

Intersecting geodesics on the modular surface
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#### Abstract

We introduce the modular intersection kernel, and we use it to study how geodesics intersect on the full modular surface $\mathbb{X}=\mathrm{PSL}_{2}(\mathbb{Z}) \backslash \mathbb{H}$. Let $C_{d}$ be the union of closed geodesics with discriminant $d$ and let $\beta \subset \mathbb{X}$ be a compact geodesic segment. As an application of Duke's theorem to the modular intersection kernel, we prove that $\left\{\left(p, \theta_{p}\right): p \in \beta \cap C_{d}\right\}$ becomes equidistributed with respect to $\sin \theta d s d \theta$ on $\beta \times[0, \pi]$ with a power saving rate as $d \rightarrow+\infty$. Here $\theta_{p}$ is the angle of intersection between $\beta$ and $C_{d}$ at $p$. This settles the main conjectures introduced by Rickards(2021).

We prove a similar result for the distribution of angles of intersections between $C_{d_{1}}$ and $C_{d_{2}}$ with a power-saving rate in $d_{1}$ and $d_{2}$ as $d_{1}+d_{2} \rightarrow \infty$. Previous works on the corresponding problem for compact surfaces do not apply to $\mathbb{X}$, because of the singular behavior of the modular intersection kernel near the cusp. We analyze the singular behavior of the modular intersection kernel by approximating it by general (not necessarily spherical) point-pair invariants on $\operatorname{PSL}_{2}(\mathbb{Z}) \backslash \mathrm{PSL}_{2}(\mathbb{R})$ and then by studying their full spectral expansion.


## 1. Introduction

Let $Y$ be a negatively curved surface of finite area. The prime geodesic theorem [Sarnak 1980] states that the number of primitive closed geodesics having length less than $L$, which we denote by $\pi(L)$, satisfies

$$
\pi(L) \sim \frac{e^{L}}{L}
$$

as $L \rightarrow \infty$. A natural problem is to understand how primitive closed geodesics of length less than $L$ are positioned or distributed in $Y$ as $L \rightarrow \infty$. In particular, one may ask
(1) how the number of transversal intersections $I\left(\alpha_{1}, \alpha_{2}\right)$ between two primitive closed geodesics $\alpha_{1}$ and $\alpha_{2}$ is distributed, or
(2) how the set of angles of intersections between $\alpha_{1}$ and $\alpha_{2}$ is distributed,

[^0]as one varies $\alpha_{1}$, or both $\alpha_{1}$ and $\alpha_{2}$ ? Bonahon [1986] defined the intersection form $i: \mathcal{C} \times \mathcal{C} \rightarrow \mathbb{R}^{+}$ on the space of currents $\mathcal{C}$ such that when $\mu_{i}(i=1,2)$ is the unique invariant measure corresponding to $\alpha_{i}$, then $i\left(\mu_{1}, \mu_{2}\right)=I\left(\alpha_{1}, \alpha_{2}\right)$. When $Y$ is compact, Pollicott and Sharp [2006] used an extension of the intersection form to understand the distribution of angles of self-intersections of closed geodesic $\alpha$ having length less than $L$, as $L \rightarrow \infty$. When $Y$ is a compact hyperbolic surface, using the intersection form, Herrera Jaramillo [2015] proved that the distribution of $I\left(\alpha_{1}, \alpha_{2}\right) /\left(l\left(\alpha_{1}\right) l\left(\alpha_{2}\right)\right)$ for closed geodesics $\alpha_{1}, \alpha_{2}$ of length $<L$, is concentrated near $1 /\left(2 \pi^{2}(g-1)\right)=2 /(\pi \operatorname{vol}(Y))$ with exponentially decaying tail, as $L \rightarrow \infty$. Here $l(\cdot)$ is the length function, and $g$ is the genus of $Y$.

In this article, we study a refined problem:
(3) How are the locations and angles of intersections between $\alpha_{1}$ and $\alpha_{2}$ jointly distributed relative to $\alpha_{2}$, as one varies $\alpha_{1}$, or both $\alpha_{1}$ and $\alpha_{2}$ ?

Let $\mathbb{X}=\mathrm{PSL}_{2}(\mathbb{Z}) \backslash \mathbb{H}$ be the full modular surface. The connection between the geometry of geodesics on $\mathbb{X}$ and number theory has a rich history. Artin [1924] discovered a relation between geometry of geodesics in $\mathbb{X}$ and continued fraction expansion. As a result, he proved that there is a hyperbolic geodesic in $\mathbb{X}$ that comes arbitrarily close to any given hyperbolic segment in $\mathbb{X}$. So this geodesic is not only dense, but dense in all directions simultaneously. Another deep connection is discovered in the spectacular work of Katok [1985]. She showed that certain holomorphic Poincaré series (introduced by Petersson) associated with closed geodesics on a Fuchsian group of the first kind, span the corresponding space of cusp forms. Moreover, she proved a formula relating the intersection angles between pairs of closed geodesics to the periods of these holomorphic Poincaré series.

On $\mathbb{X}$, primitive oriented closed geodesics are in one-to-one correspondence with conjugacy classes of primitive hyperbolic elements of $\operatorname{PSL}_{2}(\mathbb{Z})$. Moreover there is a bijection between the primitive hyperbolic conjugacy classes and the $\mathrm{SL}_{2}(\mathbb{Z})$ equivalence classes of primitive integral binary quadratic forms of nonsquare positive discriminant [Luo et al. 2009; Sarnak 1982]. So by the discriminant of a primitive closed geodesic, we mean the discriminant of the corresponding binary quadratic form. In particular, if the hyperbolic class $\gamma$ is associated to the binary quadratic form $Q$ then $\gamma^{-1}$ is associated to $-Q$.

Let $\left(x_{d}, y_{d}\right)$ be the fundamental solution of Pell's equation $x^{2}-d y^{2}=4$, and let $\varepsilon_{d}:=\frac{1}{2}\left(x_{d}+\sqrt{d} y_{d}\right)>1$. Each oriented primitive closed geodesics of discriminant $d$ has a unique lift to a closed geodesic of length $2 \log \varepsilon_{d}$ in the unit tangent bundle $S \mathbb{X}$. Let $h(d)$ be the number of inequivalent primitive integral binary quadratic forms of discriminant $d$. We denote the disjoint union of these $h(d)$ closed geodesics by $\mathscr{C}_{d} \subset S \mathbb{}$, which has total length $2 h(d) \log \varepsilon_{d}$.

Note that the closed geodesic on $\mathbb{X}$ has length $\log \varepsilon_{d}$ or $2 \log \varepsilon_{d}$ according as $Q$ is or is not equivalent to $-Q$ [Duke 1988, page 75]. We now let $C_{d}$ be the union of primitive (unoriented) closed geodesics of discriminant $d$ on $\mathbb{X}$, and note that $l\left(C_{d}\right)=h(d) \log \varepsilon_{d}$ is the total length of $C_{d}$.

Theorem 1.1. Fix $T>100$, and let $\beta$ be a compact oriented geodesic segment of length $<1$ in the region determined by $y<T$ on $\mathbb{X}$. For $0<\theta_{1}<\theta_{2}<\pi$, let $I_{\theta_{1}, \theta_{2}}\left(\beta, C_{d}\right)$ be the number of intersections between
$\beta$ and $C_{d}$ with the angle between $\theta_{1}$ and $\theta_{2}$. (Here the angle between $\beta$ and $C_{d}$ at $p \in \beta \cap C_{d}$ is measured counterclockwise from the tangent to $\beta$ at $p$ to the tangent to $C_{d}$ at $p$.)

Then we have

$$
\frac{I_{\theta_{1}, \theta_{2}}\left(\beta, C_{d}\right)}{l(\beta) l\left(C_{d}\right)}=\frac{3}{\pi^{2}} \int_{\theta_{1}}^{\theta_{2}} \sin \theta d \theta+O_{\epsilon}\left(d^{-25 / 3584+\epsilon}\right)
$$

uniformly in $\beta, \theta_{1}$, and $\theta_{2}$, under the assumption that

$$
\theta_{2}-\theta_{1} \gg d^{-25 / 7168}
$$

and that

$$
l(\beta) \gg d^{-25 / 7168}
$$

(Here and elsewhere, $A \ll_{\tau} B$ means $|A| \leq C(\tau) B$ for some constant $C(\tau)$ that depends only on $\tau$.)
Remark 1.1. This statement is false if $C_{d}$ is replaced by individual geodesics. For instance, the set of intersections between $\beta$ and a closed geodesic $\alpha$ does not necessarily become equidistributed as $l(\alpha) \rightarrow \infty$. To see this, take a finite sheeted covering $S$ of $\mathbb{X}$ whose genus is $\geq 2$. Then according to Rivin's work [2001] there are arbitrarily long simple closed geodesics on $S$. Note that these simple closed geodesics must be contained in a compact part of $S$ [Jung and Reid 2021]. This implies that there is a compact set $C \subset \mathbb{X}$ which contains arbitrarily long primitive closed geodesics. Take a geodesic segment $\beta$ in $\mathbb{X}-C$. Then there are infinitely many closed geodesics which do not intersect $\beta$.

Remark 1.2. The exponent $-\frac{25}{3584}$ can be improved slightly by refining our argument (for instance, by inputting the Weyl-like subconvex bound [Petrow and Young 2019] instead of the Burgess-like subconvex bound [Heath-Brown 1980]), but in order to keep the exposition simple, we do not discuss the optimal rate in the current article.

As an immediate consequence, we deduce that the intersection points and corresponding angles become equidistributed, resolving the main conjectures introduced by Rickards [2021].

Corollary 1.2. Fix a closed geodesic $\alpha$. Then for any fixed segment $\beta \subset \alpha$, and any fixed $0<\theta_{1}<\theta_{2}<\pi$, we have

$$
\lim _{d \rightarrow \infty} \frac{I_{\theta_{1}, \theta_{2}}\left(\beta, C_{d}\right)}{I\left(\alpha, C_{d}\right)}=\frac{l(\beta)}{l(\alpha)} \int_{\theta_{1}}^{\theta_{2}} \frac{\sin \theta}{2} d \theta
$$

Remark 1.3. Rickards's work is motivated by the work of Darmon and Vonk [2022] on the arithmetic ( $p$-arithmetic) intersection between pairs of oriented closed geodesics on the modular surfaces (Shimura curves). The arithmetic intersection between oriented closed geodesics $\alpha_{1}$ and $\alpha_{2}$ of discriminants $D_{1}$ and $D_{2}$ only depends on $D_{1}$ and $D_{2}$ and the angles of intersections between $\alpha_{1}$ and $\alpha_{2}$. Darmon and Vonk [2022, Conjecture 2] conjectured that the $p$-arithmetic intersection is algebraic and belongs to the composition of the Hilbert class field of real quadratic fields of discriminants $D_{1}$ and $D_{2}$.

To prove our main results, we introduce the modular intersection kernel. For $\delta>0$ and $\theta_{1}, \theta_{2} \in(0, \pi)$, let $k_{\delta}^{\theta_{1}, \theta_{2}}: S \sharp \times S \sharp \rightarrow \mathbb{R}$ be the integral kernel defined by

$$
k_{\delta}^{\theta_{1}, \theta_{2}}\left(\left(x_{1}, \xi_{1}\right),\left(x_{2}, \xi_{2}\right)\right)=1,
$$

if the geodesic segments of length $\delta$ from $x_{i}$ with the initial vector $\xi_{i}$ intersect at an angle $\in\left(\theta_{1}, \theta_{2}\right)$, and 0 otherwise. Under the identification $S \sharp \cong \mathrm{PSL}_{2}(\mathbb{R})$, for a given discrete subgroup $\Gamma \subset \mathrm{PSL}_{2}(\mathbb{R})$, we define the modular intersection kernel $K_{\delta}^{\theta_{1}, \theta_{2}}: \Gamma \backslash \mathrm{PSL}_{2}(\mathbb{R}) \times \Gamma \backslash \mathrm{PSL}_{2}(\mathbb{R}) \rightarrow \mathbb{R}$ by taking the average of $k_{\delta}^{\theta_{1}, \theta_{2}}$ over $\Gamma$ :

$$
K_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{1}, g_{2}\right)=\sum_{\gamma \in \Gamma} k_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{1}, \gamma g_{2}\right)
$$

The basic idea of the proof of Theorem 1.1 then is as follows. Heuristically,

$$
I_{\theta_{1}, \theta_{2}}\left(\beta, C_{d}\right)
$$

should be well approximated by

$$
\begin{equation*}
\frac{1}{2 \delta^{2}} \int_{\mathscr{C}_{d}} \int_{\tilde{\beta}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} \tag{1-1}
\end{equation*}
$$

where $\tilde{\beta} \subset S \mathbb{X}$ is a lift of $\beta$ with either of orientations of $\beta$

$$
\tilde{\beta}(t)=\left(\beta(t), \beta^{\prime}(t)\right),
$$

under assuming that $\beta(t)$ is parametrized by the arc length. As noted in [Luo et al. 2009], Duke's theorem [1988] can be extended to the equidistribution of $\mathscr{C}_{d}$ in $\mathrm{PSL}_{2}(\mathbb{Z}) \backslash \mathrm{PSL}_{2}(\mathbb{R})$ as $d \rightarrow \infty$. Observing that

$$
\begin{equation*}
\frac{1}{2 \delta^{2}} \int_{\tilde{\beta}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, g\right) d s_{1} \tag{1-2}
\end{equation*}
$$

is a compactly supported function in $g$ for compact $\beta,(1-1)$ is

$$
\sim \frac{l\left(\mathscr{C}_{d}\right)}{2 \delta^{2}} \int_{g} \int_{\tilde{\beta}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, g\right) d s_{1} d \mu_{g}
$$

which is asymptotically $\left(3 / \pi^{2}\right) l\left(C_{d}\right) l(\beta) \int_{\theta_{1}}^{\theta_{2}} \sin \alpha d \alpha$ as $\delta \rightarrow 0$, by an explicit computation.
Note that (1-2) is a discontinuous function. Therefore, in order to obtain the rate of convergence, we need a smooth approximation of (1-2), and a quantified version of Duke's theorem with explicit dependency on the test functions. To this end, we follow the argument sketched in [Luo et al. 2009] to prove:

Theorem 1.3. Assume that $f \in C_{0}^{\infty}\left(\operatorname{PSL}_{2}(\mathbb{Z}) \backslash \operatorname{PSL}_{2}(\mathbb{R})\right)$ has support in the region determined by $y<T$.
Then we have

$$
\frac{1}{l\left(\mathscr{C}_{d}\right)} \int_{\mathscr{C}_{d}} f(s) d s=\frac{3}{\pi^{2}} \int_{\mathrm{PSL}_{2}(\mathbb{Z}) \backslash \operatorname{PSL}_{2}(\mathbb{R})} f(g) d \mu_{g}+O_{\epsilon}\left(\log T d^{-25 / 512+\epsilon}\|f\|_{W^{6, \infty}}\right)
$$

Here $\|\cdot\|_{W^{k, p}}$ is the Sobolev norm:

$$
\|f\|_{W^{k, p}}=\max _{|\alpha| \leq k}\left\|\partial_{\theta}^{\alpha_{1}}\left(y \partial_{x}\right)^{\alpha_{2}}\left(y \partial_{y}\right)^{\alpha_{3}} f\right\|_{L^{p}} .
$$

Remark 1.4. The proof of Theorem 1.1 is based on the equidistribution of the lifts of $C_{d}$ in the unit tangent bundle. For this reason, one may generalize Theorem 1.1 to any surfaces and any sequence of closed geodesics whose lifts become equidistributed on the unit tangent bundle.

1A. Intersecting two closed geodesics. We now consider the number of intersections between two closed geodesics when both vary.

