letting

$$
\omega_{\nu}=j^{*}\left(\omega_{\psi, \chi}\right) \in \Omega^{1}\left(B_{v} / E_{\psi, \chi}^{\prime}\right)^{T_{v}},
$$

we have

$$
\begin{aligned}
\log _{\omega_{\psi, \chi}}\left(\boldsymbol{P}_{\psi}(\chi)\right) & =\log _{\omega_{\psi, \chi}}\left(\lambda j\left(P_{\nu}\right)\right)=\log _{\lambda^{*} \omega_{\omega_{, \chi}}}\left(j\left(P_{\nu}\right)\right)=\lambda \log _{\omega_{\psi, \chi}}\left(j\left(P_{\nu}\right)\right) \\
& =\lambda \log _{j^{*} \omega_{\psi, \chi}}\left(P_{\nu}\right)=\lambda \log _{\omega_{\nu}}\left(P_{\nu}\right) .
\end{aligned}
$$

The lemma now follows after multiplying $\omega_{\nu}$ by an appropriate scalar in $\left(E_{\psi, \chi}^{\prime}\right)^{\times}$ so that it belongs to $\Omega^{1}\left(B_{v} / E_{v}\right)^{T_{v}}$.
Remark 3.26. The original result of [Gross and Zagier 1986] is not general enough to include the situation above. However, Yuan et al. [2011] have proved a very general GZ formula with no assumptions on ramification. This formula relates the height of a Heegner point to a derivative of a Rankin-Selberg $L$-function, but
involves some extra local integrals at bad places that depend on particular choices of test vectors. To deduce that the Heegner point is nontorsion from the nonvanishing of the derivative of the $L$-function, one needs in addition to know that the local zeta integrals are nonzero; this follows in our case from the computations of [Bertolini et al. 2012b, Section 4.6].
Proposition 3.27. There exist $\omega_{\nu} \in \Omega^{1}\left(B_{v} / E_{v}\right)^{T_{v}}$ and $P_{v} \in B_{v}(K) \otimes \mathbb{Q}$ such that

$$
\begin{equation*}
\mathscr{L}_{p, \mathfrak{c}}\left(v^{*}\right)=\Omega_{p}\left(v^{*}\right)^{-1} \log _{\omega_{v}}^{2}\left(P_{\nu}\right)\left(\bmod \left(E_{\psi, \chi}^{\prime}\right)^{\times}\right), \tag{3-27}
\end{equation*}
$$

for all good pairs $(\psi, \chi)$ attached to $\nu$. The point $P_{\nu}$ is nonzero if and only if $L(v, s)$ vanishes to exact order one at $s=1$.

Proof. This follows immediately from Proposition 3.22 and Lemmas 3.23 through 3.25.

While Proposition 3.27 brings us close to Theorem 2 of the Introduction, it is somewhat more vague in that both sides of the purported equality may differ a priori by a nonzero element of the typically larger field $E_{\psi, \chi}^{\prime}$. The alert reader will also notice that this proposition is potentially vacuous for now, because the existence of a good pair for $v$ has not yet been established! The next proposition repairs this omission, and directly implies Theorem 2 of the Introduction.

Proposition 3.28. The set $S_{v}$ of pairs $(\psi, \chi)$ that are good for $v$ is nonempty. Furthermore,

$$
\begin{equation*}
\bigcap_{(\psi, x) \in S_{v}} E_{\psi, \chi}^{\prime}=E_{v} . \tag{3-28}
\end{equation*}
$$

The proof of Proposition 3.28 rests crucially on a nonvanishing result of Rohrlich [1984] and Greenberg [1985] for the central critical values of Hecke $L$-series. In order to state it, we fix a rational prime $\ell$ which is split in $K$ and let

$$
K_{\infty}^{-}=\bigcup_{n \geq 0} K_{n}^{-}
$$

be the so-called anticyclotomic $\mathbb{Z}_{\ell}$ extension of $K$; it is the unique $\mathbb{Z}_{\ell}$-extension of $K$ which is Galois over $\mathbb{Q}$ and for which $\operatorname{Gal}\left(K_{\infty}^{-} / \mathbb{Q}\right)=\mathbb{Z}_{\ell} \rtimes(\mathbb{Z} / 2 \mathbb{Z})$ is a generalized dihedral group.

Lemma 3.29 [Greenberg 1985; Rohrlich 1984]. Let $\psi_{0}$ be a self-dual Hecke character of $K$ of infinity type ( 1,0 ). Assume that the sign $w_{\psi_{0}}$ in the functional equation of $L\left(\psi_{0}, s\right)$ is equal to 1 . Then there are infinitely many finite-order characters $\chi$ of $\operatorname{Gal}\left(K_{\infty}^{-} / K\right)$ for which $L\left(\psi_{0} \chi, 1\right) \neq 0$.
Proof. Let $\mathfrak{c}^{\prime}$ be the conductor of $\psi_{0}$. In light of the hypothesis that $w_{\psi_{0}}=1$, Theorem 1 of [Greenberg 1985] implies that the Katz $p$-adic $L$-function (with $p=\ell$ ) does not vanish identically on any open $\ell$-adic neighborhood of $\psi_{0}$ in
$\Sigma_{\text {sd }}\left(\mathfrak{c}^{\prime}\right)$. (See the discussion in the first paragraph of the proof of Proposition 1 on p. 93 of [Greenberg 1985].) If $U$ is any sufficiently small such neighborhood, then:
(1) The restriction to $U$ of the Katz $p$-adic $L$-function is described by a power series with $p$-adically bounded coefficients, and therefore admits only finitely many zeros by the Weierstrass preparation theorem.
(2) The region $U$ contains a dense subset of points of the form $\psi_{0} \chi$, where $\chi$ is a finite-order character of $\operatorname{Gal}\left(K_{\infty}^{-} / K\right)$.

