# On constant terms of Eisenstein series 

by

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1. Introduction. Let $F$ be a totally real field of degree $n$, and let $M_{k}(\mathfrak{n})$ denote the space of Hilbert modular forms of level $\mathfrak{n} \subset \mathcal{O}_{F}$ and weight $k$ over $F$. Let $E_{k}(\mathfrak{n}) \subset M_{k}(\mathfrak{n})$ denote the subspace of Eisenstein series. In this paper, we generalize results of [2, Section 2] and [7] to give the constant terms of nearly all Eisenstein series $E \in E_{k}(\mathfrak{n})$ at all cusps. The space $E_{k}(\mathfrak{n})$ has a basis consisting of forms of the form $\left.E_{k}(\eta, \psi)\right|_{\mathfrak{m}}$, where $\eta$ and $\psi$ are primitive ray class characters (see $\$ 4$ ). Our formula in Theorem 4.8 gives the constant terms of these series at all cusps when $\mathfrak{m}$ is square-free and coprime to the conductors of $\eta$ and $\psi$. In fact, Theorems 4.5 and 4.8 are more general than this; in particular we handle the case where $\eta$ and $\psi$ are not necessarily primitive characters. We work with all weights $k \geq 1$. In [7], only primitive characters are considered, the level raising operator $\left.\right|_{\mathfrak{m}}$ is not applied, and the weight $k$ is taken to be at least 2 .

There are concrete arithmetic applications of our results. In [3] we prove the Brumer-Stark conjecture away from $p=2$, and in [4] we prove an exact $p$ adic formula for Brumer-Stark units. Broadly speaking, both of these results apply Ribet's method, whereby cusp forms are constructed by taking linear combinations of products of Eisenstein series [8]. Central to the advance of [3] is the method by which this cusp form is constructed. For this, we require knowledge of the constant terms at all cusps of level-raised Eisenstein series associated to possibly imprimitive characters; we also need to include weight $k=1$. Therefore the calculations of [7] are not general enough for our application, which provides the motivation for this paper.

In addition, we prove here some other results that may be of independent interest. Firstly, we provide a complete enumeration of the cusps on the Hilbert modular variety. Also, we prove that in weight $k>1$, the cuspidality

[^0]of modular forms that are ordinary at a prime $p$ is regulated by the constant terms at cusps that are unramified at $p$. We provide two proofs of this fact; one applies our results on Eisenstein series, and the other is a direct study of the $U_{p}$ operator. While these two results are likely known to the experts, we have not found a precise reference for them in the literature.

We now outline the paper and describe our results in greater detail. In $\$ 2$ we recall the definition of the space $M_{k}(\mathfrak{n})$ of Hilbert modular forms of weight $k$ and level $\mathfrak{n} \subset \mathcal{O}_{F}$, following Shimura [10]. Associated to each $\lambda$ in the narrow class group $\mathrm{Cl}^{+}(F)$ is a congruence subgroup $\Gamma_{1, \lambda}(\mathfrak{n}) \subset \mathbf{G L}_{2}^{+}(F)$. The open Hilbert modular variety corresponding to our forms has $h^{+}=$ $\# \mathrm{Cl}^{+}(F)$ components:

$$
Y=\bigsqcup_{\lambda \in \mathrm{Cl}^{+}(F)} \Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathcal{H}^{n}, \quad \mathcal{H}=\text { complex upper half-plane. }
$$

The space $M_{k}(\mathfrak{n})$ of modular forms is endowed with an action of Hecke operators described in $\$ 2.5$. Among these operators are the diamond operators $S(\mathfrak{e})$, indexed by the classes $\mathfrak{e} \in G_{\mathfrak{n}}^{+}$, the narrow ray class group of $F$ attached to the modulus $\mathfrak{n}$. The diamond operators play a central role in our applications [3], 4].