Theorem 1.4. The following estimate holds uniformly in $d_{1}, d_{2}>0$, and $0<\theta_{1}<\theta_{2}<\pi$ such that $\theta_{2}-\theta_{1} \gg\left(d_{1} d_{2}\right)^{-25 / 3072}$

$$
\frac{I_{\theta_{1}, \theta_{2}}\left(C_{d_{1}}, C_{d_{2}}\right)}{l\left(C_{d_{1}}\right) l\left(C_{d_{2}}\right)}=\frac{3}{\pi^{2}} \int_{\theta_{1}}^{\theta_{2}} \sin \theta d \theta+O_{\epsilon}\left(\left(d_{1} d_{2}\right)^{-25 / 6144+\epsilon}\right)
$$

Note that if $\Gamma$ is cocompact, then the modular intersection kernel coincides with the intersection kernel from [Lalley 2014] when $\theta=\pi$ and $\delta>0$ is sufficiently small. However, when $\Gamma \backslash \mathbb{H}$ is noncompact, then they are never the same; for instance, we have $K_{\delta}^{\theta_{1}, \theta_{2}}(g, g)=\Omega(y)$ as $y \rightarrow \infty$ (Proposition 2.2). In particular, $K_{\delta}^{\theta_{1}, \theta_{2}}$ is not a Hilbert-Schmidt kernel, so the usual spectral theory does not apply. This is the main technical difficulty of dealing with the modular intersection kernel for noncompact quotients of $\mathbb{H}$. As it will be shown in the subsequent chapters, when both $\alpha_{1}$ and $\alpha_{2}$ are closed geodesics, $I_{\theta_{1}, \theta_{2}}\left(\alpha_{1}, \alpha_{2}\right) /\left(l\left(\alpha_{1}\right) l\left(\alpha_{2}\right)\right)$ is the integral of $\delta^{-2} K_{\delta}^{\theta_{1}, \theta_{2}} /\left(l\left(\alpha_{1}\right) l\left(\alpha_{2}\right)\right)$ over $\alpha_{1} \times \alpha_{2}$. When $\alpha_{1}$ and $\alpha_{2}$ vary over closed geodesics of length $<L$, as $L \rightarrow \infty$, we expect that the integral converges to the integral of $\delta^{-2} K_{\delta}^{\theta_{1}, \theta_{2}}$ over $\Gamma \backslash S \sharp \times \Gamma \backslash S \sharp \mathbb{H}$, since $\alpha_{1} \times \alpha_{2}$ becomes equidistributed in $\Gamma \backslash S \sharp \times \Gamma \backslash S \sharp$, as $L \rightarrow \infty$. However, unboundedness of the modular intersection kernel $K$ causes issues of interchanging the limit and the integral. In particular, the argument of [Pollicott and Sharp 2006] using intersection form does not apply in this case. Hence, in order to prove Theorem 1.4, we study the full spectral expansion of $K_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{1}, g_{2}\right)$. This is similar to the existing work on the weight $m$ Selberg's trace formula [Hejhal 1976], except that we have to deal with all weights simultaneously, and that the modular intersection kernel is not diagonalizable in general. We go over this carefully in Section 5 . Once the spectral expansion is obtained, the integral of $\delta^{-2} K_{\delta}^{\theta_{1}, \theta_{2}}$ over $\alpha_{1} \times \alpha_{2}$ becomes a linear combination of the period integrals of the form

$$
\int_{\alpha_{1}} \phi_{1} d s \times \int_{\alpha_{2}} \phi_{2} d s
$$

We may now use the same estimates that we use in order to prove the effective Duke's theorem to bound these, which leads to Theorem 1.4, generalizing [Pollicott and Sharp 2006] to a noncompact hyperbolic surface.

## 2. The modular intersection kernel

2A. Parametrization. Recall that $\operatorname{PSL}_{2}(\mathbb{R})$ acts transitively on $\mathbb{H}$ and on $S \sharp$ with the fractional transformations. For $g \in \operatorname{PSL}_{2}(\mathbb{R}), z \in \mathbb{H}$ and $u \in S \sharp$ we write these actions by $g z$ and $g u$. We parametrize the points of $\mathbb{H}$ and $S \Vdash$ with $x+i y$ and $(x+i y, \exp (i \theta))$. Let

$$
\Pi((x+i y, \exp (i \theta))):=x+i y
$$

be the projection map from $S$ 乌 to $\mathbb{H}$.
Fix $z_{0}=i$ and $u_{0}=(i, \exp (i \pi / 2))$. Let $g=n a R_{\theta} \in \operatorname{PSL}_{2}(\mathbb{R})$ be the Iwasawa decomposition where

$$
n=n(x)=\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right), \quad a=a(y)=\left(\begin{array}{cc}
y^{1 / 2} & 0 \\
0 & y^{-1 / 2}
\end{array}\right), \quad \text { and } \quad R_{\theta}=\left(\begin{array}{cc}
\cos \theta & -\sin \theta \\
\sin \theta & \cos \theta
\end{array}\right)
$$

Then we have

$$
g z_{0}=x+i y \quad \text { and } \quad g u_{0}=\left(x+i y, \exp \left(i\left(\frac{\pi}{2}+2 \theta\right)\right)\right)
$$

For the rest of the paper, we identify $S \llbracket$ with $\operatorname{PSL}_{2}(\mathbb{R})$ by sending $g \in \mathrm{PSL}_{2}(\mathbb{R})$ to $g u_{0}$. We often use the following fact in our computation without mentioning.

Proposition 2.1. The image under $\gamma \in \mathrm{SL}_{2}(\mathbb{R})$ of the geodesic segment of length $\delta$ corresponding to $g=(x, \xi)$ is the geodesic segment of length $\delta$ corresponding to $\gamma g$.

We use the volume form given by $d V=(d x d y d \theta) / y^{2}$. The volume of $S \mathbb{X}$ is then $\pi^{2} / 3$.
2B. Preliminary estimates. We first recall here the definition of the modular intersection kernel described in the introduction. For $\delta>0$ and $\theta_{1}, \theta_{2} \in(0, \pi)$, we define the integral kernel

$$
k_{\delta}^{\theta_{1}, \theta_{2}}: S \mathbb{H} \times S \mathbb{H} \rightarrow \mathbb{R}
$$

by

$$
k_{\delta}^{\theta_{1}, \theta_{2}}\left(\left(x_{1}, \xi_{1}\right),\left(x_{2}, \xi_{2}\right)\right)=1,
$$

if the geodesic segment of length $\delta$ on $\mathbb{H}$ from $x_{1}$ with the initial vector $\xi_{1}$ and the segment from $x_{2}$ with the initial vector $\xi_{2}$ intersect at an angle $\in\left(\theta_{1}, \theta_{2}\right)$, and 0 otherwise. Here the angle of the intersection of geodesic segments $l_{1}$ and $l_{2}$ at $p \in l_{1} \cap l_{2}$ is measured counterclockwise from $l_{1}$ to $l_{2}$. Under the identification $S \llbracket \Vdash \cong \mathrm{PSL}_{2}(\mathbb{R})$ from Section 2 A , we note here that

$$
k_{\delta}^{\theta_{1}, \theta_{2}}\left(g g_{1}, g g_{2}\right)=k_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{1}, g_{2}\right)
$$

for any $g, g_{1}, g_{2} \in \mathrm{PSL}_{2}(\mathbb{R})$.
Now for a given discrete subgroup $\Gamma \subset \operatorname{PSL}_{2}(\mathbb{R})$, we define the modular intersection kernel $K_{\delta}^{\theta_{1}, \theta_{2}}$ : $\Gamma \backslash \mathrm{PSL}_{2}(\mathbb{R}) \times \Gamma \backslash \mathrm{PSL}_{2}(\mathbb{R}) \rightarrow \mathbb{R}$ by taking the average of $k_{\delta}^{\theta_{1}, \theta_{2}}$ over $\Gamma$ :

$$
K_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{1}, g_{2}\right)=\sum_{\gamma \in \Gamma} k_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{1}, \gamma g_{2}\right)
$$

Note that when $\Gamma$ is cocompact, and $\delta>0$ is less than a half of the injectivity radius of $\Gamma \backslash \mathbb{H}$, we have $K_{\delta}^{\theta_{1}, \theta_{2}} \leq 1$. However, when $\Gamma \backslash \mathbb{H}$ is noncompact, $K_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{1}, g_{2}\right)$ becomes arbitrarily large near the diagonal $g_{1}=g_{2}$ as $y_{1}, y_{2} \rightarrow \infty$. This is illustrated in the following proposition when $\Gamma=\operatorname{PSL}_{2}(\mathbb{Z})$.

Proposition 2.2. Fix $0<\theta<\pi$. Then for any $1>\delta>0$, we have

$$
K_{\delta}^{0, \theta}(g, g)=\Omega_{\theta}(\delta y)
$$

Proof. Consider

$$
g=\left(\operatorname{Re}^{i(\pi / 2+\alpha(\delta))}, e^{i \alpha(\delta)}\right) \in S \Vdash,
$$

where $\alpha(\delta)$ is chosen such that the geodesic segment

$$
\beta_{g}:=\left\{\operatorname{Re}^{i \theta}:\left|\theta-\frac{\pi}{2}\right|<\alpha(\delta)\right\} \subset \mathbb{H}
$$

has length $\delta$. Note that the length of the segment does not depend on $R$ and that $\alpha(\delta) \sim \delta$ as $\delta \rightarrow 0$. From this, we infer that $\beta_{g}$ and $\beta_{g}+n$ with $0<n \ll R \delta$ intersect.

The angle of intersection is explicitly given by $2 \arcsin \frac{n}{R}$. So for all sufficiently small $0<\delta<\theta$, we have

$$
k_{\delta}^{\theta_{1}, \theta_{2}}\left(g,\left(\begin{array}{ll}
1 & n \\
0 & 1
\end{array}\right) g\right)=1
$$

for $0<n \ll R \delta$. This implies that

$$
K_{\delta}^{\theta_{1}, \theta_{2}}(g, g) \gg \delta R \gg \delta y .
$$

In view of Proposition 2.2, the following proposition provides a nice upper bound of the modular intersection kernel.

Proposition 2.3. Let $\Gamma=\mathrm{PSL}_{2}(\mathbb{Z})$ and let $1>\delta>0$. Let h be a compactly supported function on $S \sharp$, where we assume that $h((\cdot, \xi))$ is supported in $B_{\delta}(i)$ for any $\xi \in S^{1}$. Define $H: \Gamma \backslash S \sharp \times \Gamma \backslash S \llbracket$ by

$$
H\left(g_{1}, g_{2}\right)=\sum_{\gamma \in \Gamma} h\left(g_{1}^{-1} \gamma g_{2}\right)
$$

for $g_{1}, g_{2} \in \Gamma \backslash \mathrm{PSL}_{2}(\mathbb{R})$. Then for $g_{i}=\left(z_{i}, \xi_{i}\right)$ with $\operatorname{dist}_{\Gamma \backslash \mathbb{H}}\left(z_{1}, z_{2}\right)>2 \delta$, we have

$$
H\left(g_{1}, g_{2}\right)=0
$$

When $y_{1}>0$ and $y_{2}>0$ are sufficiently large, we have

$$
H\left(g_{1}, g_{2}\right) \ll \delta \sqrt{y_{1} y_{2}}\|h\|_{L^{\infty}}
$$

Proof. If $H>0$, then there exists $\gamma \in \Gamma$ such that

$$
h\left(g_{1}^{-1} \gamma g_{2}\right)>0
$$

This implies that the balls of radius $\delta$ centered at $z_{1}$ and $\gamma z_{2}$ intersect, hence

$$
\operatorname{dist}_{\mathbb{H}}\left(z_{1}, \gamma z_{2}\right)<2 \delta,
$$

which contradicts the assumption.
Now to prove the second estimate, we first note that when $y_{2}$ is sufficiently large, we have $y\left(\gamma g_{2}\right)<1$ unless $\gamma=\left(\begin{array}{ll}1 & n \\ 0 & 1\end{array}\right)$. Therefore $h\left(g_{1}^{-1} \gamma g_{2}\right)>0$ only if $\gamma=\left(\begin{array}{ll}1 & n \\ 0 & 1\end{array}\right)$. Note that $h\left(g_{1}^{-1} \gamma g_{2}\right)=1$ holds only if $\operatorname{dist}_{H-H}\left(z_{1}, n+z_{2}\right)<2 \delta$. This is equivalent to

$$
\operatorname{arccosh}\left(1+\frac{\left(n+x_{2}-x_{1}\right)^{2}+\left(y_{1}-y_{2}\right)^{2}}{y_{1} y_{2}}\right)<2 \delta,
$$

and so

$$
\left(n+x_{2}-x_{1}\right)^{2}<y_{1} y_{2}(\cosh (2 \delta)-1)-\left(y_{1}-y_{2}\right)^{2} \leq y_{1} y_{2}(\cosh (2 \delta)-1)
$$

from which we infer that there are at most $\ll \delta \sqrt{y_{1} y_{2}}$ choices of $\gamma$ which makes $h\left(g_{1}, \gamma g_{2}\right)>0$.
Now we analyze the modular intersection kernel when one variable is assumed to be contained in a compact set. We first note that if $\delta$ is less than half of the injectivity radius of $g_{0}$ in $\Gamma \backslash S \llbracket$, then for each $g \in S \sharp$, there is at most one $\gamma \in \Gamma$ such that

$$
k_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{0}, \gamma g\right) \neq 0 .
$$

Therefore $K_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{0}, \cdot\right)$ coincides with $k_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{0}, \cdot\right)$ in the $2 \delta$-neighborhood of $g_{0}$, which is a translation of $k_{\delta}^{\theta_{1}, \theta_{2}}((i, i), \cdot)$ around $(i, i)$.
Lemma 2.4. For $0<\theta_{1}<\theta_{2}<\pi$, we have

$$
\int_{\mathbb{H}} k_{\delta}^{\theta_{1}, \theta_{2}}((i, i), g) d V=\left(\cos \theta_{1}-\cos \theta_{2}\right) \delta^{2} .
$$

Assume that $0<\delta<1$. Then for any $\varepsilon=o(\delta)$ and $\varepsilon=o\left(\theta_{2}-\theta_{1}\right)$ there exist a smooth majorant $M_{\delta}^{\theta_{1}, \theta_{2}}$ and a smooth minorant $m_{\delta}^{\theta_{1}, \theta_{2}}$, i.e.,

$$
0 \leq m_{\delta}^{\theta_{1}, \theta_{2}} \leq k_{\delta}^{\theta_{1}, \theta_{2}}((i, i), \cdot) \leq M_{\delta}^{\theta_{1}, \theta_{2}}
$$

such that

$$
\int m_{\delta}^{\theta_{1}, \theta_{2}} d V \text { and } \int M_{\delta}^{\theta_{1}, \theta_{2}} d V
$$

are both

$$
\left(\cos \theta_{1}-\cos \theta_{2}\right) \delta^{2}(1+O(\varepsilon))
$$

and that

$$
\left\|m_{\delta}^{\theta_{1}, \theta_{2}}\right\|_{W^{k, \infty}}+\left\|M_{\delta}^{\theta_{1}, \theta_{2}}\right\|_{W^{k, \infty}}=O_{k}\left(\varepsilon^{-k}\right) .
$$

Proof. Note that the action of the geodesic flow of time $t$ on $S \llbracket=\mathrm{PSL}_{2}(\mathbb{R})$ is the multiplication from the right by $a\left(e^{t}\right)$. For given $\varphi \in\left(\theta_{1}, \theta_{2}\right)$, we describe the collection of $g \in \operatorname{PSL}_{2}(\mathbb{R})$ for which the corresponding geodesic segment of length $\delta$ intersects $\left\{i y: e^{\delta}>y>1\right\}$ transversally at angle $\varphi$. Note that this happens only when

$$
g a\left(e^{t_{2} / 2}\right)=\left\{\begin{array}{l}
a\left(e^{t_{1} / 2}\right) R_{\varphi / 2} \\
a\left(e^{t_{1} / 2}\right) R_{(\varphi+\pi) / 2}
\end{array}\right.
$$

for some $0<t_{1}, t_{2}<\delta$. Hence

$$
g=\left\{\begin{array}{l}
a\left(e^{t_{1} / 2}\right) R_{\varphi / 2} a\left(e^{-t_{2} / 2}\right), \\
a\left(e^{t_{1} / 2}\right) R_{(\varphi+\pi) / 2} a\left(e^{-t_{2} / 2}\right)
\end{array}\right.
$$

Consider $\Psi: A K A \rightarrow \operatorname{PSL}_{2}(\mathbb{R})$ given by

$$
\left(t_{1}, \varphi, t_{2}\right) \mapsto a\left(e^{t_{1} / 2}\right) R_{\varphi / 2} a\left(e^{-t_{2} / 2}\right)
$$

The determinant of the Jacobian of $\Psi$ is a nonzero multiple of $|\sin \varphi|$ (we refer the readers to the Appendix for the computation), and so this defines a local diffeomorphism away from $\varphi=0$ and $\pi$. Observe that $\Psi$ is injective away from $\varphi=0$ and $\pi$. From this we infer that the support of $k_{\delta}^{\theta_{1}, \theta_{2}}((i, i), g)$ is the image of the open box

$$
\left\{\left(t_{1}, \varphi, t_{2}\right): 0<t_{1}, t_{2}<\delta, \theta_{1}<\varphi<\theta_{2} \text { or } \theta_{1}+\pi<\varphi<\theta_{2}+\pi\right\}
$$

under $\Psi$, and

$$
\begin{aligned}
\int_{\mathbb{H}} k_{\delta}^{\theta_{1}, \theta_{2}}((i, i), g) d V & =\frac{1}{2} \int_{0}^{\delta} \int_{0}^{\delta} \int_{\theta_{1}}^{\theta_{2}}|\sin (\varphi)| d \varphi d t_{1} d t_{2}+\frac{1}{2} \int_{0}^{\delta} \int_{0}^{\delta} \int_{\theta_{1}+\pi}^{\theta_{2}+\pi}|\sin (\varphi)| d \varphi d t_{1} d t_{2} \\
& =\left(\cos \theta_{1}-\cos \theta_{2}\right) \delta^{2}
\end{aligned}
$$

where we used $d V=\frac{1}{2}|\sin \varphi| d \varphi d t_{1} d t_{2}((\mathrm{~A}-1))$.
Note that the support of $k_{\delta}^{\theta_{1}, \theta_{2}}((i, i), \cdot)$ is an open set which has a piecewise smooth boundary. Therefore, under the assumption that $\varepsilon=o(\delta)$ and $\varepsilon=o\left(\theta_{2}-\theta_{1}\right)$, there exist smooth majorant and minorant whose $L^{1}$ norms are $\left(\cos \theta_{1}-\cos \theta_{2}\right) \delta^{2}(1+O(\varepsilon))$, and whose $k$-th derivatives are $O_{k}\left(\varepsilon^{-k}\right)$.