Lemma 3.29 follows directly from these two facts.
Proof of Proposition 3.28. Let $\bar{S}_{v} \supset S_{v}$ be the set of pairs satisfying conditions (1)-(3) in the definition of a good pair, but without necessarily requiring the more subtle fourth condition. The proof of Proposition 3.28 will be broken down into four steps.
Step 1. The set $\bar{S}_{v}$ is nonempty.
To see this, let $\psi$ be any Hecke character of $K$ of infinity type $(1,0)$ and conductor $\mathfrak{m}$, where $\mathfrak{m}$ is a cyclic $\mathbb{O}_{K}$-ideal prime to $\mathfrak{c}$. Setting $\chi=\psi v^{-1}$, the pair $(\psi, \chi)$ satisfies conditions (1) and (3) by construction, and (2) as well on account of Remark 3.20. Therefore, the pair $(\psi, \chi)$ belongs to $\bar{S}_{v}$.

Step 2. Given $(\psi, \chi) \in \bar{S}_{v}$, there exist $\left(\psi_{1}, \chi_{1}\right)$ and $\left(\psi_{2}, \chi_{2}\right) \in S_{v}$ with

$$
E_{\psi_{1}, \chi_{1}}^{\prime} \cap E_{\psi_{2}, \chi_{2}}^{\prime} \subset E_{\psi, \chi}^{\prime}
$$

To see this, let $\ell=\lambda \bar{\lambda}$ be a rational prime which splits in $K$ and is relatively prime to the class number of $K$ and the conductors of $\psi$ and $\chi$, and which is unramified in $E_{\psi, \chi}^{\prime} / \mathbb{Q}$. For such a prime, let

$$
K_{\infty}=\bigcup_{n \geq 0} K_{n}, \quad K_{\infty}^{\prime}=\bigcup_{n \geq 0} K_{n}^{\prime}
$$

be the unique $\mathbb{Z}_{\ell}$-extensions of $K$ which are unramified outside of $\lambda$ and $\bar{\lambda}$ respectively, with $\left[K_{n}: K\right]=\ell^{n}$ and likewise for $K_{n}^{\prime}$. The condition that $\ell$ does not divide the class number of $K$ implies that the fields $K_{n}$ and $K_{n}^{\prime}$ are totally ramified at $\lambda$ and $\bar{\lambda}$ respectively. If $\alpha$ is any character of $\operatorname{Gal}\left(K_{\infty} / K\right)$, the pair $\left(\psi_{1}, \chi_{1}\right):=(\psi \alpha, \chi \alpha)$ still belongs to $\bar{S}_{v}$, with $\mathfrak{m}$ in condition 1 replaced by $\mathfrak{m} \lambda^{n}$ for a suitable $n \geq 0$. Furthermore,

$$
\begin{equation*}
L\left(\psi_{1}^{*} \chi_{1}^{-1} \boldsymbol{N}_{K}^{-1}, 0\right)=L\left(\psi^{*} \chi^{-1} \boldsymbol{N}_{K}^{-1} \cdot\left(\alpha^{*} / \alpha\right), 0\right) \tag{3-29}
\end{equation*}
$$

The character $\alpha^{*} / \alpha$ is an anticyclotomic character of $K$ of $\ell$-power order and conductor, and all such characters can be obtained by choosing $\alpha$ appropriately. The fact that ( $\psi, \chi$ ) satisfies conditions (1)-(3) of a good pair implies (see Remark 3.21) that the sign $w_{\psi^{*} \chi^{-1}}$ is equal to +1 . Hence, by Lemma 3.29, there exists a choice
of $\alpha$ for which the $L$-value appearing on the right of (3-29) is nonvanishing. The corresponding pair ( $\psi_{1}, \chi_{1}$ ) belongs to $S_{v}$ and satisfies

$$
E_{\psi_{1}, \chi_{1}}^{\prime} \subset E_{\psi, \chi, \ell, n}^{\prime}:=E_{\psi, \chi}^{\prime} \mathbb{Q}\left(\zeta_{\ell^{n}}\right) K_{n} K_{n}^{\prime}
$$

for some $n$. Note that the extension $E_{\psi, \chi, \ell, n}^{\prime} / E_{\psi, \chi}^{\prime}$ has degree dividing $\ell^{\infty}(\ell-1)$. Repeating the same construction with a different rational prime $\ell^{\prime}$ in place of $\ell$ such that $\ell^{\prime}-1$ is prime to $\ell$ yields a second pair $\left(\psi_{2}, \chi_{2}\right) \in S_{v}$ and a corresponding extension $E_{\psi, \chi, \ell^{\prime}, n^{\prime}}^{\prime}$, whose degree over $E_{\psi, \chi}^{\prime}$ divides $\ell^{\prime \infty}\left(\ell^{\prime}-1\right)$, and such that

$$
E_{\psi_{2}, x_{2}}^{\prime} \subset E_{\psi, \chi, \ell^{\prime}, n^{\prime}}^{\prime} .
$$

Let

$$
E^{\prime \prime}:=E_{\psi, \chi, \ell, n}^{\prime} \cap E_{\psi, \chi, \ell^{\prime}, n^{\prime}}^{\prime} .
$$