As an immediate application, we have the following corollary.
Corollary 2.5. Fix a compact subset $C \subset \Gamma \backslash S \sharp$, and assume that $\delta$ is less than the half of the infimum of injectivity radius of $g \in C$ in $\Gamma \backslash S \sharp$. Then for any given compact geodesic segment $\beta \subset C$, and for any given $\varepsilon>0$ which is $o(\delta)$ and $o\left(\theta_{2}-\theta_{1}\right)$,

$$
\int_{\beta} K_{\delta}^{\theta_{1}, \theta_{2}}(s, \cdot) d s
$$

admits a smooth majorant $M_{\beta, \delta}^{\theta_{1}, \theta_{2}}$ and a smooth minorant $m_{\beta, \delta}^{\theta_{1}, \theta_{2}}$ such that

$$
\left\|m_{\beta, \delta}^{\theta_{1}, \theta_{2}}\right\|_{L^{1}},\left\|M_{\beta, \delta}^{\theta_{1}, \theta_{2}}\right\|_{L^{1}}=l(\beta)\left(\cos \theta_{1}-\cos \theta_{2}\right) \delta^{2}(1+O(\varepsilon))
$$

and that

$$
\left\|m_{\beta, \delta}^{\theta_{1}, \theta_{2}}\right\|_{W^{k, \infty}}+\left\|M_{\beta, \delta}^{\theta_{1}, \theta_{2} i}\right\|_{W^{k, \infty}}=O_{k}\left(l(\beta) \varepsilon^{-k}\right)
$$

2C. Intersection numbers. In this section, we prove formulas relating the number of intersections between two geodesics to the integral of the modular intersection kernel over the two geodesics.

Lemma 2.6. Let $\alpha_{i}=\left\{\alpha_{i}(t): t \in\left[0, l\left(\alpha_{i}\right)\right)\right\}$ be closed geodesics in $\Gamma \backslash \mathbb{H}$ parametrized by the arc length, and let $\widetilde{\alpha}_{i}=\left\{\left(\alpha_{i}(t), \alpha_{i}^{\prime}(t)\right): t \in\left[0, l\left(\alpha_{i}\right)\right)\right\} \subset S \sharp$ be the lifts of $\alpha_{i}$ for $i=1,2$. Then for any $\delta>0$,

$$
I_{\theta_{1}, \theta_{2}}\left(\alpha_{1}, \alpha_{2}\right)=\frac{1}{\delta^{2}} \int_{\widetilde{\alpha}_{2}} \int_{\widetilde{\alpha}_{1}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}
$$

Remark 2.1. For each $\alpha_{i}$, there are two choices of parametrization by the arc length, namely $\alpha_{i}(t)$ and $\alpha_{i}(-t)$, but the integral does not depend on the choice of the parametrizations.

Proof. By abuse of notations, we think of each $\alpha_{i}$ with $t \in\left[0, l\left(\alpha_{i}\right)\right)$ a geodesic segment in $\mathbb{H}$ and accordingly $\widetilde{\alpha}_{i}$ a corresponding curve in $S \sharp \Vdash$. For a geodesic segment $\alpha \subset \mathbb{H}$ parametrized by $t \in[a, b]$, let $[\alpha] \subset \mathbb{H}$ be the biinfinite geodesic $\{\alpha(t): t \in \mathbb{R}\}$ that contains $\alpha$. Then we express the integral as follows:

$$
\begin{aligned}
\frac{1}{\delta^{2}} \int_{\widetilde{\alpha}_{2}} \int_{\widetilde{\alpha}_{1}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} & =\sum_{\gamma \in \Gamma} \frac{1}{\delta^{2}} \int_{\gamma \widetilde{\alpha}_{2}} \int_{\widetilde{\alpha}_{1}} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} \\
& =\sum_{\gamma \in \Gamma / \Gamma_{\left[\alpha_{2}\right]}} \frac{1}{\delta^{2}} \int_{\left.\gamma \widetilde{\alpha_{2}}\right]} \int_{\widetilde{\alpha}_{1}} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} \\
& =\sum_{\left.\gamma \in \Gamma_{\left[\alpha_{1}\right]} \backslash \Gamma / \Gamma_{\left[\alpha_{2}\right]}\right]} \sum_{\gamma^{\prime} \in \Gamma_{\left[\alpha_{1}\right]}} \frac{1}{\delta^{2}} \int_{\gamma^{\prime} \gamma\left[\widetilde{\left.\alpha_{2}\right]}\right.} \int_{\widetilde{\alpha}_{1}} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} \\
& =\sum_{\gamma \in \Gamma_{\left[\alpha_{1}\right]} \backslash \Gamma / \Gamma_{\left[\alpha_{2}\right]}} \frac{1}{\delta^{2}} \int_{\gamma\left[\widetilde{\left.\alpha_{2}\right]}\right.} \int_{\widetilde{\left.\alpha_{1}\right]}} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} .
\end{aligned}
$$

Here $\Gamma_{\left[\alpha_{i}\right]}$ is the stabilizer subgroup of $\Gamma$ with respect to $\left[\alpha_{i}\right]$.
Now because two geodesics in $\mathbb{H}$ may intersect at most once, for each intersection point $p \in \alpha_{1} \cap \alpha_{2}$ on $\Gamma \backslash \mathbb{H}$, there exists a unique $\gamma \in \Gamma / \Gamma_{\left[\alpha_{2}\right]}$ such that $\alpha_{1}$ and $\gamma\left[\alpha_{2}\right]$ intersect at a lift of $p$. Also, because $\left[\alpha_{1}\right]$ is a disjoint union of $\gamma^{\prime} \alpha_{1}$ with $\gamma^{\prime} \in \Gamma_{\left[\alpha_{1}\right]}$, each $\left\{\gamma^{\prime} \gamma: \gamma^{\prime} \in \Gamma_{\left[\alpha_{1}\right]}\right\}$ contains at most one $\gamma^{\prime} \gamma$ such that $\gamma^{\prime} \gamma\left[\alpha_{2}\right]$ intersects $\alpha_{1}$.

Therefore the intersections of $\alpha_{1}$ and $\alpha_{2}$ are in one-to-one correspondence with $\gamma \in \Gamma_{\left[\alpha_{1}\right]} \backslash \Gamma / \Gamma_{\left[\alpha_{2}\right]}$ such that $\gamma\left[\alpha_{2}\right]$ intersects $\left[\alpha_{1}\right]$. We complete the proof by observing that

$$
\int_{\gamma\left[\widetilde{\alpha_{2}}\right]} \int_{\left[\widetilde{\alpha_{1}}\right]} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}=1,
$$

if $\left[\alpha_{1}\right]$ and $\gamma\left[\alpha_{2}\right]$ intersect at an angle $\in\left(\theta_{1}, \theta_{2}\right)$, and $=0$ otherwise.
Now let $\beta=\{\beta(t): t \in[0, l(\beta))\}$ be a compact geodesic segment in $\Gamma \backslash \mathbb{H}$, and let $\alpha_{2}$ be a closed geodesic as before. Then

$$
\frac{1}{\delta^{2}} \int_{\widetilde{\alpha}_{2}} \int_{\tilde{\beta}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}
$$

does not always give $I\left(\beta, \alpha_{2}\right)$. Instead, it is a weighted sum over the intersections of $\beta_{0}:=\{\beta(t): t \in$ $[0, l(\beta)+\delta)\}$ and $\alpha_{2}$. We prove the following.

Lemma 2.7. With the same notations as above, assume that $0<\delta<l(\beta)$ and that $\beta_{0}$ has no self intersection. For $0<\theta_{1}<\theta_{2}<\pi$, let $S\left(\beta_{0}, \alpha_{2}\right)_{\theta_{1}, \theta_{2}}$ be the set of intersections between $\beta_{0}$ and $\alpha_{2}$ where the intersection angle is $\in\left(\theta_{1}, \theta_{2}\right)$. Then we have

$$
\frac{1}{\delta^{2}} \int_{\tilde{\alpha}_{2}} \int_{\tilde{\beta}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}=\sum_{p \in S\left(\beta_{0}, \alpha_{2}\right)_{\theta_{1}, \theta_{2}}} \min \left\{\frac{\beta^{-1}(p)}{\delta}, 1, \frac{l(\beta)+\delta-\beta^{-1}(p)}{\delta}\right\}
$$

Proof. As in the proof of Lemma 2.6, we first have

$$
\begin{aligned}
\frac{1}{\delta^{2}} \int_{\widetilde{\alpha}_{2}} \int_{\tilde{\beta}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} & =\sum_{\gamma \in \Gamma} \frac{1}{\delta^{2}} \int_{\gamma \widetilde{\alpha}_{2}} \int_{\tilde{\beta}} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} \\
& =\sum_{\gamma \in \Gamma / \Gamma_{\left[\alpha_{2}\right]}} \frac{1}{\delta^{2}} \int_{\gamma\left[\widetilde{\alpha_{2}}\right]} \int_{\tilde{\beta}} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}
\end{aligned}
$$

Note that because we assumed that $\beta_{0}$ has no self-intersection, $p \in S\left(\beta_{0}, \alpha_{2}\right)_{\theta_{1}, \theta_{2}}$ is in one-to-one correspondence with $\gamma \in \Gamma / \Gamma_{\left[\alpha_{2}\right]}$ such that $\beta_{0}$ and $\gamma\left[\widetilde{\alpha_{2}}\right]$ intersect at $p$ at an angle $\in\left(\theta_{1}, \theta_{2}\right)$. We denote by $\gamma_{p}$ the $\gamma$ corresponding to $p$. Observe that

$$
\int_{\gamma\left[\widetilde{\alpha_{2}}\right]} \int_{\tilde{\beta}} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}=0,
$$

if $\gamma\left[\widetilde{\alpha_{2}}\right] \cap \beta_{0}=\varnothing$. So it is sufficient to prove that

$$
\frac{1}{\delta^{2}} \int_{\gamma_{p}\left[\widetilde{\alpha_{2}}\right]} \int_{\tilde{\beta}} k_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}=\min \left\{\frac{\beta^{-1}(p)}{\delta}, 1, \frac{l(\beta)+\delta-\beta^{-1}(p)}{\delta}\right\}
$$

This follows by observing that

$$
k_{\delta}^{\theta_{1}, \theta_{2}}\left(\left(\beta\left(t_{1}\right), \beta^{\prime}\left(t_{1}\right)\right),\left(\gamma_{p} \alpha_{2}\left(t_{2}\right),\left(\gamma_{p} \alpha_{2}\right)^{\prime}\left(t_{2}\right)\right)\right)=1
$$

for

$$
\left(t_{1}, t_{2}\right) \in\left(\beta^{-1}(p)-\delta, \beta^{-1}(p)\right) \times\left(\alpha_{2}^{-1}(p)-\delta, \alpha_{2}^{-1}(p)\right)
$$

and 0 otherwise, whereas the integral over $\tilde{\beta}$ is over the range $t_{1} \in(0, l(\beta))$.

## 3. Spectral theory

3A. Spectral expansion. We first go over the spectral decomposition of $L^{2}(S \rtimes)$. Readers may find more details on the subject in [Kubota 1973; Lang 1985]. On $G=\operatorname{PSL}_{2}(\mathbb{R})$, there is a differential operator of order 2 that commutes with the $G$ action,

$$
\Omega=y^{2} \partial_{x}^{2}+y^{2} \partial_{y}^{2}+y \partial_{x} \partial_{\theta}
$$

which is called the Casimir operator. An equivariant eigenfunction of $\Omega$ is a function $f \in C^{\infty}(S \mathbb{X})$ that satisfies

$$
\Omega f=\lambda f
$$

for some $\lambda \in \mathbb{R}$, and

$$
\begin{equation*}
f\left(g R_{\theta}\right)=e^{-i m \theta} f(g) \tag{3-1}
\end{equation*}
$$

for some $m \in 2 \mathbb{Z}$. We say that a function has weight $m$ if it satisfies (3-1).
Each irreducible (cuspidal) subrepresentation of the right regular representation

$$
\rho_{g}: f(h) \mapsto f(h g)
$$

on $L^{2}(S \rtimes)$ is generated by an equivariant eigenfunction of $\Omega$.
We let $\boldsymbol{E}^{+}$and $\boldsymbol{E}^{-}$to be the raising and lowering operator acting on equivariant functions on $L^{2}(S \mathbb{X})$, which are given by [Jakobson 1994]

$$
\begin{equation*}
\boldsymbol{E}^{+}=e^{-2 i \theta}\left(2 i y \partial_{x}+2 y \partial_{y}+i \partial_{\theta}\right) \quad \text { and } \quad \boldsymbol{E}^{-}=e^{2 i \theta}\left(2 i y \partial_{x}-2 y \partial_{y}+i \partial_{\theta}\right) . \tag{3-2}
\end{equation*}
$$

Note that $\boldsymbol{E}^{+}$(resp. $\boldsymbol{E}^{-}$) maps a weight $m$ eigenfunction of $\Omega$ to a weight $m+2$ (resp. $m-2$ ) eigenfunction of $\Omega$.

For an even integer $m$ let

$$
\psi_{s, m}(g)=y^{s} e^{-i m \theta}
$$

Note that $\psi_{s, m}$ is invariant under the action of the unipotent upper triangular matrices. The weight $m$ Eisenstein series is then given by

$$
E_{m}(g, s)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \psi_{s, m}(\gamma g)
$$

where $\Gamma_{\infty}=\left\{\left(\begin{array}{ll}1 & n \\ 0 & 1\end{array}\right): n \in \mathbb{Z}\right\}$ is the stabilizer subgroup of $\Gamma$ with respect to the cusp $i \infty$. Although the right-hand side of the equation is absolutely convergent only for $\operatorname{Re}(s)>1$, the weight $m$ Eisenstein series has a meromorphic continuation to the entire complex plane.

Let $\Theta$ be the closure of

$$
\left\{\int_{-\infty}^{\infty} h(t) E_{m}\left(g, \frac{1}{2}+i t\right) d t: h(t) \in C_{0}^{\infty}(\mathbb{R}), m \in 2 \mathbb{Z}\right\}
$$

in $L^{2}(S \mathbb{X})$, and let

$$
L_{\text {cusp }}^{2}(S \mathbb{X})=\left\{f \in L^{2}(S \mathbb{X}): \int_{0}^{1} f(n(x) g) d x=0 \text { for almost every } g \in S \mathbb{X}\right\}
$$

be the space of cusp forms. Then we have the decomposition

$$
L^{2}(S \mathbb{X})=\langle\{1\}\rangle \oplus \Theta \oplus L_{\text {cusp }}^{2}(S \mathbb{X}),
$$

where $\langle\{1\}\rangle$ is the subspace spanned by a constant function.