We see then using degrees that $E^{\prime \prime} / E_{\psi, \chi}^{\prime}$ has degree dividing $(\ell-1)$, and hence $E^{\prime \prime} \subseteq E_{\psi, \chi}^{\prime} \mathbb{Q}\left(\zeta_{\ell}\right)$. Since $\ell$ is unramified in $E_{\psi, \chi}^{\prime} / \mathbb{Q}$, the extension $E^{\prime \prime} / E_{\psi, \chi}^{\prime}$ must be totally ramified at the primes above $\ell$. On the other hand, being a subextension of $E_{\psi, \chi, \ell^{\prime}, n^{\prime}}^{\prime} / E_{\psi, \chi}^{\prime}$, it is also unramified at the primes above $\ell$, and hence $E^{\prime \prime}=E_{\psi, \chi}^{\prime}$. It follows that $E_{\psi_{1}, \chi_{1}}^{\prime} \cap E_{\psi_{2}, \chi_{2}}^{\prime} \subset E_{\psi, \chi}^{\prime}$.
Thanks to Step 2, we are reduced to showing that

$$
\begin{equation*}
\bigcap_{(\psi, x) \in \bar{S}_{v}} E_{\psi, \chi}^{\prime}=E_{\nu} . \tag{3-30}
\end{equation*}
$$

The next step shows that the fields $E_{\psi, \chi}^{\prime}$ can be replaced by $E_{\psi, \chi}$ in this equality. Step 3. For all $(\psi, \chi) \in \bar{S}_{\nu}$, there exists a finite-order character $\alpha$ of $G_{K}$ such that the pair $(\psi \alpha, \chi \alpha)$ belongs to $\bar{S}_{v}$ and

$$
\begin{equation*}
E_{\psi, \chi}^{\prime} \cap E_{\psi \alpha, \chi \alpha}^{\prime} \subseteq E_{\psi, \chi} . \tag{3-31}
\end{equation*}
$$

To see this, note that the finite-order character $\chi$ has cyclic image, isomorphic to $\mathbb{Z} / n \mathbb{Z}$ say. Pick $\alpha$ such that conditions (i)-(iii) below are satisfied:
(i) $\alpha$ has order $n$ and is ramified at a single prime $\lambda$ of $K$ which lies over a rational prime $\ell$ that is split in $K$.
(ii) $\lambda$ is prime to the conductors of $\chi$ and $\chi^{*}$.
(iii) $\ell$ is unramified in $E_{\psi, \chi}^{\prime} / \mathbb{Q}$.

Conditions (i) and (ii) imply:
(iv) The field $H_{\chi \alpha} / K$ is totally ramified at $\lambda$ and unramified at $\lambda^{*}$ while $H_{\chi^{*} \alpha^{*}}$ is unramified at $\lambda$ and totally ramified at $\lambda^{*}$.

Taking

$$
L=E_{\psi, \chi}, \quad M_{1}=H_{\chi} H_{\chi^{*}}, \quad M_{2}=H_{\chi \alpha} H_{\chi^{*} \alpha^{*}}
$$

we see from (iii) and (iv) that:
(v) $L M_{1} / L$ is unramified at all primes above $\ell$.
(vi) Any subextension of $L M_{2} / L$ is ramified at some prime above $\lambda$ or $\lambda^{*}$.

Thus, $L M_{1} \cap L M_{2}=L$. But $L M_{1}=E_{\psi, \chi} H_{\chi} H_{\chi^{*}}=E_{\psi, \chi}^{\prime}$. Also, since $\alpha$ has order $n$, we have $E_{\psi^{\prime}, \chi^{\prime}}=E_{\psi, \chi}$ and

$$
L M_{2}=E_{\psi, \chi} H_{\chi \alpha} H_{\chi^{*} \alpha^{*}}=E_{\psi^{\prime}, \chi^{\prime}} H_{\chi \alpha} H_{\chi^{*} \alpha^{*}}=E_{\psi^{\prime}, \chi^{\prime}}^{\prime}
$$

so (3-31) follows.
Step 4. We are now reduced to showing

$$
\begin{equation*}
\bigcap_{(\psi, \chi) \in \bar{S}_{v}} E_{\psi, \chi}=E_{v} \tag{3-32}
\end{equation*}
$$

We will do this by showing: There exists a pair $(\psi, \chi) \in \bar{S}_{v}$ such that $E_{\psi, \chi}=E_{v}$.

We begin by choosing an ideal $\mathfrak{m}_{0}$ of $\mathbb{O}_{K}$ with the property that $\mathbb{O}_{K} / \mathfrak{m}_{0}=\mathbb{Z} / M \mathbb{Z}$ is cyclic, and an odd quadratic Dirichlet character $\varepsilon_{M}$ of conductor dividing $M$. Let $\psi_{0}$ be any Hecke character satisfying

$$
\psi_{0}((a))=\varepsilon_{M}\left(a \bmod \mathfrak{m}_{0}\right) a
$$

on principal ideals $(a)$ of $K$. Such a $\psi_{0}$ satisfies condition (1) in Definition 3.19, and therefore, after letting $\chi_{0}$ be the finite-order character satisfying

$$
v^{*}=\psi_{0}^{-1} \chi_{0} \mathrm{~N}_{K}
$$

it follows that $\left(\psi_{0}, \chi_{0}\right)$ belongs to $\bar{S}_{v}$. Furthermore, the restriction of $\psi_{0}$ to the group of principal ideals of $K$ takes values in $K$, and therefore

$$
\begin{equation*}
\chi_{0}(\sigma) \in E_{\nu}, \quad \text { for all } \sigma \in G_{H}:=\operatorname{Gal}(\bar{K} / H) \tag{3-34}
\end{equation*}
$$

The character $\psi_{0}$ itself takes values in a CM field of degree $\left[H: K\right.$ ], denoted $E_{0}$, which need not be contained in $E_{v}$ in general. To remedy this problem, let $H_{0}$ be the abelian extension of the Hilbert class field $H$ cut out by the character $\chi_{0}$. Next, let $H_{0}^{\prime}$ be any abelian extension of $K$ containing $H$ such that:
(1) There is an isomorphism $u: \operatorname{Gal}\left(H_{0}^{\prime} / K\right) \rightarrow \operatorname{Gal}\left(H_{0} / K\right)$ of abstract groups such that the diagram

commutes, where the dotted arrows indicate the isomorphisms induced by $u$ and the other arrows are the canonical maps of Galois theory.
(2) The relative discriminant of $H_{0}^{\prime}$ over $K$ is relatively prime to its conjugate (and therefore to the discriminant of $K$, in particular).