We express the cuspidal subspace as a direct sum of subspaces generated by Maass forms and modular forms as in [Luo et al. 2009, (1.10)],

$$
L_{\mathrm{cusp}}^{2}(S \mathbb{X})=\sum_{j=1}^{\infty} W_{\pi_{j}^{0}} \bigoplus \sum_{m \geq 12} \sum_{j=1}^{d_{m}}\left(W_{\pi_{j}^{m}} \oplus W_{\pi_{j}^{-m}}\right)
$$

where each $W_{\pi_{j}^{m}}$ corresponds to a $G$ and Hecke irreducible subspace of a right regular representation on $L_{\text {cusp }}^{2}$. Here $d_{m}$ is the dimension of the space of holomorphic cusp forms of weight $m$ for $\mathrm{PSL}_{2}(\mathbb{Z})$. Each $\pi_{j}^{0}$ corresponds to a Maass-Hecke cusp form which we denote by $\phi_{j}^{0}$. For $m>0, \pi_{j}^{m}$ corresponds to a holomorphic Hecke cusp form $\phi_{j}^{m}$. We identify a weight $m$ function on $\Gamma \backslash \mathbb{H}$

$$
f(\gamma z)=(c z+d)^{m} f(z) \quad \text { for }\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)=\gamma \in \Gamma
$$

with a weight $m \Gamma$-invariant function $F$ on $\operatorname{PSL}_{2}(\mathbb{R})$ via

$$
\begin{equation*}
F(g)=y^{m / 2} f(z) e^{-i m \theta} \tag{3-3}
\end{equation*}
$$

When $m \geq 0$, viewing $\phi_{j}^{m}$ as a function on $S \mathbb{X}$, each $W_{\pi_{j}^{m}}$ is spanned by

$$
\ldots, \quad\left(\boldsymbol{E}^{-}\right)^{3} \phi_{j}^{m}, \quad\left(\boldsymbol{E}^{-}\right)^{2} \phi_{j}^{m}, \quad \boldsymbol{E}^{-} \phi_{j}^{m}, \quad \phi_{j}^{m}, \quad \boldsymbol{E}^{+} \phi_{j}^{m}, \quad\left(\boldsymbol{E}^{+}\right)^{2} \phi_{j}^{m}, \quad\left(\boldsymbol{E}^{+}\right)^{3} \phi_{j}^{m}, \quad \ldots
$$

Note that when $m>0, \boldsymbol{E}^{-} \phi_{j}^{m}=0$.
For $m<0$, we set

$$
W_{\pi_{j}^{-m}}=\overline{W_{\pi_{j}^{m}}}=\left\{\bar{f}: f \in W_{\pi_{j}^{m}}\right\}
$$

Now let

$$
U_{\pi_{j}^{0}}=W_{\pi_{j}^{0}} \quad \text { and } \quad U_{\pi_{j}^{m}}=W_{\pi_{j}^{m}} \oplus W_{\pi_{j}^{-m}}
$$

when $m>0$. We specify an orthonormal basis of each $U_{\pi_{j}^{m}}$ as follows.
The Maass cusp form case $m=0$ : Let $-\left(\frac{1}{4}+t_{j}^{2}\right)$ be the Laplacian eigenvalue of $\phi_{j}^{0}{ }^{\dagger}$ for some real $t_{j}$. We set $\phi_{j, 0}^{0}=\phi_{j}^{0}$, and define $\phi_{j, l}^{0}$ for $l \in 2 \mathbb{Z}$ inductively by

$$
\begin{equation*}
\boldsymbol{E}^{-} \phi_{j, l}^{0}=\left(l+1-2 i t_{j}\right) \phi_{j, l-2}^{0} \quad \text { and } \quad \boldsymbol{E}^{+} \phi_{j, l}^{0}=\left(l+1+2 i t_{j}\right) \phi_{j, l+2}^{0} \tag{3-4}
\end{equation*}
$$

The holomorphic Hecke cusp form case $m>0$ : We set $\phi_{j, m}^{m}=\phi_{j}^{m}$ and $\phi_{j,-m}^{m}=\overline{\phi_{j}^{m}}$, and define $\phi_{j, l}^{m}$ for $l \in 2 \mathbb{Z}$ inductively by

$$
\begin{equation*}
\boldsymbol{E}^{-} \phi_{j, l}^{m}=(l-m) \phi_{j, l-2}^{m} \quad \text { and } \quad \boldsymbol{E}^{+} \phi_{j, l}^{m}=(l+m) \phi_{j, l+2}^{m} \tag{3-5}
\end{equation*}
$$

Finally, note that we have the following relation among the weight $m$ Eisenstein series.

$$
\begin{aligned}
& \boldsymbol{E}^{-} E_{m}\left(g, \frac{1}{2}+i t\right)=(m+1-2 i t) E_{m-2}\left(g, \frac{1}{2}+i t\right), \quad \text { and } \\
& \boldsymbol{E}^{+} E_{m}\left(g, \frac{1}{2}+i t\right)=(m+1+2 i t) E_{m+2}\left(g, \frac{1}{2}+i t\right)
\end{aligned}
$$

[^1]With these notations, we have:
Proposition 3.1. Let $f \in L^{2}(S \backslash)$. Then we have

$$
\begin{aligned}
& f(g)= \\
& \quad \frac{3}{\pi^{2}} \int_{S X} f\left(g_{1}\right) d g_{1}+\sum_{\substack{m \geq 0 \\
2 \mid m}} \sum_{j=1}^{d_{m}} \sum_{\substack{l \in 2 \mathbb{Z} \\
|l| \geq m}}\left\langle f, \phi_{j, l}^{m}\right\rangle_{S X} \phi_{j, l}^{m}(g)+\sum_{m \in 2 \mathbb{Z}} \frac{1}{4 \pi} \int_{-\infty}^{\infty}\left\langle f, E_{m}\left(\cdot, \frac{1}{2}+i t\right)\right\rangle_{S X} E_{m}\left(g, \frac{1}{2}+i t\right) d t,
\end{aligned}
$$

where we set $d_{0}=+\infty$.

## 4. Effective equidistribution

4A. Invariant linear form. Define $\mu_{d}$ to be the integral over discriminant $d$ oriented closed geodesics on $S \mathbb{X}$,

$$
\mu_{d}(F):=\int_{\mathscr{C}_{d}} F(s) d s=\sum_{\operatorname{disc}(q)=d} \int_{C(q)} F(s) d s .
$$

where $C(q) \subset S \mathbb{X}$ is the oriented closed geodesic associated to the binary quadratic form $q$ [Luo et al. 2009, 2.3]. Then for any $F \in U_{\pi_{j}^{m}}$, we have

$$
\mu_{d}(F)=\mu_{d}\left(\phi_{j}^{m}\right) \eta_{j}^{m}(F)
$$

for some linear form $\eta_{j}^{m}$ on $U_{\pi_{j}^{m}}$ invariant under the diagonal action [loc. cit., Section 3.7.1], which we describe below following [loc. cit., Section 3.2]. (Note that the parameter $s$ in [loc. cit.] is replaced by $2 i t$ in this article for consistency.)
 have

$$
\begin{equation*}
\eta_{j}^{0}\left(\phi_{j, l}^{0}\right)=\eta_{j}^{0}\left(\phi_{j,-l}^{0}\right)=\frac{\left(1-2 i t_{j}\right)\left(5-2 i t_{j}\right) \cdots\left(l-3-2 i t_{j}\right)}{\left(3+2 i t_{j}\right)\left(7+2 i t_{j}\right) \cdots\left(l-1+2 i t_{j}\right)}, \tag{4-1}
\end{equation*}
$$

and $\eta_{j}^{0}\left(\phi_{j, l}^{0}\right)$ is identically 0 if $l \equiv 2(\bmod 4)$. Note that $\left\{\phi_{j, l}^{0}\right\}_{l \in 2 \mathbb{Z}}$ is an orthogonal basis of $U_{\pi_{j}^{0}}$, and normalized so that,

$$
\left\|\phi_{j, l}^{0}\right\|_{L^{2}}=\left\|\phi_{j}^{0}\right\|_{L^{2}} .
$$

The holomorphic Hecke cusp form case $m>0$ : Let $\phi_{j, l}^{m}$ be the holomorphic Hecke cusp form defined by $(3-5)$. When $m \equiv 2(\bmod 4), \eta_{j}^{m}$ is identically 0 .

When $m \equiv 0(\bmod 4)$, for $l \geq 4$ with $4 \mid l$,

$$
\begin{equation*}
\eta_{j}^{m}\left(\phi_{j, m+l}^{m}\right)=\eta_{j}^{m}\left(\phi_{j,-m-l}^{m}\right)=\frac{1 \cdot 3 \cdot 5 \cdots(l / 2-1)}{(m+1)(m+3) \cdots(m+l / 2-1)}, \tag{4-2}
\end{equation*}
$$

and $\eta_{j}^{m}\left(\phi_{j, m+l}^{m}\right)$ vanishes for $l \equiv 2(\bmod 4)$.
Note that $\left\{\phi_{j, l}^{m}\right\}_{l \in 2 \mathbb{Z},|l| \geq m}$ is an orthogonal basis of $U_{\pi_{j}^{m}}$, and normalized so that

$$
\left\|\phi_{j, l}^{m}\right\|_{L^{2}}=\left\|\phi_{j}^{m}\right\|_{L^{2}}
$$

for $l \in 2 \mathbb{Z},|l| \geq m$.

Eisenstein series case: By the above identities and following [Luo et al. 2009, Section 3], we have

$$
\mu_{d}\left(E_{m}\left(g, \frac{1}{2}+i t\right)\right)=\eta(m, t) \mu_{d}\left(E_{0}\left(g, \frac{1}{2}+i t\right)\right)
$$

where for $m \geq 4$ such that $4 \mid m$,

$$
\begin{equation*}
\eta(m, t)=\eta(-m, t)=\frac{(1-2 i t)(5-2 i t) \cdots(2 m-3-2 i t)}{(3+2 i t)(7+2 i t) \cdots(2 m-1+2 i t)} \tag{4-3}
\end{equation*}
$$

and $\eta(m, t)$ is identically 0 if $m \equiv 2(\bmod 4)$.

## 4B. Period integrals.

4B1. Holomorphic cusp forms. In this section, we give an upper bound on the period integrals of holomorphic forms. We first use the results of Shintani to relate the period integrals of holomorphic cusp forms to the Fourier coefficients of half integral holomorphic forms. We then apply the result of Kohnen and Zagier [1981] which gives an explicit version of the Waldspurger's formula for the Fourier coefficients of half integral holomorphic forms. An upper bound on these period integrals is deduced by using the subconvexity bounds on the central value of the $L$-functions and the Ramanujan bound on the Fourier coefficients of holomorphic modular forms.

Note that $c(d)$ is identically zero when $m \equiv 2(\bmod 4)$, and so we assume that $4 \mid m$. Let $\hat{\phi}_{j}^{m}$ be a normalization of the Hecke holomorphic cusp form $\phi_{j}^{m}$ of weight $m$ such that $a_{1}=1$. Let

$$
c(d):=\sum_{\operatorname{disc}(q)=d} \int_{C(q)} \hat{\phi}_{j}^{m}(z) q(z, 1)^{m / 2-1} d z,
$$

where $\hat{\phi}_{j}^{m}(z)$ is the associated holomorphic modular form defined on the upper half plane and the integration is on the upper half plane (3-3). By [Luo et al. 2009, (2.4) page 14], we have

$$
\begin{equation*}
|c(d)|=|d|^{m / 4-1 / 2}\left|\mu_{d}\left(\hat{\phi}_{j}^{m}\right)\right| . \tag{4-4}
\end{equation*}
$$

Let

$$
\theta\left(z, \hat{\phi}_{j}^{m}\right):=\sum_{d \geq 1} c(d) e(d z)
$$

By [Shintani 1975, Theorem 2], $\theta\left(z, \phi_{j}^{m}\right)$ is a Hecke holomorphic cusp form of weight $(m+1) / 2$ and level $\Gamma_{0}(4)$. By [Luo et al. 2009, (6.2), page 37], we have the following explicit version of Rallis inner product formula

$$
\left\langle\theta\left(\hat{\phi}_{j}^{m}\right), \theta\left(\hat{\phi}_{j}^{m}\right)\right\rangle=\frac{(m / 2-1)!}{2^{m} \pi^{m / 2}} L\left(\frac{1}{2}, \phi_{j}^{m}\right)\left\langle\hat{\phi}_{j}^{m}, \hat{\phi}_{j}^{m}\right\rangle
$$

Suppose that $d=D b^{2}$ with $D$ a fundamental discriminant. By [Kohnen and Zagier 1981, Theorem 1], for $D$ a fundamental discriminant with $D>0$ and $4 \mid m$, we have

$$
\frac{c(D)^{2}}{\left\langle\theta\left(\hat{\phi}_{j}^{m}\right), \theta\left(\hat{\phi}_{j}^{m}\right)\right\rangle}=\frac{(m / 2-1)!}{\pi^{m / 2}} D^{(m-1) / 2} \frac{L\left(1 / 2, \phi_{j}^{m} \otimes \chi_{D}\right)}{\left\langle\hat{\phi}_{j}^{m}, \hat{\phi}_{j}^{m}\right\rangle},
$$

which implies that

$$
|c(D)|=D^{(m-1) / 4} \frac{(m / 2-1)!}{2^{m / 2} \pi^{m / 2}}\left(L\left(\frac{1}{2}, \phi_{j}^{m}\right) L\left(\frac{1}{2}, \phi_{j}^{m} \otimes \chi_{D}\right)\right)^{1 / 2}
$$

By using the Ramanujan bound on the Fourier coefficients of integral weight cusp forms and the above, we have

$$
|c(d)| \ll_{\epsilon} b^{m-1 / 2+\epsilon}|c(D)|<_{\epsilon} d^{m-1 / 4+\epsilon} \frac{(m / 2-1)!}{2^{m / 2} \pi^{m / 2}}\left(L\left(\frac{1}{2}, \phi_{j}^{m}\right) L\left(\frac{1}{2}, \phi_{j}^{m} \otimes \chi_{D}\right)\right)^{1 / 2}
$$

and so

$$
\left|\mu_{d}\left(\hat{\phi}_{j}^{m}\right)\right| \lll<|d|^{1 / 4+\epsilon} \frac{(m / 2-1)!}{2^{m / 2} \pi^{m / 2}}\left(L\left(\frac{1}{2}, \phi_{j}^{m}\right) L\left(\frac{1}{2}, \phi_{j}^{m} \otimes \chi_{D}\right)\right)^{1 / 2}
$$

by (4-4).
We now use the convexity bound

$$
L\left(\frac{1}{2}, \phi_{j}^{m}\right) \ll_{\epsilon} m^{1 / 2+\epsilon},
$$

and the subconvexity bound [Blomer et al. 2007, Theorem 1]

$$
L\left(\frac{1}{2}, \phi_{j}^{m} \otimes \chi_{D}\right) \ll_{\epsilon} m^{(75+12 \theta) / 16} D^{1 / 2-(1 / 8)(1-2 \theta)+\epsilon}
$$

where $\theta=\frac{7}{64}$ is the best exponent toward Ramanujan conjecture for Maass forms, to see that

$$
\left|\mu_{d}\left(\hat{\phi}_{j}^{m}\right)\right| \lll \epsilon d^{1 / 4+\epsilon} \frac{(m / 2-1)!}{2^{m / 2} \pi^{m / 2}} m^{2.64} D^{1 / 4-25 / 512}
$$

It is well-known that

$$
\left\langle\hat{\phi}_{j}^{m}, \hat{\phi}_{j}^{m}\right\rangle=\frac{\Gamma(m)}{(4 \pi)^{m}} L\left(1, \operatorname{sym}^{2} \phi_{j}^{m}\right)
$$

up to a constant. Hence, by Stirling's approximation

$$
\begin{equation*}
\left|\mu_{d}\left(\phi_{j}^{m}\right)\right| \lll d^{1 / 4+\epsilon} m^{2.9} D^{1 / 4-25 / 512} \ll d^{1 / 2-25 / 512+\epsilon} m^{2.9} . \tag{4-5}
\end{equation*}
$$

4B2. Maass forms. In this section, we give an upper bound on the period integrals of Maass forms. We first recall some results of Katok and Sarnak [1993] that generalize the work of Shintani [1975] to Maass forms and related the period integrals to the Fourier coefficients of half integral weight Maass forms. Then we use an explicit version of the Waldspurger formula [Baruch and Mao 2010] and give a nontrivial bound on these period integrals by using the subconvexity bound on the central value of the $L$-functions and the best bound toward Ramanujan conjecture for Maass forms.

Let $\phi_{j}^{0}$ be a Hecke-Maass form with $\left\langle\phi_{j}^{0}, \phi_{j}^{0}\right\rangle=1$ and with the Laplacian eigenvalue $-\left(\frac{1}{4}+t_{j}^{2}\right)$. For $d>0$, let

$$
\rho(d):=\frac{1}{\sqrt{8} \pi^{1 / 4} d^{3 / 4}} \sum_{\operatorname{disc}(q)=d} \int_{C(q)} \phi_{j}^{0} d s
$$

be the associated period integral, and let

$$
\theta\left((u+i v), \phi_{j}^{0}\right):=\sum_{d \neq 0} \rho(d) W_{\operatorname{sgn}(d) / 4, i t_{j} / 2}(4 \pi|d| v) e(d u),
$$

where $W_{\operatorname{sgn}(d) / 4, i t_{j} / 2}$ is the usual Whittaker function. Here $\rho(d)$ for $d<0$ is the sum of $\phi_{j}^{0}$ over the CM points with the discriminant $d$ appropriately normalized; see [Katok and Sarnak 1993, page 197] or [Sardari 2021, Section 3.3] for a detailed discussion.

Note from [Katok and Sarnak 1993] that $\theta\left((u+i v), \phi_{j}^{0}\right)$ is a weight $\frac{1}{2}$ Hecke-Maass form with the Laplacian eigenvalue $-\left(\frac{1}{4}+\frac{t_{j}^{2}}{4}\right)$. By [Katok and Sarnak 1993, (5.6), page 224] or [Luo et al. 2009, (6.4), page 38], we have the following version of the Rallis inner product formula

$$
\left\langle\theta\left(\phi_{j}^{0}\right), \theta\left(\phi_{j}^{0}\right)\right\rangle=\frac{3}{2} \Lambda\left(\frac{1}{2}, \phi_{j}^{0}\right),
$$

where

$$
\Lambda\left(s, \phi_{j}^{0}\right)=\pi^{-s} \Gamma\left(\frac{s+i t_{j}}{2}\right) \Gamma\left(\frac{s-i t_{j}}{2}\right) L\left(s, \phi_{j}^{0}\right)
$$

is the completed $L$-function.
By an explicit form of Waldspurger formula [Baruch and Mao 2010, Theorem 1.4], and the best exponent toward the Ramanujan conjecture [Lester and Radziwiłł 2020, Corollary 6.1], we have

$$
\frac{\rho(d)}{\left\langle\theta\left(\phi_{j}^{0}\right), \theta\left(\phi_{j}^{0}\right)\right\rangle^{1 / 2}} \ll \epsilon \frac{1}{\sqrt{|d|}}\left(\frac{L\left(1 / 2, \phi_{j}^{0} \otimes \chi_{D}\right)}{L\left(1, \operatorname{sym}^{2} \phi_{j}^{0}\right)}\right)^{1 / 2} b^{7 / 64+\epsilon}\left|t_{j}\right|^{-\operatorname{sgn}(d) / 4} e^{\pi\left|t_{j}\right| / 4}
$$

where $d=D b^{2}$ with $D$ a fundamental discriminant. Note from Stirling's formula that

$$
\Gamma\left(\frac{1 / 2+i t_{j}}{2}\right) \Gamma\left(\frac{1 / 2-i t_{j}}{2}\right) \ll\left|t_{j}\right|^{-1 / 2} e^{-\pi\left|t_{j}\right| / 2}
$$

from which we infer that

$$
\begin{aligned}
\mu_{d}\left(\phi_{j}^{0}\right) & \lll d^{3 / 4}|\rho(d)| \\
& \lll \epsilon d^{1 / 4}\left(\Lambda\left(\frac{1}{2}, \phi_{j}^{0}\right)\right)^{1 / 2}\left(\frac{L\left(1 / 2, \phi_{j}^{0} \otimes \chi_{D}\right)}{L\left(1, \operatorname{sym}^{2} \phi_{j}^{0}\right)}\right)^{1 / 2} b^{7 / 64+\epsilon}\left|t_{j}\right|^{-\operatorname{sgn}(d) / 4} e^{\pi\left|t_{j}\right| / 4} \\
& \lll \epsilon d^{1 / 4}\left(L\left(\frac{1}{2}, \phi_{j}^{0}\right) L\left(\frac{1}{2}, \phi_{j}^{0} \otimes \chi_{D}\right)\right)^{1 / 2} b^{7 / 64+\epsilon}\left|t_{j}\right|^{-((\operatorname{sgn}(d)+1) / 4)+\epsilon}
\end{aligned}
$$

We now use the convexity bound,

$$
L\left(\frac{1}{2}, \phi_{j}^{0}\right) \ll_{\epsilon}\left|t_{j}\right|^{1 / 2+\epsilon},
$$

and the subconvexity bound [Blomer et al. 2007, Theorem 1],

$$
L\left(\frac{1}{2}, \phi_{j}^{0} \otimes \chi_{D}\right) \ll_{\epsilon}\left|t_{j}\right|^{(31+4 \theta+\epsilon) / 16} D^{1 / 2-(1-2 \theta) / 8+\epsilon},
$$

to conclude that

$$
\begin{equation*}
\mu_{d}\left(\phi_{j}^{0}\right) \lll \epsilon d^{1 / 4+\epsilon}\left|t_{j}\right|^{3 / 4} b^{7 / 64+\epsilon} D^{1 / 4-25 / 512} \ll d^{1 / 2-25 / 512+\epsilon}\left|t_{j}\right|^{3 / 4} \tag{4-6}
\end{equation*}
$$

4B3. Eisenstein series. For a nonsquare integer $d \equiv 0,1(\bmod 4)$, let $d=D b^{2}$ where $D$ is a fundamental discriminant. Then we have the following explicit formula for the period integral of the Eisenstein series [Zagier 1981, page 282]: ${ }^{\dagger}$

$$
\begin{equation*}
\mu_{d}\left(E_{0}(\cdot, s)\right)=\frac{\Gamma(s / 2)^{2} d^{s / 2} L(s, d)}{\Gamma(s) \zeta(2 s)} \tag{4-7}
\end{equation*}
$$

where

$$
\begin{equation*}
L(s, d)=L\left(s, \chi_{D}\right)\left(\sum_{a \mid b} \mu(a)\left(\frac{D}{a}\right) a^{-s} \sigma_{1-2 s}\left(\frac{b}{a}\right)\right) \tag{4-8}
\end{equation*}
$$

Here $L\left(s, \chi_{D}\right)$ is the Dirichlet $L$-function attached to the quadratic Dirichlet character $\chi_{D}(\cdot)=\left(\frac{D}{\cdot}\right)$, $\mu(\cdot)$ is the Möbius function, and $\sigma_{v}(\cdot)=\sum_{a \mid} \cdot a^{v}$ is the divisor function.

Now assume that $s=\frac{1}{2}+i t$ for some $t \in \mathbb{R}$. By Stirling's formula, we have

$$
\frac{\Gamma(s / 2)^{2}}{\Gamma(s)} \ll|t|^{-1 / 2}
$$

By the zero free region of $\zeta(2 s)$ around $2 s=1+2 i t$, we have

$$
|\zeta(2 s)| \ggg_{\epsilon} t^{-\epsilon} .
$$

We also have the convexity bound

$$
\zeta(s) \ll|t|^{1 / 4}
$$

and we know from [Heath-Brown 1980] that

$$
L\left(\frac{1}{2}+i t, \chi_{D}\right) \ll_{\epsilon}((|t|+1) D)^{3 / 16+\epsilon} .
$$

Finally, observe that we have

$$
\sum_{a \mid b} \mu(a)\left(\frac{D}{a}\right) a^{-s} \sigma_{1-2 s}\left(\frac{b}{a}\right) \ll_{\epsilon} d^{\epsilon}
$$

Combining all these estimates, we deduce the following estimate from (4-7) for $s=\frac{1}{2}+i t$ :

$$
\begin{equation*}
\mu_{d}\left(E_{0}(\cdot, s)\right) \lll \epsilon d^{1 / 2-1 / 16+\epsilon} \tag{4-9}
\end{equation*}
$$

4C. Proof of Theorem 1.3. For any compactly supported smooth function $F \in C_{0}^{\infty}(S \mathbb{X})$, recall from Proposition 3.1 that we have

$$
\begin{aligned}
& F(g) \\
& =\frac{3}{\pi^{2}} \int_{S X} F\left(g_{1}\right) d g_{1}+\sum_{\substack{m \geq 0 \\
2 \mid m}} \sum_{j=1}^{d_{m}} \sum_{\substack{l \in \mathbb{Z} \\
|l| \geq m}}\left\langle F, \phi_{j, l}^{m}\right\rangle_{S X} \phi_{j, l}^{m}(g)+\sum_{m \in 2 \mathbb{Z}} \frac{1}{4 \pi} \int_{-\infty}^{\infty}\left\langle F, E_{m}\left(\cdot, \frac{1}{2}+i t\right)\right\rangle_{S X} E_{m}\left(g, \frac{1}{2}+i t\right) d t,
\end{aligned}
$$

[^2]and so from the discussion of Section 4A, we have
\[

$$
\begin{aligned}
\mu_{d}(F)=\mu_{d}\left(\frac{3}{\pi^{2}}\right) \int_{S \mathbb{X}} F(g) d g & +\sum_{\substack{m \geq 0 \\
4 \mid m}} \sum_{j=1}^{d_{m}} \mu_{d}\left(\phi_{j}^{m}\right) \sum_{\substack{l \in 4 \mathbb{Z} \\
|l| \geq m}}\left\langle F, \phi_{j, l \mid}^{m}\right\rangle_{S \mathbb{X}} \eta_{j}^{m}\left(\phi_{j, l}^{m}\right) \\
& +\sum_{m \in 4 \mathbb{Z}} \frac{1}{4 \pi} \int_{-\infty}^{\infty}\left\langle F, E_{m}\left(\cdot, \frac{1}{2}+i t\right)\right\rangle_{S \mathbb{X}} \eta\left(m, \frac{1}{2}+i t\right) \mu_{d}\left(E_{0}\left(\cdot, \frac{1}{2}+i t\right)\right) d t
\end{aligned}
$$
\]

Firstly, we have from (4-1), (4-2), and (4-3) that $\eta_{j}^{m}\left(\phi_{j, l}^{m}\right)$ and $\eta\left(m, \frac{1}{2}+i t\right)$ are both $O(1)$. Note by successive integration by parts and Cauchy-Schwarz inequality, we have for all $N \geq 1$,

$$
\left\langle F, \phi_{j, l}^{m}\right\rangle \ll_{N}\left(|l|^{2}+1\right)^{-N}\|F\|_{W^{2 N, 2}(S X)}
$$

when $m>0$, and

$$
\left\langle F, \phi_{j, l}^{0}\right\rangle \ll N_{N}\left(|l|^{2}+\left|t_{j}\right|^{2}+1\right)^{-N}\|F\|_{W^{2 N, 2}(S X)} .
$$

Likewise, assuming that the support of $F$ is contained in $y<T$, we have

$$
\left\langle F, E_{m}\left(\cdot, \frac{1}{2}+i t\right)\right\rangle_{S X}<_{N}\left(|m|^{2}+t^{2}+1\right)^{-N}\|F\|_{W^{2 N, 2}(S X)} \log T,
$$

where we used [Kubota 1973, (6.1.6)] and [Jakobson 1994, (1.6), (1.7)].
Now for $m>0$, we take $N=3$ and apply (4-5) to see that

$$
\sum_{\substack{m>0 \\ 4 \backslash m}} \sum_{j=1}^{d_{m}} \mu_{d}\left(\phi_{j}^{m}\right) \sum_{\substack{l \in 4 \mathbb{Z} \\|l| \geq m}}\left\langle F, \phi_{j, l}^{m}\right\rangle_{S X} \eta_{j}^{m}\left(\phi_{j, l}^{m}\right) \ll_{\epsilon} d^{1 / 2-25 / 512+\epsilon}\|F\|_{W^{6,2}(S X)},
$$

and for $m=0$, we take $N=2$ and apply (4-6) to deduce

$$
\sum_{j=1}^{\infty} \mu_{d}\left(\phi_{j}^{0}\right) \sum_{l \in 4 \mathbb{Z}}\left\langle F, \phi_{j, l}^{0}\right\rangle_{S X} \eta_{j}^{0}\left(\phi_{j, l}^{0}\right) \ll_{\epsilon} d^{1 / 2-25 / 512+\epsilon}\|F\|_{W^{4,2}(S X)}
$$

For the Eisenstein series contribution, we take $N=2$ and apply (4-9) to see

$$
\sum_{m \in 4 \mathbb{Z}} \frac{1}{4 \pi} \int_{-\infty}^{\infty}\left\langle F, E_{m}\left(\cdot, \frac{1}{2}+i t\right)\right\rangle_{S X} \eta\left(m, \frac{1}{2}+i t\right) \mu_{d}\left(E_{0}\left(\cdot, \frac{1}{2}+i t\right)\right) d t \ll_{\epsilon} \log T d^{7 / 16+\epsilon}\|F\|_{W^{4,2}(S X)}
$$

Therefore Theorem 1.3 will follow once we establish the following lower bound for the total length of $\mathscr{C}_{d}$ :

$$
\begin{equation*}
l\left(\mathscr{C}_{d}\right)=2 h(d) \log \epsilon_{d} \gg_{\epsilon} d^{1 / 2-\epsilon} . \tag{4-10}
\end{equation*}
$$

To see this, let $d=D b^{2}$ where $D$ is a fundamental discriminant. Then by Dirichlet class number formula [Davenport 1967, page 50] for binary quadratic forms discriminant $d$ (or by letting $s \rightarrow 1$ in (4-7)), we have

$$
h(d) \log \left(\epsilon_{d}\right)=d^{1 / 2} L(1, d)
$$

with the same $L(\cdot, d)$ given in (4-8), i.e.,

$$
L(1, d)=L\left(1, \chi_{D}\right)\left(\sum_{a \mid b} \mu(a)\left(\frac{D}{a}\right) a^{-1} \sigma_{-1}\left(\frac{b}{a}\right)\right)
$$

Note that

$$
\sum_{a \mid b} \mu(a)\left(\frac{D}{a}\right) a^{-1} \sigma_{-1}\left(\frac{b}{a}\right)=\sum_{c a \mid b} \mu(a)\left(\frac{D}{a}\right) \frac{c}{b}=\frac{1}{b} \sum_{e \mid b} e \prod_{p \mid e}\left(1-\left(\frac{D}{p}\right) p^{-1}\right),
$$

where $e=a c$, and that

$$
\frac{1}{b} \sum_{e \backslash b} e \prod_{p \mid e}\left(1-\left(\frac{D}{p}\right) p^{-1}\right) \gg b^{-\epsilon} .
$$

Now (4-10) follows by using Siegel's lower bound [1935]

$$
L\left(1, \chi_{D}\right) \ggg_{\epsilon} D^{-\epsilon},
$$

and this completes the proof of Theorem 1.3.
4D. Proof of Theorem 1.1. We are now ready to prove Theorem 1.1. Assume that $\beta:[0, l(\beta)] \rightarrow \mathbb{X}$ is a sufficiently short compact geodesic segment in the region determined by $y<T$ such that $\beta([-l(\beta), 2 l(\beta)])$ has no self intersection. (We fix $T$ for simplicity, but it is possible to vary $T$ with $d$.) For $\delta=d^{-a}$ with $a>0$ to be chosen later, such that $l(\beta) \gg \delta$, let

$$
\beta_{1}:=\{\beta(t): t \in[0, l(\beta)-\delta]\} \quad \text { and } \quad \beta_{2}:=\{\beta(t): t \in[-\delta, l(\beta)]\} .
$$

Then from Lemma 2.7, we have

$$
\frac{1}{\delta^{2}} \int_{\widetilde{\alpha}_{2}} \int_{\widetilde{\beta}_{1}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} \leq I^{\theta_{1}, \theta_{2}}\left(\beta, \alpha_{2}\right) \leq \frac{1}{\delta^{2}} \int_{\widetilde{\alpha}_{2}} \int_{\widetilde{\beta}_{1}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}
$$

for any closed geodesic $\alpha_{2}$. Now define $f_{1}, f_{2} \in C_{0}^{\infty}(S \mathbb{X})$ using Lemma 2.4 by

$$
f_{1}(g)=\frac{1}{\delta^{2}} \int_{\widetilde{\beta_{1}}} m_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}^{-1} g\right) d s_{1} \quad \text { and } \quad f_{2}(g)=\frac{1}{\delta^{2}} \int_{\widetilde{\beta_{2}}} M_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}^{-1} g\right) d s_{1},
$$

with $\varepsilon=d^{-2 a}$, where we assume that $\theta_{2}-\theta_{1} \gg d^{-a}$. Note that $m\left(g_{1}^{-1} g_{2}\right)$ and $M\left(g_{1}^{-1} g_{2}\right)$ are minorant and majorant of $K_{\delta}^{\theta_{1}, \theta_{2}}\left(g_{1}, g_{2}\right)$ for $g_{1} \in \beta_{i}, g_{2} \in S \mathbb{\text { for all sufficiently large } d \text { . Hence, for all sufficiently }}$ large $d$ (independent of $\alpha_{2}$ ), we have

$$
\int_{\widetilde{\alpha}_{2}} f_{1}(s) d s \leq I^{\theta_{1}, \theta_{2}}\left(\beta, \alpha_{2}\right) \leq \int_{\widetilde{\alpha}_{2}} f_{2}(s) d s,
$$

and so

$$
\begin{equation*}
\int_{\mathscr{C}_{d}} f_{1}(s) d s \leq 2 I^{\theta_{1}, \theta_{2}}\left(\beta, C_{d}\right) \leq \int_{\mathscr{C}_{d}} f_{2}(s) d s, \tag{4-11}
\end{equation*}
$$

where the factor 2 amounts to the fact that $\mathscr{C}_{d}$ is a double cover of $C_{d}$.