If the bottom exact sequence of groups in (3-35) is split, then the extension $H_{0}^{\prime}$ is readily produced, using class field theory. To handle the general case, we follow an approach that is suggested by the proof of Proposition 2.1.7 in [Serre 1992]. Let $\tilde{\Phi}:=\operatorname{Gal}\left(H_{0} / K\right)$ and let $\Psi: G_{K} \rightarrow \tilde{\Phi}$ be the homomorphism attached to the extension $H_{0}$. Since $H$ is everywhere unramified over $K$, the restriction $\Psi_{v}$ of $\Psi$ to a decomposition group at any prime $v$ of $K$ maps the inertia subgroup $I_{v}$ to $\mathfrak{C}:=\operatorname{Gal}\left(H_{0} / H\right)$. After viewing $\mathfrak{C}$ as a module of finite cardinality endowed with the trivial action of $G_{K}$, let $H_{S}^{1}(K, \mathfrak{C}):=\operatorname{Hom}\left(G_{K, S}, \mathfrak{C}\right)$ denote the group of homomorphisms from $G_{K}$ to $\mathfrak{C}$ which are unramified outside a given finite set $S$ of primes of $K$, and let $H_{[S]}^{1}\left(K, \mathfrak{C}^{*}\right)$ denote the dual Selmer group attached to $H_{S}^{1}(K, \mathfrak{C})$ in the sense of, for example, Theorem 2.18 of [Darmon et al. 1997]. Here $\mathfrak{C}^{*}:=\operatorname{Hom}\left(\mathfrak{C}, \boldsymbol{G}_{m}\right)$ is the Kummer dual of $\mathfrak{C}$, which is isomorphic to $\mu_{n}$ when $\mathfrak{C}=\mathbb{Z} / n \mathbb{Z}$ is cyclic of order $n$. Kummer theory (along with the nondegeneracy of the local Tate pairing) identifies $H_{[S]}^{1}\left(K, \mu_{n}\right)$ with the subgroup of $K^{\times} /\left(K^{\times}\right)^{n}$ consisting of elements $\alpha$ for which

$$
\operatorname{ord}_{v}(\alpha)=0(\bmod n) \text { for all } v, \quad \operatorname{res}_{v}(\alpha) \in\left(K_{v}^{\times}\right)^{n} \text { for all } v \in S .
$$

Let $S$ be any finite set of primes of $K$ at which $\Psi$ is unramified, satisfying the further conditions

$$
\begin{equation*}
v \in S \Rightarrow \bar{v} \notin S \quad \text { and } \quad H_{[S]}^{1}\left(K, \mathfrak{C}^{*}\right)=0 \tag{3-36}
\end{equation*}
$$

The existence of such a set $S$ follows from the statement that for any

$$
\alpha \in K^{\times}-\left(K^{\times}\right)^{n},
$$

there is a set of primes $v$ of $K$ of positive Dirichlet density for which the image of $\alpha$ in $K_{v}^{\times}$is not an $n$-th power. (This statement follows in turn from the Chebotarev density theorem applied to the extension $K\left(\mu_{n}, \alpha^{1 / n}\right)$.) Now let $T$ be any finite set of places which is disjoint from $S$. Comparing the statement of Theorem 2.18 of
[Darmon et al. 1997] in the case $M=\mathfrak{C}$ and $\mathscr{L}=S$ and $\mathscr{L}=S \cup T$ respectively, and noting that both $H_{[S]}^{1}\left(K, \mathfrak{C}^{*}\right)$ and (a fortiori) $H_{[S \cup T]}^{1}\left(K, \mathfrak{C}^{*}\right)$ are trivial, gives

$$
\frac{\# H_{S \cup T}^{1}(K, \mathfrak{C})}{\# H_{S}^{1}(K, \mathfrak{C})}=\prod_{v \in T} \frac{\# H^{1}\left(K_{v}, \mathfrak{C}\right)}{\# \mathfrak{C}}=\prod_{v} \# \operatorname{Hom}\left(I_{v}, \mathfrak{C}\right) .
$$

It follows that the rightmost arrow in the tautological exact sequence

$$
0 \rightarrow H_{S}^{1}(K, \mathfrak{C}) \rightarrow H_{S \cup T}^{1}(K, \mathfrak{C}) \rightarrow \prod_{c \in T} \operatorname{Hom}\left(I_{v}, \mathfrak{C}\right)
$$

is surjective. Letting $T$ be the set of places at which $\Psi$ is ramified, it follows that there is a homomorphism $\epsilon: G_{K} \rightarrow \mathfrak{C}$ satisfying

$$
\epsilon_{v}=\Psi_{v} \quad \text { on } I_{v}, \quad \text { for all } v \notin S .
$$

After possibly enlarging the set $S$ satisfying (3-36) and translating $\epsilon$ by a suitable homomorphism unramified outside $S$, we may further assume that the homomorphism $\Psi \epsilon^{-1}$ maps $G_{K}$ surjectively onto $\tilde{\Phi}$; the field $H_{0}^{\prime}$ can then be obtained as the fixed field of the kernel of the homomorphism $\Psi \epsilon^{-1}$.