We now apply Theorem 1.3 to see that

$$
\frac{1}{l\left(\mathscr{C}_{d}\right)} \int_{\mathscr{C}_{d}} f_{i}(s) d s=\frac{3}{\pi^{2}} \int_{S X} f_{i}(g) d \mu_{g}+O_{\epsilon}\left(d^{-25 / 512+\epsilon}\left\|f_{i}\right\|_{W^{6, \infty}}\right) .
$$

Because of the choice of $f_{1}$ and $f_{2}$, we have

$$
\left\|f_{i}\right\|_{W^{6, \infty}} \ll \varepsilon^{-6} l(\beta) \ll d^{12 a} l(\beta)
$$

and

$$
\int_{S X} f_{i}(g) d \mu_{g}=\left(\cos \theta_{1}-\cos \theta_{2}\right)(l(\beta)+O(\delta))(1+O(\varepsilon))=\left(\cos \theta_{1}-\cos \theta_{2}\right) l(\beta)\left(1+O\left(d^{-2 a}\right)\right)
$$

by Lemma 2.4. Now we complete the proof of Theorem 1.1 for sufficiently short geodesic segments by choosing $a=\frac{25}{7168}$ and applying these estimates to (4-11). This then implies Theorem 1.1 for any geodesic segment of length $<1$ by dividing the segment into finitely many sufficiently short geodesic segments, and then applying Theorem 1.1 to each of them.

## 5. Selberg's pretrace Formula for $\mathrm{PSL}_{2}(\mathbb{R})$

Let $k \in C_{0}^{\infty}\left(\operatorname{PSL}_{2}(\mathbb{R})\right)$, and let $K$ be the integral kernel on $S \mathbb{X}$ defined by

$$
K\left(g_{1}, g_{2}\right)=\sum_{\gamma \in \Gamma} k\left(g_{1}, \gamma g_{2}\right),
$$

where $k\left(g_{1}, g_{2}\right)=k\left(g_{1}^{-1} g_{2}\right)$. The corresponding integral operator $T_{K}$ acts on $f \in L^{2}(S \mathbb{X})$ by

$$
T_{K}(f):=\int_{S \mathbb{X}} K\left(g_{1}, g_{2}\right) f\left(g_{2}\right) d g_{2}=\int_{\mathrm{PSL}_{2}(\mathbb{R})} k\left(g_{1}^{-1} g_{2}\right) f\left(g_{2}\right) d g_{2}
$$

It follows that $T_{K}(f) \in L^{2}(S X)$. In this section, we study the spectral expansion of $K$ in terms of the equivariant eigenfunctions of the Casimir operator, which are explicitly described in Section 3A. In other words, we derive Selberg's pretrace formula for $\mathrm{PSL}_{2}(\mathbb{Z}) \backslash \mathrm{PSL}_{2}(\mathbb{R})$.

5A. Cuspidal spectrum. In this section, we describe explicitly the spectrum of $T_{K}$ acting on the cuspidal subspace $L_{\text {cusp }}^{2}(S \mathbb{X})$. Let $R_{g}(f)(x)=f(x g)$ be the right regular action of $\mathrm{PSL}_{2}(\mathbb{R})$ on

$$
\left.L_{\text {cusp }}^{2}\left(\Gamma \backslash \mathrm{PSL}_{2}(\mathbb{R})\right)=L_{\text {cusp }}^{2}: S \mathbb{X}\right)
$$

Lemma 5.1. Let $\pi$ be an irreducible unitary representation of $\operatorname{PSL}_{2}(\mathbb{R})$. Then for any $f \in W_{\pi} \subset$ $L_{\text {cusp }}^{2}(S \rtimes)$, we have

$$
T_{K}(f) \in W_{\pi}
$$

Proof. Observe that

$$
T_{K}(f)\left(g_{1}\right)=\int_{\mathrm{PSL}_{2}(\mathbb{R})} k\left(g_{1}, g_{2}\right) f\left(g_{2}\right) d g_{2}=\int_{\mathrm{PSL}_{2}(\mathbb{R})} k\left(g_{1}^{-1} g_{2}\right) f\left(g_{2}\right) d g_{2}=\int_{\mathrm{PSL}_{2}(\mathbb{R})} k(u) f\left(g_{1} u\right) d u
$$

where $u=g_{1}^{-1} g_{2}$. Hence, we have

$$
T_{K}(f)=\int_{\operatorname{PSL}_{2}(\mathbb{R})} k(u) R_{u}(f) d u
$$

and because $R_{u}(f) \in W_{\pi}$ for every $u$, we conclude that $T_{K}(f) \in W_{\pi}$.
From 5 A , for an abstract irreducible unitary representation $\pi$ of $\mathrm{PSL}_{2}(\mathbb{R})$ and $f \in W_{\pi}$, we define the action of $k$ on $f$ by

$$
k * f=\int_{\mathrm{PSL}_{2}(\mathbb{R})} k(u) \pi(u)(f) d u
$$

which agrees with $T_{K}(f)$ when $W_{\pi}$ is a subspace of $L_{\text {cusp }}^{2}(S X)$.
Let $\psi: W_{\pi} \rightarrow W_{\pi^{\prime}}$ be an isomorphism of representations $\pi$ and $\pi^{\prime}$. Note that for $f \in W_{\pi}$ and $f^{\prime} \in W_{\pi^{\prime}}$ with $\psi(f)=f^{\prime}$, we have $\psi(k * f)=k * f^{\prime}$. We denote by $\phi_{m} \in W_{\pi}$ the unique (up to a unit scalar) vector of norm 1 and weight $m$. We fix the unit scalar except for the spherical or the lowest weight vector, by using the normalized lowering and raising operator that we introduced in (3-4) and (3-5).

Now let

$$
\begin{equation*}
h(k, m, n, \pi):=\left\langle k * \phi_{m}, \phi_{n}\right\rangle, \tag{5-1}
\end{equation*}
$$

and let $M_{\pi}(m, n)(g)=\left\langle\pi(g) \phi_{m}, \phi_{n}\right\rangle$ be the matrix coefficient of $\pi$. We note that $h(k, m, n, \pi)$ and $M_{\pi}(m, n)(g)$ do not depend on the choice of the unit scalar of the spherical or the lowest weight vector.

We recall some properties of $M_{\pi}(m, n)(g)$ in the following lemma.
Lemma 5.2. We have for every $g \in \operatorname{PSL}_{2}(\mathbb{R})$,

$$
\left|M_{\pi}(m, n)(g)\right| \leq 1,
$$

and

$$
M_{\pi}(m, n)\left(R_{\theta^{\prime}} g R_{\theta}\right)=e^{-i m \theta} e^{-i n \theta^{\prime}} M_{\pi}(m, n)(g)
$$

Proof. We have

$$
1=\left|\pi(g) \phi_{m}\right|^{2}=\sum_{n}\left\langle\pi(g) \phi_{m}, \phi_{n}\right\rangle^{2}
$$

from which it is immediate that $\left|M_{\pi}(m, n)(g)\right| \leq 1$. For the second identity, we have

$$
M_{\pi}(m, n)\left(R_{\theta^{\prime}} g R_{\theta}\right)=\left\langle\pi(g) \pi\left(R_{\theta}\right) \phi_{m}, \pi\left(R_{-\theta^{\prime}}\right) \phi_{n}\right\rangle=e^{-i m \theta} e^{-i n \theta^{\prime}} M_{\pi}(m, n)(g)
$$

Define $k_{m, n} \in C_{0}^{\infty}\left(\operatorname{PSL}_{2}(\mathbb{R})\right)$ by

$$
\begin{equation*}
k_{m, n}(g):=\frac{1}{4 \pi^{2}} \int_{0}^{2 \pi} \int_{0}^{2 \pi} k\left(R_{\theta^{\prime}} g R_{\theta}\right) e^{-i n \theta^{\prime}-i m \theta} d \theta^{\prime} d \theta \tag{5-2}
\end{equation*}
$$

Note that

$$
\begin{equation*}
k_{m, n}\left(R_{\theta_{1}} g R_{\theta_{2}}\right)=e^{i n \theta_{1}} k_{m, n}(g) e^{i m \theta_{2}} \tag{5-3}
\end{equation*}
$$

The following lemma holds for every unitary irreducible representation of $\mathrm{PSL}_{2}(\mathbb{R})$.

Lemma 5.3. We have

$$
h(k, m, n, \pi)=\int_{\mathrm{PSL}_{2}(\mathbb{R})} k_{m, n}(u) M_{\pi}(m, n)(u) d u
$$

and for all nonnegative integers $N_{1}, N_{2}$, we have the following estimate

$$
h(k, m, n, \pi) \lll N_{N=N_{1}+N_{2}}(1+|m|)^{-N_{1}}(1+|n|)^{-N_{2}}\|k\|_{W^{N, 1}} .
$$

Proof. Recall from the definition that

$$
h(k, m, n, \pi)=\int_{\mathrm{PSL}_{2}(\mathbb{R})} k(u)\left\langle\pi(u) \phi_{m}, \phi_{n}\right\rangle d u=\int_{\mathrm{PSL}_{2}(\mathbb{R})} k(u) M_{\pi}(m, n)(u) d u
$$

and so

$$
\begin{aligned}
h(k, m, n, \pi) & =\int_{\mathrm{PSL}_{2}(\mathbb{R})} k(u) M_{\pi}(m, n)(u) d u . \\
& =\frac{1}{4 \pi^{2}} \int_{\mathrm{PSL}_{2}(\mathbb{R})} \int_{\theta} \int_{\theta^{\prime}} k\left(R_{\theta^{\prime}} u R_{\theta}\right) M_{\pi}(m, n)\left(R_{\theta^{\prime}} u R_{\theta}\right) d \theta d \theta^{\prime} d u \\
& =\frac{1}{4 \pi^{2}} \int_{\mathrm{PSL}_{2}(\mathbb{R})} M_{\pi}(m, n)(u) \int_{\theta} \int_{\theta^{\prime}} k\left(R_{\theta^{\prime}} u R_{\theta}\right) e^{-i m \theta} e^{-i n \theta^{\prime}} d \theta d \theta^{\prime} d u \\
& =\int_{\mathrm{PSL}_{2}(\mathbb{R})} k_{m, n}(u) M_{\pi}(m, n)(u) d u .
\end{aligned}
$$

Therefore, by integration by parts, we have

$$
\begin{aligned}
h(k, m, n, \pi) & \leq \int_{\mathrm{PSL}_{2}(\mathbb{R})}\left|k_{m, n}(u)\right| d u \\
& =\int_{\mathrm{PSL}_{2}(\mathbb{R})}\left|\frac{1}{4 \pi^{2}} \int_{\theta} \int_{\theta^{\prime}} k\left(R_{\theta^{\prime}} u R_{\theta}\right) e^{-i m \theta} e^{-i n \theta^{\prime}} d \theta d \theta^{\prime}\right| d u \\
& \ll N_{N}(1+|m|)^{-N_{1}}(1+|n|)^{-N_{2}}\|k\|_{W^{N, 1}}
\end{aligned}
$$

where we used $\left|M_{\pi}(m, n)(u)\right| \leq 1$ from Lemma 5.2. This completes the proof of our lemma.
5A1. Principal series representation of $\mathrm{SL}_{2}(\mathbb{R})$. For our application in the subsequent chapters, we need a refined estimate for $h(k, m, n, \pi)$ when $\pi$ is a unitary principal series representation. We first give an explicit representation of $h(k, m, n, \pi)$.

Lemma 5.4. Let $W_{\pi}$ be a unitary principal series representation of $\mathrm{SL}_{2}(\mathbb{R})$ with the parameter $\frac{1}{2}+$ it [Knapp 2001, Chapter VII]. Let

$$
\begin{equation*}
h(k, m, n, t):=\int_{\operatorname{PSL}_{2}(\mathbb{R})} k_{m, n}(g) y^{1 / 2+i t} e^{-i m \theta} d g \tag{5-4}
\end{equation*}
$$

where $g=n a(y) R_{\theta}$. Then we have

$$
h(k, m, n, \pi)=h(k, m, n, t) .
$$

Proof. We note that principal series representations are induced from the unitary characters of the upper triangular matrices to $\mathrm{PSL}_{2}(\mathbb{R})$ [Knapp 2001, Chapter VII]. In this model, a dense subspace of a representation is given by

$$
\left\{f: \mathrm{PSL}_{2}(\mathbb{R}) \rightarrow \mathbb{C} \text { continuous : } f(\text { xan })=e^{(i t+1 / 2) \log (a)} f(x)\right\}
$$

with the norm

$$
|f|^{2}=\frac{1}{2 \pi} \int_{\theta}\left|f\left(R_{\theta}\right)\right|^{2} d \theta
$$

and the $\mathrm{PSL}_{2}(\mathbb{R})$ action is given by

$$
\pi(g) f(x)=f\left(g^{-1} x\right)
$$

The weight $m$ unit vectors are explicitly given by

$$
\phi_{m}\left(R_{\theta} a(y) n\right)=e^{i m \theta} y^{-(1 / 2+i t)}
$$

Note that the orthonormal basis $\left\{\phi_{m}\right\}$ is normalized as our convention in (3-4), i.e.,

$$
\boldsymbol{E}^{-} \phi_{m}=(m+1-2 i t) \phi_{m-2} \quad \text { and } \quad \boldsymbol{E}^{+} \phi_{m}=(m+1+2 i t) \phi_{m+2}
$$

With these, we first see that

$$
\begin{aligned}
k * \phi_{m}\left(R_{\theta^{\prime}}\right) & =\int_{\mathrm{PSL}_{2}(\mathbb{R})} k(u) y\left(u^{-1} R_{\theta^{\prime}}\right)^{-(1 / 2+i t)} e^{i m \theta\left(u^{-1} R_{\left.\theta^{\prime}\right)}\right.} d u \\
& =\int_{\mathrm{PSL}_{2}(\mathbb{R})} k\left(R_{\theta^{\prime}} v^{-1}\right) y(v)^{-(1 / 2+i t)} e^{i m \theta(v)} d v,
\end{aligned}
$$

where $v=u^{-1} R_{\theta^{\prime}}$ and $v=R_{\theta(v)} a(y(v)) n(v)$. We therefore have

$$
\begin{aligned}
h(k, m, n, \pi) & =\left\langle k * f_{m}, f_{n}\right\rangle \\
& =\frac{1}{2 \pi} \int_{\theta^{\prime}} k * f_{m}\left(R_{\theta^{\prime}}\right) \bar{f}_{n}\left(R_{\theta^{\prime}}\right) d \theta^{\prime} \\
& =\frac{1}{2 \pi} \int_{\theta^{\prime}} e^{-i n \theta^{\prime}} \int_{\mathrm{PSL}_{2}(\mathbb{R})} k\left(R_{\theta^{\prime}} v^{-1}\right) y(v)^{-(1 / 2+i t)} e^{i m \theta(v)} d v d \theta^{\prime} \\
& =\frac{1}{2 \pi} \int_{\mathrm{PSL}_{2}(\mathbb{R})} y^{1 / 2+i t} \int_{\theta^{\prime}} e^{-i n \theta^{\prime}} e^{-i m \theta} k\left(R_{\theta^{\prime}} w\right) d \theta^{\prime} d w \\
& =\int_{\mathrm{PSL}_{2}(\mathbb{R})} k_{m, n}(w) y^{1 / 2+i t} e^{-i m \theta} d w
\end{aligned}
$$

where $w=v^{-1}$ and $w=n a(y) R_{\theta}$. Note that $y=y(v)^{-1}$ and $\theta=-\theta(v)$.
We now prove that $h(k, m, n, t)$ decays fast in all parameters uniformly.
Lemma 5.5. Suppose that $k$ is supported inside the compact subset $C \subset \mathrm{SL}_{2}(\mathbb{R})$. Then we have

$$
\int_{\mathrm{PSL}_{2}(\mathbb{R})} k_{m, n}(g) y^{1 / 2+i t} e^{-i m \theta} d g<_{N, C}(1+|m|)^{-N_{1}}(1+|n|)^{-N_{2}}(1+|t|)^{-N_{3}}\|k\|_{W^{N, \infty}}
$$

for any $N_{1}, N_{2}, N_{3} \geq 0$, where $N=N_{1}+N_{2}+N_{3}$.