With the extension $H_{0}^{\prime}$ in hand, let $\alpha: \operatorname{Gal}\left(H_{0}^{\prime} / K\right) \rightarrow E_{\chi}^{\times}$be the finite-order Hecke character given by

$$
\alpha(\sigma)=\chi_{0}(u(\sigma))^{-1},
$$

and set $(\psi, \chi)=\left(\psi_{0} \alpha, \chi_{0} \alpha\right)$. By construction, $(\psi, \chi)$ belongs to $\bar{S}_{\nu}$. We claim that $\chi$ and $\psi$ take values in $E_{\nu}$. Since $\nu^{*}=\psi^{-1} \chi \mathrm{~N}_{K}$, it is enough to prove this statement for $\chi$. Observe that, for all integral ideals $\mathfrak{a}$ prime to the conductors of $\chi_{0}, \chi$, and $\psi$, we have

$$
\chi(\mathfrak{a})=\chi_{0}\left(\sigma_{\mathfrak{a}}\right) / \chi_{0}\left(u\left(\sigma_{\mathfrak{a}}\right)\right)=\chi_{0}\left(\sigma_{\mathfrak{a}} u\left(\sigma_{\mathfrak{a}}\right)^{-1}\right) .
$$

But the element $\sigma_{\mathfrak{a}} u\left(\sigma_{\mathfrak{a}}\right)^{-1}$ belongs to $\operatorname{Gal}\left(H_{0} / H\right)$ by construction, and hence $\chi_{0}\left(\sigma_{\mathfrak{a}}^{-1} u\left(\sigma_{\mathfrak{a}}\right)\right)$ belongs to $E_{v}$ by (3-34). It follows that $\psi$ and $\chi$ are $E_{\nu}$-valued, and therefore $E_{\psi, \chi}=E_{\nu}$, as claimed in (3-33).

3F. Elliptic curves with complex multiplication. Theorem 2 of the Introduction admits an alternate formulation involving algebraic points on elliptic curves with complex multiplication rather than $K$-rational points on the CM abelian varieties $B_{v}$ of Theorem 2.5. The goal of this section is to describe this variant. As in the Introduction, we just write $\mathscr{L}_{p}$ for the $p$-adic $L$-function $\mathscr{L}_{p, \mathfrak{c}}$, where $\mathfrak{c}$ is the conductor of $\nu$.

We begin by reviewing the explicit construction of $B_{v}$ in terms of CM elliptic curves. The reader is referred to $\S 4$ of [Goldstein and Schappacher 1981], whose treatment we largely follow, for a more detailed exposition. Let $F$ be any abelian
extension of $K$ for which

$$
\begin{equation*}
v_{F}:=v \circ \mathrm{~N}_{F / K} \tag{3-37}
\end{equation*}
$$

becomes $K$-valued. There exists an elliptic curve $A / F$ with complex multiplication by $0_{K}$ whose associated Grossencharacter is $v_{F}$. (See Theorem 6 of [Shimura 1971] and its corollary on p. 512.) Let

$$
B:=\operatorname{Res}_{F / K}(A) .
$$

It is an abelian variety over $K$ of dimension $d:=[F: K]$. Let $G:=\operatorname{Gal}(F / K)=$ $\operatorname{Hom}_{K}(F, \overline{\mathbb{Q}})$, where the natural identification between these two sets arises from the distinguished embedding of $F$ into $\overline{\mathbb{Q}}$ that was fixed from the outset. By definition of the restriction of scalars functor, there are natural isomorphisms

$$
B_{/ F}=\prod_{\sigma \in G} A^{\sigma}, \quad B(\bar{K})=A\left(\bar{K} \otimes_{K} F\right)=\prod_{\sigma \in G} A^{\sigma}(\bar{K})
$$

of algebraic groups over $F$ and abelian groups respectively. In particular, a point of $B(\bar{K})$ is described by a $d$-tuple $\left(P_{\tau}\right)_{\tau \in G}$, with $P_{\tau} \in A^{\tau}(\bar{K})$. Relative to this identification, the Galois group $G_{K}$ acts on $B(\bar{K})$ on the left by the rule

$$
\xi\left(P_{\tau}\right)_{\tau}=\left(\xi P_{\tau}\right)_{\xi \tau}, \quad \text { for all } \xi \in G_{K} .
$$

Consider the "twisted group ring"

$$
\begin{equation*}
T:=\bigoplus_{\sigma \in G} \operatorname{Hom}_{F}\left(A, A^{\sigma}\right)=\left\{\sum_{\sigma \in G} a_{\sigma} \sigma \mid a_{\sigma} \in \operatorname{Hom}_{F}\left(A, A^{\sigma}\right)\right\} \tag{3-38}
\end{equation*}
$$

with multiplication given by

$$
\begin{equation*}
\left(a_{\sigma} \sigma\right)\left(a_{\tau} \tau\right)=a_{\sigma} a_{\tau}^{\sigma} \sigma \tau \tag{3-39}
\end{equation*}
$$

where the isogeny $a_{\tau}^{\sigma}$ belongs to $\operatorname{Hom}_{F}\left(A^{\sigma}, A^{\sigma \tau}\right)$ and the composition of isogenies in (3-39) is to be taken from left to right. The right action of $T$ on $B(\bar{K})$ defined by

$$
\begin{equation*}
\left(P_{\tau}\right)_{\tau} *\left(a_{\sigma} \sigma\right):=\left(a_{\sigma}^{\tau}\left(P_{\tau}\right)\right)_{\tau \sigma} \tag{3-40}
\end{equation*}
$$