Proof. From the definition, we have

$$
\int_{\mathrm{PSL}_{2}(\mathbb{R})} k_{m, n}(g) y^{1 / 2+i t} e^{-i m \theta} d g=\frac{1}{4 \pi} \int_{\mathbb{H}} \int_{0}^{2 \pi} \int_{0}^{2 \pi} k\left(R_{\theta_{1}^{\prime}} n(x) a(y) R_{\theta_{2}^{\prime}}\right) y^{1 / 2+i t} e^{-i n \theta_{1}^{\prime}-i m \theta_{2}^{\prime}} d \theta_{1}^{\prime} d \theta_{2}^{\prime} \frac{d x d y}{y^{2}},
$$

and so the statement follows from integration by parts.
5B. Continuous spectrum. For $k_{m, n}$ given by (5-2), let

$$
\begin{equation*}
K_{m, n}\left(g_{1}, g_{2}\right):=\sum_{\gamma \in \Gamma} k_{m, n}\left(g_{1}^{-1} \gamma g_{2}\right) . \tag{5-5}
\end{equation*}
$$

Then we infer from (5-3) that

$$
K_{m, n}\left(g_{1} R_{\theta_{1}}, g_{2} R_{\theta_{2}}\right)=e^{-i n \theta_{1}} K_{m, n}\left(g_{1}, g_{2}\right) e^{i m \theta_{2}}
$$

and so it defines an integral operator that maps weight $m$ forms to weight $n$ forms. Denote by $S^{m} \subset$ $L^{2}\left(\Gamma \backslash \operatorname{PSL}_{2}(\mathbb{R})\right)$ the space of weight $m$ forms and by $S_{\text {cusp }}^{m}$ the space of weight $m$ forms in $L_{\text {cusp }}^{2}(S \mathbb{X})$. We first recall the following result regarding the decomposition of $K_{m, m}$.

Theorem 5.6 [Hejhal 1976]. The integral kernel

$$
K_{m, m}\left(g_{1}, g_{2}\right)-\frac{1}{4 \pi} \int_{-\infty}^{\infty} h(k, m, m, t) E_{m}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{m}\left(g_{2}, \frac{1}{2}+i t\right)} d t
$$

defines a compact operator $S_{\mathrm{cusp}}^{m} \rightarrow S_{\mathrm{cusp}}^{m}$ that acts trivially on $\Theta$. (Here $h(k, m, m, t)$ is given by (5-4).)
We define $\boldsymbol{E}^{a}$ to be $\left(\boldsymbol{E}^{+}\right)^{a}$ if $a>0$, and $\left(\boldsymbol{E}^{-}\right)^{|a|}$ if $a<0$. We have

$$
\overline{\boldsymbol{E}^{a}}=(-\boldsymbol{E})^{-a}
$$

which follows directly from (3-2). Let $c_{m, n}$ be given by

$$
\boldsymbol{E}^{n-m} E_{m}(g, s)=c_{m, n}(s) E_{n}(g, s)
$$

Observe that

$$
\boldsymbol{E}^{n-m} y^{s} e^{-i m \theta}=c_{m, n}(s) y^{s} e^{-i n \theta}
$$

and that

$$
\begin{equation*}
\overline{c_{m, n}\left(\frac{1}{2}+i t\right)}=c_{n, m}\left(\frac{1}{2}+i t\right) \tag{5-6}
\end{equation*}
$$

for $t \in \mathbb{R}$.
Theorem 5.7. For $m, n \in 2 \mathbb{Z}$,

$$
K_{m, n}\left(g_{1}, g_{2}\right)-\frac{1}{4 \pi} \int_{-\infty}^{\infty} h(k, m, n, t) E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{m}\left(g_{2}, \frac{1}{2}+i t\right)} d t
$$

defines a compact operator $S_{\text {cusp }}^{m} \rightarrow S_{\text {cusp }}^{n}$ that acts trivially on $\Theta$.

Proof. Note that

$$
\int \boldsymbol{E}_{g_{2}}^{m-n}\left(K\left(g_{1}, g_{2}\right) f\left(g_{2}\right)\right) d g_{2}=0
$$

for every $g_{1}, m \neq n$, and $f \in C_{0}^{\infty}\left(\Gamma \backslash \operatorname{PSL}_{2}(\mathbb{R})\right)$. Hence

$$
T_{K} \boldsymbol{E}^{m-n}: C_{0}^{\infty}\left(\Gamma \backslash \mathrm{PSL}_{2}(\mathbb{R})\right) \rightarrow C_{0}^{\infty}\left(\Gamma \backslash \mathrm{PSL}_{2}(\mathbb{R})\right)
$$

is an integral operator with the integral kernel

$$
K^{\prime}\left(g_{1}, g_{2}\right)=\sum_{\gamma \in \Gamma} k^{\prime}\left(g_{1}^{-1} \gamma g_{2}\right)
$$

where

$$
k^{\prime}(g)=(-\boldsymbol{E})^{m-n} k(g)=\overline{\boldsymbol{E}^{n-m}} k(g) .
$$

Then by Theorem 5.6, we see that

$$
K^{\prime \prime}\left(g_{1}, g_{2}\right)=K_{n, n}^{\prime}\left(g_{1}, g_{2}\right)-\frac{1}{4 \pi} \int_{-\infty}^{\infty} h\left(k^{\prime}, n, n, t\right) E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{n}\left(g_{2}, \frac{1}{2}+i t\right)} d t
$$

defines a compact operator $T_{K^{\prime \prime}}: S_{\text {cusp }}^{n} \rightarrow S_{\text {cusp }}^{n}$ that acts trivially on $\Theta$. Note that

$$
\begin{aligned}
\int_{-\infty}^{\infty} h\left(k^{\prime}, n, n, t\right) E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{n}\left(g_{2}, \frac{1}{2}+i t\right)} d t & \\
& =\int_{-\infty}^{\infty} \frac{h\left(k^{\prime}, n, n, t\right)}{\overline{c_{m, n}(1 / 2+i t)}} E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{\boldsymbol{E}^{n-m} E_{m}\left(g_{2}, \frac{1}{2}+i t\right)} d t
\end{aligned}
$$

Let

$$
K^{\prime \prime \prime}\left(g_{1}, g_{2}\right):=K_{m, n}\left(g_{1}, g_{2}\right)-\frac{1}{4 \pi} \int_{-\infty}^{\infty} \frac{h\left(k^{\prime}, n, n, t\right)}{\overline{c_{m, n}(1 / 2+i t)}} E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{m}\left(g_{2}, \frac{1}{2}+i t\right)} d t
$$

Note that

$$
T_{K^{\prime \prime}}=T_{K^{\prime \prime \prime}} \circ \boldsymbol{E}^{m-n} .
$$

Firstly, since $\boldsymbol{E}^{m-n}$ does not annihilate the Eisenstein series, $T_{K^{\prime \prime \prime}}$ acts trivially on $\Theta$.
If $m>n \geq 0$ or $m<n \leq 0$, then as a map $S_{\text {cusp }}^{n} \rightarrow S_{\text {cusp }}^{m}, \operatorname{ker}\left(\boldsymbol{E}^{m-n}\right)$ is empty, and we may decompose $S_{\text {cusp }}^{m}$ as

$$
S_{\text {cusp }}^{m}=\Im\left(\boldsymbol{E}^{m-n}\right) \oplus R
$$

where $R$ is a finite dimensional subspace of $S_{\text {cusp }}^{m}$ spanned by modular forms of weight $>n$ and their images under raising operators in $S_{\text {cusp }}^{m}$. Note that

$$
\left(\boldsymbol{E}^{m-n}\right)^{-1}: \Im\left(\boldsymbol{E}^{m-n}\right) \rightarrow S_{\text {cusp }}^{n}
$$

is a bounded operator, hence

$$
\left.T_{K^{\prime \prime \prime}}\right|_{I m\left(\boldsymbol{E}^{m-n}\right)}=T_{K^{\prime \prime}} \circ\left(\boldsymbol{E}^{m-n}\right)^{-1}
$$

is a compact operator. This implies that $T_{K^{\prime \prime \prime}}$ is a direct sum of a compact operator and finite dimensional linear operator, which is a compact operator.

If $n>m \geq 0$ or $n<m \leq 0$, then $\boldsymbol{E}^{m-n}: S_{\text {cusp }}^{n} \rightarrow S_{\text {cusp }}^{m}$ is surjective, and so we may define a bounded operator

$$
\left(\boldsymbol{E}^{m-n}\right)^{-1}: S_{m} \rightarrow\left(\operatorname{ker}\left(\boldsymbol{E}^{m-n}\right)\right)^{\perp}
$$

from which it follows that

$$
T_{K^{\prime \prime \prime}}=T_{K^{\prime \prime}} \circ\left(\boldsymbol{E}^{m-n}\right)^{-1}
$$

is a compact operator.
If $n>0>m$ or $m>0>n$, then we further decompose $T_{K^{\prime \prime}}$ to

$$
S_{\text {cusp }}^{n} \xrightarrow{\boldsymbol{E}^{-n}} S_{\text {cusp }}^{0} \xrightarrow{E^{m}} S_{\text {cusp }}^{m} \xrightarrow{T_{K^{\prime \prime \prime}}} S_{\text {cusp }}^{n},
$$

and then combine the above arguments to see that $T_{K^{\prime \prime}}$ is a compact operator.
Finally, observe that

$$
h\left(k^{\prime}, n, n, t\right)=\int_{\mathrm{PSL}_{2}(\mathbb{R})}\left(\overline{\boldsymbol{E}^{n-m}} k(g)\right) y^{\frac{1}{2}+i t} e^{i n \theta} d g=c_{n, m}\left(\frac{1}{2}+i t\right) \int_{\mathrm{PSL}_{2}(\mathbb{R})} k(g) y^{\frac{1}{2}+i t} e^{i m \theta} d g
$$

and we complete the proof using (5-6).
5C. General case. We are now ready to describe Selberg's pretrace formula for $\mathrm{PSL}_{2}(\mathbb{R})$.
Theorem 5.8. For $k \in C_{0}^{\infty}\left(\operatorname{PSL}_{2}(\mathbb{R})\right)$, let $K$ be the integral kernel on $S \mathbb{X}$ defined by

$$
K\left(g_{1}, g_{2}\right)=\sum_{\gamma \in \Gamma} k\left(g_{1}, \gamma g_{2}\right) .
$$

Then we have

$$
\begin{aligned}
K\left(g_{1}, g_{2}\right)=\frac{9}{\pi^{4}} \iint K\left(g_{1}, g_{2}\right) d g_{1} d g_{2} & +\sum_{\substack{e \geq 0 \\
2\lceil e}} \sum_{j=1}^{d_{e}} \sum_{\substack{m, n \in 2 \mathbb{Z} \\
|m|,|n| \geq e}} h\left(k, m, n, \pi_{j}^{e}\right) \phi_{j, n}^{e}\left(g_{1}\right) \overline{\phi_{j, m}^{e}\left(g_{2}\right)} \\
& +\frac{1}{4 \pi} \sum_{m, n \in 2 \mathbb{Z}} \int_{-\infty}^{\infty} h(k, m, n, t) E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{m}\left(g_{2}, \frac{1}{2}+i t\right)} d t
\end{aligned}
$$

where $\pi_{j}^{e}$ is the irreducible unitary representation of $\operatorname{PSL}_{2}(\mathbb{R})$ associated to $\phi_{j}^{e}$.

Proof. We first note from (5-2) and (5-5) that

$$
\begin{aligned}
K_{m, n}\left(g_{1}, g_{2}\right) & =\sum_{\gamma \in \Gamma} \frac{1}{4 \pi^{2}} \int_{0}^{2 \pi} \int_{0}^{2 \pi} k\left(R_{\theta_{1}^{\prime}} g_{1}^{-1} \gamma g_{2} R_{\theta_{2}^{\prime}}\right) e^{-i n \theta_{1}^{\prime}-i m \theta_{2}^{\prime}} d \theta_{1}^{\prime} d \theta_{2}^{\prime} \\
& =\frac{1}{4 \pi^{2}} \int_{0}^{2 \pi} \int_{0}^{2 \pi} \sum_{\gamma \in \Gamma} k\left(R_{-\theta_{1}^{\prime}} g_{1}^{-1} \gamma g_{2} R_{\theta_{2}^{\prime}}\right) e^{i n \theta_{1}^{\prime}-i m \theta_{2}^{\prime}} d \theta_{1}^{\prime} d \theta_{2}^{\prime} \\
& =\frac{1}{4 \pi^{2}} \int_{0}^{2 \pi} \int_{0}^{2 \pi} K\left(g_{1} R_{\theta_{1}^{\prime}}, g_{2} R_{\theta_{2}^{\prime}}\right) e^{i n \theta_{1}^{\prime}-i m \theta_{2}^{\prime}} d \theta_{1}^{\prime} d \theta_{2}^{\prime} \\
& =\frac{1}{4 \pi^{2}} \int_{0}^{2 \pi} \int_{0}^{2 \pi} K\left(\left(x_{1}, y_{1}, \theta_{1}^{\prime}\right),\left(x_{2}, y_{2}, \theta_{2}^{\prime}\right)\right) e^{i n \theta_{1}^{\prime}-i m \theta_{2}^{\prime}} d \theta_{1}^{\prime} d \theta_{2}^{\prime} e^{-i n \theta_{1}+i m \theta_{2}} .
\end{aligned}
$$

Therefore, we have the Fourier expansion of $K$,

$$
K\left(g_{1}, g_{2}\right)=\sum_{n, m \in 2 \mathbb{Z}} K_{m, n}\left(g_{1}, g_{2}\right)
$$

where the summation is uniform for $g_{1}$ and $g_{2}$ in compacta.
We infer from Theorem 5.7 that

$$
K_{m, n}\left(g_{1}, g_{2}\right)-\frac{1}{4 \pi} \int_{-\infty}^{\infty} h(k, m, n, t) E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{m}\left(g_{2}, \frac{1}{2}+i t\right)} d t
$$

defines a compact operator acting on $L_{\text {cusp }}$ that acts trivially on $\Theta$. Because it only acts nontrivially on weight $m$ forms, we see that

$$
\begin{aligned}
& K_{m, n}\left(g_{1}, g_{2}\right)-\frac{1}{4 \pi} \int_{-\infty}^{\infty} h(k, m, n, t) E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{m}\left(g_{2}, \frac{1}{2}+i t\right)} d t \\
&=\frac{9}{\pi^{4}} \iint K_{m, n}\left(g_{1}, g_{2}\right) d g_{1} d g_{2}+\sum_{\substack{e \geq 0 \\
2\rceil e}}^{\min \{|m|,|n|\}} \sum_{j=1}^{d_{e}} h\left(k, m, n, \pi_{j}^{e}\right) \phi_{j, n}^{e}\left(g_{1}\right) \overline{\phi_{j, m}^{e}\left(g_{2}\right)},
\end{aligned}
$$

where we used (5-1), and the fact that

$$
\int_{-\infty}^{\infty} h(k, m, n, t) E_{n}\left(g_{1}, \frac{1}{2}+i t\right) \overline{E_{m}\left(g_{2}, \frac{1}{2}+i t\right)} d t
$$

acts trivially on $L_{\text {cusp }}^{2}$. Note that the integral on the right-hand side of the equation vanishes unless $m=n=0$, in which case it is identical to

$$
\frac{9}{\pi^{4}} \iint K\left(g_{1}, g_{2}\right) d g_{1} d g_{2}
$$

5D. Proof of Theorem 1.4. We now present a proof of Theorem 1.4. By Theorem 5.8, we have

$$
\frac{1}{l\left(\mathscr{C}_{d_{1}}\right) l\left(\mathscr{C}_{d_{2}}\right)} \int_{\mathscr{C}_{d_{2}}} \int_{\mathscr{C}_{d_{1}}} K\left(s_{1}, s_{2}\right) d s_{1} d s_{2}=M+D+\frac{1}{4 \pi} E
$$

where

$$
\begin{aligned}
M & =\frac{9}{\pi^{4}} \iint K\left(g_{1}, g_{2}\right) d g_{1} d g_{2}, \\
D & =\sum_{\substack{e \geq 0 \\
2 \mid e}} \sum_{j=1}^{d_{e}} \sum_{\substack{d_{n}, n \in 2 \mathbb{Z} \\
|m|,|n| \geq e}} h\left(k, m, n, \pi_{j}^{e}\right) \frac{\mu_{d_{1}}\left(\phi_{j, n}^{e}\right)}{l\left(\mathscr{C}_{d_{1}}\right)} \frac{\overline{\mu_{d_{2}}\left(\phi_{j, m}^{e}\right)}}{l\left(\mathscr{C}_{d_{2}}\right)} \\
& =\sum_{\substack{e \geq 0 \\
4 \mid e}} \sum_{j=1}^{d_{e}} \frac{\mu_{d_{1}}\left(\phi_{j}^{e}\right)}{l\left(\mathscr{C}_{d_{1}}\right)} \frac{\overline{\mu_{d_{2}}\left(\phi_{j}^{e}\right)}}{l\left(\mathscr{C}_{d_{2}}\right)} \sum_{\substack{m, n \in 4 \mathbb{Z} \\
|m|,|n| \geq e}} h\left(k, m, n, \pi_{j}^{e}\right) \eta_{j}^{e}\left(\phi_{j, n}^{e}\right) \overline{\eta_{j}^{e}\left(\phi_{j, m}^{e}\right)},
\end{aligned}
$$

and

$$
\begin{aligned}
E & =\sum_{m, n \in 2 \mathbb{Z}} \int_{-\infty}^{\infty} h(k, m, n, t) \frac{\mu_{d_{1}}\left(E_{n}(\cdot, 1 / 2+i t)\right)}{l\left(\mathscr{C}_{d_{1}}\right)} \frac{\overline{\mu_{d_{2}}\left(E_{m}(\cdot, 1 / 2+i t)\right)}}{l\left(\mathscr{C}_{d_{2}}\right)} d t \\
& =\sum_{m, n \in 4 \mathbb{Z}} \int_{-\infty}^{\infty} h(k, m, n, t) \frac{\mu_{d_{1}}\left(E_{0}(\cdot, 1 / 2+i t)\right)}{l\left(\mathscr{C}_{d_{1}}\right)} \frac{\overline{\mu_{d_{2}}\left(E_{0}(\cdot, 1 / 2+i t)\right)}}{l\left(\mathscr{C}_{d_{2}}\right)} \eta\left(n, \frac{1}{2}+i t\right) \overline{\eta\left(m, \frac{1}{2}+i t\right)} d t .
\end{aligned}
$$

For $D$ with $e>0$, we use (4-2), (4-5), Lemma 5.3 with $N_{1}=N_{2}=5$, and (4-10) to see that

$$
\begin{aligned}
\sum_{\substack{e>0 \\
4 \mid e}} \sum_{j=1}^{d_{e}} \frac{\mu_{d_{1}}\left(\phi_{j}^{e}\right)}{l\left(\mathscr{C}_{d_{1}}\right)} \frac{\overline{\mu_{d_{2}}\left(\phi_{j}^{e}\right)}}{l\left(\mathscr{C}_{d_{2}}\right)} & \sum_{\substack{m, n \in 4 \mathbb{Z} \\
|m|,|n| \geq e}} h\left(k, m, n, \pi_{j}^{e}\right) \eta_{j}^{e}\left(\phi_{j, n}^{e}\right) \overline{\eta_{j}^{e}\left(\phi_{j, m}^{e}\right)} \\
& <_{\epsilon} \sum_{\substack{e>0 \\
4 \mid e}} e^{6.8}\left(d_{1} d_{2}\right)^{-25 / 512+\epsilon} \sum_{\substack{m, n \in 4 \mathbb{Z} \\
|m|,|n| \geq e}}|m|^{-5}|n|^{-5}\|k\|_{W^{10, \infty}} \\
& \ll\left(d_{1} d_{2}\right)^{-25 / 512+\epsilon}\|k\|_{W^{10, \infty}} .
\end{aligned}
$$

For $D$ with $e=0$, we use (4-1), (4-6), Lemma 5.5 with $N_{1}=N_{2}=2$ and $N_{3}=4$, and (4-10) to see that

$$
\begin{aligned}
\sum_{j=1}^{\infty} \frac{\mu_{d_{1}}\left(\phi_{j}^{0}\right)}{l\left(\mathscr{C}_{d_{1}}\right)} \overline{\frac{\mu_{d_{2}}\left(\phi_{j}^{0}\right)}{l\left(\mathscr{C}_{d_{2}}\right)}} & \sum_{m, n \in 4 \mathbb{Z}} h\left(k, m, n, \pi_{j}^{0}\right) \eta_{j}^{0}\left(\phi_{j, n}^{0}\right) \overline{\eta_{j}^{0}\left(\phi_{j, m}^{0}\right)} \\
& \ll{ }_{\epsilon} \sum_{j=1}^{\infty}\left(d_{1} d_{2}\right)^{-25 / 512+\epsilon}\left|t_{j}\right|^{3 / 2} \sum_{m, n \in 4 \mathbb{Z}}(1+|m|)^{-2}(1+|n|)^{-2}\left(1+\left|t_{j}\right|\right)^{-4}\|k\|_{W^{8, \infty}} \\
& \ll\left(d_{1} d_{2}\right)^{-25 / 512+\epsilon}\|k\|_{W^{8, \infty}} .
\end{aligned}
$$

For $E$, we use (4-3), (4-9), Lemma 5.5 with $N_{1}=N_{2}=2$ and $N_{3}=3$, and (4-10) to see that

$$
\begin{gathered}
\sum_{m, n \in 4 \mathbb{Z}} \int_{-\infty}^{\infty} h(k, m, n, t) \frac{\mu_{d_{1}}\left(E_{0}(\cdot, 1 / 2+i t)\right)}{l\left(\mathscr{C}_{d_{1}}\right)} \frac{\overline{\mu_{d_{2}}\left(E_{0}(\cdot, 1 / 2+i t)\right)}}{l\left(\mathscr{C}_{d_{2}}\right)} \eta\left(n, \frac{1}{2}+i t\right) \overline{\eta\left(m, \frac{1}{2}+i t\right)} d t \\
<_{\epsilon} \sum_{m, n \in 4 \mathbb{Z}} \int_{-\infty}^{\infty}\left(d_{1} d_{2}\right)^{-1 / 16+\epsilon}(1+|m|)^{-2}(1+|n|)^{-2}(|t|+1)^{-2}\|k\|_{W^{7, \infty}} d t \\
\ll\left(d_{1} d_{2}\right)^{-1 / 16+\epsilon}\|k\|_{W^{7}, \infty} .
\end{gathered}
$$

Now observe that

$$
\iint K\left(g_{1}, g_{2}\right) d g_{1} d g_{2}=\int_{S X} \int_{S H} k\left(g_{1}^{-1} g_{2}\right) d g_{2} d g_{1}=\frac{\pi^{2}}{3} \int_{S \sharp} k(g) d g,
$$

and so

$$
M=\frac{3}{\pi^{2}} \int_{S \bigoplus} k(g) d g .
$$

So far, we proved the following:
Theorem 5.9. For any $k \in C_{0}^{\infty}(S \sharp)$, we have

$$
\frac{1}{l\left(\mathscr{C}_{d_{1}}\right) l\left(\mathscr{C}_{d_{2}}\right)} \int_{\mathscr{C}_{d_{2}}} \int_{\mathscr{C}_{d_{1}}} K\left(s_{1}, s_{2}\right) d s_{1} d s_{2}=\frac{3}{\pi^{2}} \int_{S H} k(g) d g+O_{\epsilon}\left(\left(d_{1} d_{2}\right)^{-25 / 512+\epsilon}\|k\|_{W^{10, \infty}}\right) .
$$

Remark 5.1. Note that this is not the same as equidistribution of $\mathscr{C}_{d_{1}} \times \mathscr{C}_{d_{2}}$ in $S \mathbb{X} \times \mathbb{}$. For instance, if we replace $K$ with any compactly supported smooth function in $S \mathbb{X} \times S \mathbb{X}$, then the equality may not hold when $d_{1}$ is fixed and $d_{2}$ tends to $\infty$.

In order to prove Theorem 1.4, we make specific choices of $k$ in Theorem 5.9. We let $K_{1}$ and $K_{2}$ to be the kernel corresponding to $k=m_{\delta}^{\theta_{1}, \theta_{2}}$ and $k=M_{\delta}^{\theta_{1}, \theta_{2}}$ defined in Lemma 2.4, respectively. Then by Lemma 2.4, we have

$$
\begin{aligned}
\frac{1}{l\left(\mathscr{C}_{d_{1}}\right) l\left(\mathscr{C}_{d_{2}}\right)} \int_{\mathscr{C}_{d_{2}}} \int_{\mathscr{C}_{d_{1}}} K_{1}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} & \leq \frac{1}{l\left(\mathscr{C}_{d_{1}}\right) l\left(\mathscr{C}_{d_{2}}\right)} \int_{\mathscr{C}_{d_{2}}} \int_{\mathscr{C}_{d_{1}}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2} \\
& \leq \frac{1}{l\left(\mathscr{C}_{d_{1}}\right) l\left(\mathscr{C}_{d_{2}}\right)} \int_{\mathscr{C}_{d_{2}}} \int_{\mathscr{C}_{d_{1}}} K_{2}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}
\end{aligned}
$$

while we know from Lemma 2.6 that

$$
\int_{\mathscr{C}_{d_{2}}} \int_{\mathscr{C}_{d_{1}}} K_{\delta}^{\theta_{1}, \theta_{2}}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}=4 \delta^{2} I_{\theta_{1}, \theta_{2}}\left(C_{d_{1}}, C_{d_{2}}\right)
$$

We now apply Theorem 5.9 and Lemma 2.4 to see that

$$
\frac{1}{l\left(\mathscr{C}_{d_{1}}\right) l\left(\mathscr{C}_{d_{2}}\right)} \int_{\mathscr{C}_{d_{2}}} \int_{\mathscr{C}_{d_{1}}} K_{i}\left(s_{1}, s_{2}\right) d s_{1} d s_{2}=\frac{3}{\pi^{2}}\left(\cos \theta_{1}-\cos \theta_{2}\right) \delta^{2}(1+O(\varepsilon))+O_{\epsilon}\left(\left(d_{1} d_{2}\right)^{-25 / 512+\epsilon} \varepsilon^{-10}\right)
$$

Therefore, we have

$$
\frac{I_{\theta_{1}, \theta_{2}}\left(C_{d_{1}}, C_{d_{2}}\right)}{l\left(C_{d_{1}}\right) l\left(C_{d_{2}}\right)}=\frac{3}{\pi^{2}}\left(\cos \theta_{1}-\cos \theta_{2}\right)\left(1+O\left(\delta^{2}\right)\right)(1+O(\varepsilon))+O_{\epsilon}\left(\left(d_{1} d_{2}\right)^{-25 / 512+\epsilon} \varepsilon^{-10} \delta^{-2}\right)
$$

and by choosing $\delta^{2}=\varepsilon=\left(d_{1} d_{2}\right)^{-25 / 6144}$, we complete the proof of Theorem 1.4.

## Appendix: Jacobian computation

Recall that $\Psi: A K A \rightarrow \mathrm{SL}_{2}(\mathbb{R})$ is given by

$$
\left(t_{1}, \varphi, t_{2}\right) \mapsto\left(\begin{array}{cc}
e^{t_{1} / 2} & 0 \\
0 & e^{-t_{1} / 2}
\end{array}\right) R_{\frac{\varphi}{2}}\left(\begin{array}{cc}
e^{-t_{2} / 2} & 0 \\
0 & e^{t_{2} / 2}
\end{array}\right)=\left(\begin{array}{cc}
e^{\left(t_{1}-t_{2}\right) / 2} \cos \frac{\varphi}{2} & -e^{\left(t_{1}+t_{2}\right) / 2} \sin \frac{\varphi}{2} \\
e^{\left(-t_{1}-t_{2}\right) / 2} \sin \frac{\varphi}{2} & e^{\left(t_{2}-t_{1}\right) / 2} \cos \frac{\varphi}{2}
\end{array}\right)
$$

In this section, we compute the pullback of $d V=d x d y d \theta / y^{2}$ under $\Psi$. We start with the identity

$$
\left(\begin{array}{cc}
e^{\left(t_{1}-t_{2}\right) / 2} \cos \frac{\varphi}{2} & -e^{\left(t_{1}+t_{2}\right) / 2} \sin \frac{\varphi}{2} \\
e^{\left(-t_{1}-t_{2}\right) / 2} \sin \frac{\varphi}{2} & e^{\left(t_{2}-t_{1}\right) / 2} \cos \frac{\varphi}{2}
\end{array}\right)=n(x) a(y) R_{\theta}=\left(\begin{array}{cc}
* & * \\
\frac{\sin \theta}{\sqrt{y}} & \frac{\cos \theta}{\sqrt{y}}
\end{array}\right)
$$

By comparing the image of $i \in \mathbb{H}$, we have

$$
x+i y=\frac{e^{\left(t_{1}-t_{2}\right) / 2} \cos \frac{\varphi}{2} i-e^{\left(t_{1}+t_{2}\right) / 2} \sin \frac{\varphi}{2}}{e^{\left(-t_{1}-t_{2}\right) / 2} \sin \frac{\varphi}{2} i+e^{\left(t_{2}-t_{1}\right) / 2} \cos \frac{\varphi}{2}}
$$

and for simplicity, we write this as $\frac{A}{B}$. By comparing the second row of each matrix, we have

$$
\frac{e^{i \theta}}{\sqrt{y}}=B
$$

From a quick computation, we see that
$A_{t_{1}}=\frac{A}{2}, \quad B_{t_{1}}=-\frac{B}{2}$,
$A_{t_{2}}=\frac{\bar{A}}{2}$,
$B_{t_{2}}=\frac{\bar{B}}{2}$,
$A_{\varphi}=-\frac{e^{t_{1}}}{2} B, \quad B_{\varphi}=\frac{e^{-t_{1}}}{2} A$,
$\mathfrak{\Im} A \bar{B}=1$,
$y=\frac{1}{|B|^{2}}$.

We use these to express the Jacobian matrix in terms of $A$ and $B$ as follows:

$$
\frac{\partial(x, y, \theta)}{\partial\left(t_{1}, t_{2}, \varphi\right)}=\left(\begin{array}{ccc}
\operatorname{Re} \frac{A}{B} & \Im \frac{1}{B^{2}} & \operatorname{Re}\left(-\frac{e^{t_{1}}}{2}-\frac{e^{-t_{1}}}{2} \frac{A^{2}}{B^{2}}\right) \\
\Im \frac{A}{B} & -\operatorname{Re} \frac{1}{B^{2}} & \Im\left(-\frac{e^{t_{1}}}{2}-\frac{e^{-t_{1}}}{2} \frac{A^{2}}{B^{2}}\right) \\
0 & \frac{1}{2} \Im \frac{B}{B} & \frac{e^{-t_{1}}}{2|B|^{2}}
\end{array}\right)
$$

From this, we have

$$
\begin{aligned}
\frac{1}{y^{2}}\left|\frac{\partial(x, y, \theta)}{\partial\left(t_{1}, t_{2}, \varphi\right)}\right| & =|B|^{4}\left|\frac{\partial(x, y, \theta)}{\partial\left(t_{1}, t_{2}, \varphi\right)}\right| \\
& =\left|-\frac{1}{2} e^{-t_{1}} \operatorname{Re}\left(\frac{\bar{A}}{B}\right)+\frac{1}{4} \Im\left(\bar{B}^{2}\right) \Im\left(\bar{A} B\left(e^{t_{1}}+e^{-t_{1}} \frac{A^{2}}{B^{2}}\right)\right)\right| \\
& =\left|\frac{e^{t_{1}}}{2} \Im\left(B^{2}\right)+\frac{e^{-t_{1}}}{4|B|^{2}}\left(-2 \operatorname{Re}(A B)-|A|^{2} \Im\left(B^{2}\right)\right)\right| .
\end{aligned}
$$

Now we use the definition of $A$ and $B$ to compute each term explicitly as follows

$$
\begin{aligned}
2 \operatorname{Re}(A B) & =-\left(e^{t_{2}}+e^{-t_{1}}\right) \sin \varphi \\
e^{t_{1}} \Im\left(B^{2}\right) & =\sin \varphi \\
e^{-t_{1}}|A|^{2} & =e^{t_{2}} \sin ^{2} \frac{\varphi}{2}+e^{-t_{2}} \cos ^{2} \frac{\varphi}{2} \\
e^{t_{1}}|B|^{2} & =e^{-t_{2}} \sin ^{2} \frac{\varphi}{2}+e^{t_{2}} \cos ^{2} \frac{\varphi}{2},
\end{aligned}
$$

and so

$$
\frac{1}{y^{2}}\left|\frac{\partial(x, y, \theta)}{\partial\left(t_{1}, t_{2}, \varphi\right)}\right|=\frac{1}{2}|\sin \varphi|
$$

Therefore, we conclude that

$$
\begin{equation*}
d V=\frac{1}{2}|\sin \varphi| d t_{1} d t_{2} d \varphi . \tag{A-1}
\end{equation*}
$$

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[^1]:    ${ }^{\dagger}$ Formally, it is the eigenvalue of the Laplace-Beltrami operator on $\mathbb{X}$ that corresponds to $\phi_{j}^{0}$.

[^2]:    ${ }^{\dagger}$ When $b=1$, this is a classical result due to Hecke [Siegel 1965, page 88].