commutes with the Galois action described in Section 3F, and corresponds to a natural inclusion $T \hookrightarrow \operatorname{End}_{K}(B)$. The $K$-algebra $E:=T \otimes_{\mathbb{Z}} \mathbb{Q}$ is isomorphic to a finite product

$$
E=\prod_{i} E_{i}
$$

of CM fields, and $\operatorname{dim}_{K}(E)=\operatorname{dim}(B)$. Therefore, the pair $(B, E)$ is a CM abelian variety in the sense of Definition 2.4. The compatible system of $\ell$-adic Galois
representations attached to $(B, E)$ corresponds to an $E$-valued algebraic Hecke character $\tilde{v}$ in the sense of Definition 2.3, satisfying the relation

$$
\begin{equation*}
\sigma_{\mathfrak{a}}(P)=P * \tilde{v}(\mathfrak{a}), \quad \text { for all } \mathfrak{a} \in I_{\mathfrak{f} \ell} \text { and } P \in B(\bar{K})_{\ell \infty}, \tag{3-41}
\end{equation*}
$$

where $\sigma_{\mathfrak{a}} \in G_{K}^{\text {ab }}$ denotes as before the Artin symbol attached to $\mathfrak{a} \in I_{\mathfrak{f} \ell}$.
The element $\tilde{v}(\mathfrak{a}) \in T$ is of the form $\varphi_{\mathfrak{a}} \sigma_{\mathfrak{a}}$, where

$$
\begin{equation*}
\varphi_{\mathfrak{a}}: A \rightarrow A^{\sigma_{\mathfrak{a}}} \tag{3-42}
\end{equation*}
$$

is an isogeny of degree Na satisfying

$$
\begin{equation*}
\varphi_{\mathfrak{a}}(P)=P^{\sigma_{\mathfrak{a}}}, \tag{3-43}
\end{equation*}
$$

for any $P \in A[\mathfrak{g}]$ with $(\mathfrak{g}, \mathfrak{a})=1$. Note that the isogenies $\varphi_{\mathfrak{a}}$ satisfy the following cocycle condition:

$$
\begin{equation*}
\varphi_{\mathfrak{a b}}=\varphi_{\mathfrak{b}}^{\sigma_{\mathfrak{a}}} \circ \varphi_{\mathfrak{a}} \tag{3-44}
\end{equation*}
$$

The following proposition relates the Hecke characters $\tilde{v}$ and $\nu$.
Proposition 3.30. Given any homomorphism $j \in \operatorname{Hom}_{K}(E, \mathbb{C})$, let $v_{j}:=j \circ \tilde{v}$ be the corresponding $\mathbb{C}$-valued Hecke character of $K$ of infinity type $(1,0)$. The assignment $j \mapsto v_{j}$ gives a bijection from $\operatorname{Hom}_{K}(E, \mathbb{C})$ to the set $\Sigma_{v, F}$ of Hecke characters $v^{\prime}$ of $K$ (of infinity type $(1,0)$ ) satisfying

$$
v^{\prime} \circ \mathbf{N}_{F / K}=v \circ \mathbf{N}_{F / K} .
$$

Proposition 3.30 implies that there is a unique homomorphism $j_{\nu} \in \operatorname{Hom}_{K}(E, \mathbb{C})$ satisfying $j_{\nu} \circ \tilde{\nu}=\nu$. In particular, $j_{\nu}$ maps $E$ to $E_{\nu}$ and $T$ to a finite-index subring of $T_{v}$. The abelian variety $B_{v}$ attached to $v$ in Theorem 2.5 can now be defined as the quotient $B \otimes_{T, j_{v}} T_{\nu}$. In subsequent constructions, it turns out to be more useful to realize $B_{v}$ as a subvariety of $B$, which can be done by setting

$$
\begin{equation*}
B_{v}:=B\left[\operatorname{ker} j_{\nu}\right] . \tag{3-45}
\end{equation*}
$$

The natural action of $T$ on $B_{v}$ factors through the quotient $T / \operatorname{ker}\left(j_{v}\right)$, an integral domain having $E_{\nu}$ as field of fractions.

Consider the inclusion

$$
\begin{equation*}
i_{v}: B_{v}(K) \hookrightarrow B(K)=A(F), \tag{3-46}
\end{equation*}
$$

where the last identification arises from the functorial property of the restriction of scalars. The following proposition gives an explicit description of the image of $\left(B_{\nu}(K) \otimes E_{\nu}\right)^{T_{v}}$ in $A(F) \otimes_{O_{K}} E_{v}$ under the inclusion $i_{v}$ obtained from (3-46).

Proposition 3.31. Let $\tilde{E}$ be any field containing $E_{\nu}$. The inclusion $i_{v}$ of (3-46) identifies $\left(B_{v}(K) \otimes \tilde{E}\right)^{T_{v}}$ with
$\left(A(F) \otimes_{\odot_{K}} \tilde{E}\right)^{v}:=\left\{P \in A(F) \otimes_{\odot_{K}} \tilde{E}\right.$ such that $\varphi_{\mathfrak{a}}(P)=v(\mathfrak{a}) P^{\sigma_{\mathfrak{a}}}$, for all $\left.\mathfrak{a} \in I_{\mathfrak{f}}\right\}$.
Proof. It follows from the definitions that $B(K)$ is identified with the set of $\left(P_{\tau}\right)$, with $P_{\tau} \in A^{\tau}(\bar{K})$ satisfying

$$
\begin{equation*}
\xi P_{\tau}=P_{\xi \tau}, \quad \text { for all } \xi \in G_{K} . \tag{3-47}
\end{equation*}
$$

Furthermore, if such a $\left(P_{\tau}\right)$ belongs to $\left(B_{v}(K) \otimes E_{\nu}\right)^{T_{v}}$, then after setting $\tilde{v}(\mathfrak{a})=$ $\varphi_{\mathfrak{a}} \sigma_{\mathfrak{a}}$ as in (3-42), we also have

$$
\begin{equation*}
\left(\varphi_{\mathfrak{a}}^{\tau}\left(P_{\tau}\right)\right)_{\tau \sigma_{\mathfrak{a}}}=\left(P_{\tau}\right)_{\tau} * \tilde{v}(\mathfrak{a})=\left(\nu(\mathfrak{a}) P_{\tau}\right)_{\tau} . \tag{3-48}
\end{equation*}
$$

Equating the $\sigma_{\mathfrak{a}}$-components of these two vectors gives

$$
\varphi_{\mathfrak{a}}\left(P_{1}\right)=v(\mathfrak{a}) P_{\sigma_{\mathfrak{a}}}=v(\mathfrak{a}) \sigma_{\mathfrak{a}} P_{1},
$$

where 1 is the identity embedding of $F$ and the last equality follows from (3-47). The proposition follows directly from this, after noting that the identification of $B(K)$ with $A(F)$ is simply the one sending $\left(P_{\tau}\right)_{\tau}$ to $P_{1}$.

Given a global field $F$ as in (3-37), let $F_{\nu}$ denote the subfield of $\overline{\mathbb{Q}}$ generated by $F$ and $E_{\nu}$. Recall that $\omega_{A} \in \Omega^{1}(A / F)$ is a nonzero differential and that $\Omega_{p}(A)$ is the associated $p$-adic period.

Theorem 3.32. There exists a point $P_{A, \nu} \in\left(A(F) \otimes_{\Theta_{K}} E_{\nu}\right)^{v}$ such that

$$
\mathscr{L}_{p}\left(v^{*}\right)=\Omega_{p}(A)^{-1} \log _{\omega_{A}}^{2}\left(P_{A, v}\right)\left(\bmod F_{v}^{\times}\right) .
$$

The point $P_{A, \nu}$ is nonzero if and only if $L^{\prime}(\nu, 1) \neq 0$.
Proof. Theorem 2 of the Introduction asserts that

$$
\begin{equation*}
\mathscr{L}_{p}\left(v^{*}\right)=\Omega_{p}\left(v^{*}\right)^{-1} \log _{\omega_{v}}^{2}\left(P_{v}\right), \tag{3-49}
\end{equation*}
$$

for some point $P_{v} \in B_{v}(K) \otimes \mathbb{Q}$ which is nontrivial if and only if $L^{\prime}(\nu, 1) \neq 0$. By Lemma 2.14, we find

$$
\begin{equation*}
\Omega_{p}\left(v^{*}\right)^{-1}=\Omega_{p}(A)^{-1}\left(\bmod F_{v}^{\times}\right) . \tag{3-50}
\end{equation*}
$$

Also, by Proposition 3.31, we can view $P_{\nu}$ as a point $P_{A, v} \in\left(A(F) \otimes_{\odot_{K}} E_{\nu}\right)^{\nu}$, and we have

$$
\begin{equation*}
\log _{\omega_{v}}\left(P_{\nu}\right)=\log _{\omega_{A}}\left(P_{A, \nu}\right)\left(\bmod F_{v}^{\times}\right) . \tag{3-51}
\end{equation*}
$$

Theorem 3.35 now follows by rewriting (3-49) using (3-50) and (3-51).

3G. A special case. This section is devoted to a more detailed and precise treatment of Theorem 3.32 under the following special assumptions:
(1) The quadratic imaginary field $K$ has class number one, odd discriminant, and unit group of order two. This implies that $K=\mathbb{Q}(\sqrt{-D})$, where $D:=$ $-\operatorname{Disc}(K)$ belongs to the finite set

$$
S:=\{7,11,19,43,67,163\} .
$$

(2) $\psi_{0}$ is the Hecke character of $K$ of infinity type $(1,0)$ given by the formula

$$
\begin{equation*}
\psi_{0}((a))=\varepsilon_{K}\left(a \bmod \mathfrak{d}_{K}\right) a \tag{3-52}
\end{equation*}
$$

The character $\psi_{0}$ determines (uniquely, up to an isogeny) an elliptic curve $A / \mathbb{Q}$ satisfying

$$
\operatorname{End}_{K}(A)=\mathfrak{O}_{K}, \quad L(A / \mathbb{Q}, s)=L\left(\psi_{0}, s\right)
$$

After fixing $A$, we will also write $\psi_{A}$ instead of $\psi_{0}$. It can be checked that the conductor of $\psi_{A}$ is equal to $\mathfrak{d}_{K}$, and that

$$
\psi_{A}^{*}=\bar{\psi}_{A}, \quad \psi_{A} \psi_{A}^{*}=\mathrm{N}_{K}, \quad \varepsilon_{\psi_{A}}=\varepsilon_{K}
$$

Remark 3.33. The rather stringent assumptions on $K$ that we have imposed exclude the arithmetically interesting, but somewhat idiosyncratic, cases where $K=\mathbb{Q}(\sqrt{-3}), \mathbb{Q}(i)$, and $\mathbb{Q}(\sqrt{-2})$.

With the above assumptions, the character $\psi_{A}$ can be used to give an explicit description of the set $\Sigma_{\mathrm{sd}}\left(c \mathfrak{d}_{K}\right)$ :

Lemma 3.34. Let $c$ be an integer prime to $D$, and let $v$ be a Hecke character in $\Sigma_{\mathrm{sd}}\left(c \mathfrak{d}_{K}\right)$. Then $v$ is of the form

$$
v=\psi_{A} \chi^{-1}
$$

where $\chi$ is a finite-order ring class character of $K$ of conductor $c$.
Proof. The fact that $v$ and $\psi_{A}$ both have central character $\varepsilon_{K}$ implies that $\chi$ is a ring class character that is unramified at $\mathfrak{d}_{K}$, and hence has conductor exactly $c$.

Given a ring class character $\chi$ of conductor $c$ as above with values in a field $E_{\chi}$, let

$$
\begin{align*}
& \left(A\left(H_{c}\right) \otimes_{\mathfrak{O}_{K}} E_{\chi}\right)^{\chi}  \tag{3-53}\\
& :=\left\{P \in A\left(H_{c}\right) \otimes_{\mathfrak{O}_{K}} E_{\chi} \text { such that } \sigma P=\chi(\sigma) P, \text { for all } \sigma \in \operatorname{Gal}\left(H_{c} / K\right)\right\}
\end{align*}
$$

Finally, choose a nonzero differential $\omega_{A} \in \Omega^{1}(A / K)$, and write $\Omega_{p}(A)$ for the $p$-adic period attached to this choice as in Section 3A. Since $A=B_{\psi_{0}}$ is the abelian variety attached to $\psi_{0}$, it follows that $\Omega_{p}\left(\psi_{A}^{*}\right)=\Omega_{p}(A)$.

The following theorem is a more precise variant of Theorem 3.32.
Theorem 3.35. Let $\chi$ be a ring class character of $K$ of conductor prime to $\mathfrak{d}_{K}$. Then there exists a point $P_{A}(\chi) \in\left(A\left(H_{\chi}\right) \otimes_{\varrho_{K}} E_{\chi}\right)^{\chi}$ such that

$$
\mathscr{L}_{p}\left(\psi_{A}^{*} \chi\right)=\Omega_{p}(A)^{-1} \mathfrak{g}(\chi) \log _{\omega_{A}}^{2}\left(P_{A}(\chi)\right)\left(\bmod E_{\chi}^{\times}\right) .
$$

The point $P_{A}(\chi)$ is nonzero if and only if $L^{\prime}\left(\psi_{A} \chi^{-1}, 1\right) \neq 0$.
Proof. By Theorem 2 of the Introduction,

$$
\begin{equation*}
\mathscr{L}_{p}\left(\psi_{A}^{*} \chi\right)=\mathscr{L}_{p}\left(v^{*}\right)=\Omega_{p}\left(v^{*}\right)^{-1} \log _{\omega_{v}}^{2}\left(P_{v}\right)\left(\bmod E_{v}^{\times}\right), \tag{3-54}
\end{equation*}
$$

for some point $P_{v} \in B_{v}(K) \otimes \mathbb{Q}$ which is nontrivial if and only if $L^{\prime}\left(\psi_{A} \chi^{-1}, 1\right) \neq 0$. Since $\chi^{*-1}=\chi$ and $E_{\nu}=E_{\chi}$, we find from Lemma 2.14 that

$$
\begin{equation*}
\Omega_{p}\left(v^{*}\right)^{-1}=\Omega_{p}\left(\psi_{A}^{*} \chi^{*-1}\right)^{-1}=\Omega_{p}(A)^{-1} \mathfrak{g}(\chi)^{-1}\left(\bmod E_{\chi}^{\times}\right) \tag{3-55}
\end{equation*}
$$

After noting that, as in (2-7), $\left.B_{v}=B_{\psi, \chi}=\left(A \otimes_{\Theta_{K}} T_{\chi}\right)\right)^{-1}$ as abelian varieties over $K$, we observe that $\omega_{\nu}=\omega_{\psi, \chi}$ and that the point $P_{\nu} \in B_{\nu}(K)$ can be written as

$$
P_{\nu}=\sum_{\sigma \in G} P^{\sigma} \otimes \chi^{-1}(\sigma),
$$

for some $P \in A\left(H_{c}\right) \otimes \mathbb{Q}$. Letting $P_{A, \chi}$ be the corresponding element in

$$
A\left(H_{c}\right) \otimes_{\varrho_{K}} E_{\chi}
$$

given by

$$
P_{A, \chi}=\sum_{\sigma \in G} \chi^{-1}(\sigma) P^{\sigma}
$$

we have

$$
\begin{equation*}
\log _{\omega_{\nu}}^{2}\left(P_{\nu}\right)=\log _{\omega_{\nu, \chi}}^{2}\left(P_{\nu}\right)=\mathfrak{g}(\chi)^{2} \log _{\omega_{\psi, \chi}^{0}}^{2}\left(P_{\nu}\right)=\mathfrak{g}(\chi)^{2} \log _{\omega_{A}}^{2}\left(P_{A, \chi}\right)\left(\bmod E_{\chi}^{\times}\right) \tag{3-56}
\end{equation*}
$$

where the second equality follows from Lemma 2.8 and the last from Lemma 3.23. Theorem 3.35 now follows by rewriting (3-54) using (3-55) and (3-56).

In the special case where $\chi$ is a quadratic ring class character of $K$, cutting out an extension $L=K(\sqrt{a})$ of $K$, we obtain

$$
\begin{equation*}
\mathscr{L}_{p}\left(\psi_{A}^{*} \chi\right)=\Omega_{p}(A)^{-1} \sqrt{a} \log _{\omega_{A}}^{2}\left(P_{A, L}^{-}\right)\left(\bmod K^{\times}\right), \tag{3-57}
\end{equation*}
$$

where $P_{A, L}^{-}$is a $K$-vector space generator of the trace 0 elements in $A(L) \otimes \mathbb{Q}$. Since in this case $\psi_{A} \chi$ is the Hecke character attached to a CM elliptic curve over $\mathbb{Q}$, from (3-57) one recovers Rubin's Theorem 1 of the Introduction.

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