In $\S 3$, we study the set of cusps associated to $Y$ :

$$
\operatorname{cusps}(\mathfrak{n})=\bigsqcup_{\lambda \in \mathrm{Cl}^{+}(F)} \Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbb{P}^{1}(F)
$$

We provide an explicit enumeration of this set. For $\mathfrak{m} \mid \mathfrak{n}$, let $Q_{\mathfrak{m}, \mathfrak{n}}$ denote the quotient of $G_{\mathfrak{m}}^{+} \times G_{\mathfrak{n} / \mathfrak{m}}^{+}$by the subgroup generated by diagonally embedded principal ideals $(x)$, where $x \in \mathcal{O}_{F}$ is congruent to 1 modulo $\mathfrak{n}$. The following result proved in 3.3 is already implicit in [12, pp. 422-423].

Theorem 1. There is a stratification $\operatorname{cusps}(\mathfrak{n})=\bigsqcup_{\mathfrak{m} \mid \mathfrak{n}} \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}$ with $\# \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}$ $=\# Q_{\mathfrak{m}, \mathfrak{n}}$. Each $\mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}$ is stable under the action of $G_{\mathfrak{n}}^{+}$via the diamond operators.

In $\$ 4$ we study the Eisenstein series in $M_{k}(\mathfrak{n})$ and calculate their constant terms at all cusps. This generalizes the results of [2, Proposition 2.1] and [7]. We work in a more general setting in this paper by considering all cusps and allowing for Eisenstein series associated to imprimitive characters. We normalize our constant terms (see (4) below) so that they are independent of the choice of representatives (up to sign). Furthermore, with these normalizations the constant terms exhibit nice integrality properties that are studied by Silliman 11.

For an ideal $\mathfrak{b} \mid \mathfrak{n}$, define

$$
\begin{equation*}
C_{\infty}(\mathfrak{b}, \mathfrak{n})=\bigsqcup_{\mathfrak{b} \mid \mathfrak{m}} \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}} \tag{1}
\end{equation*}
$$

In Theorem 4.5 we prove the following.
TheOrem 2. Let $k>1$, and let $\chi_{1}$ and $\chi_{2}$ be narrow ray class characters of $F$ with associated signs $q_{1}, q_{2} \in(\mathbb{Z} / 2 \mathbb{Z})^{n}$, respectively. Assume that $\chi_{2}$ is primitive of conductor $\mathfrak{b}$. Then the constant term of $E_{k}\left(\chi_{1}, \chi_{2}\right)$ vanishes at any cusp not lying in $C_{\infty}(\mathfrak{b}, \mathfrak{n})$. Furthermore, if

$$
\mathcal{A} \in \Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbb{P}^{1}(F)
$$

is represented by $\alpha / \gamma \in \mathbb{P}^{1}(F)$ and lies in $C_{\infty}(\mathfrak{b}, \mathfrak{n})$, then the constant term of $E_{k}\left(\chi_{1}, \chi_{2}\right)$ at $\mathcal{A}$ normalized as in (4) is given by

$$
\begin{align*}
\frac{1}{2^{n}} \frac{\tau\left(\chi_{1} \chi_{2}^{-1}\right)}{\tau\left(\chi_{2}^{-1}\right)}\left(\frac{\mathrm{Nb}}{\mathrm{Nf}}\right)^{k} \operatorname{sgn}(-\gamma)^{q_{1}} & \operatorname{sgn}(\alpha)^{q_{2}} \chi_{1}\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}\right) \chi_{2}^{-1}\left(\mathfrak{a}_{\mathcal{A}}\right)  \tag{2}\\
& \times L\left(\chi^{-1}, 1-k\right) \prod_{\mathfrak{q}}\left(1-\chi(\mathfrak{q}) \mathrm{Nq}^{-k}\right)
\end{align*}
$$

Here $\chi$ denotes the primitive character associated to $\chi_{1} \chi_{2}^{-1}, \mathfrak{f}=\operatorname{cond}(\chi)$, and $\mathfrak{q}$ runs through all primes dividing $\mathfrak{n}$ but not $\mathfrak{f}$. The integral ideals $\mathfrak{a}_{\mathcal{A}}$ and $\mathfrak{c}_{\mathcal{A}}$ associated to $\mathcal{A}$ are defined in 13 , and the condition $\mathcal{A} \in C_{\infty}(\mathfrak{b}, \mathfrak{n})$ implies that $\mathfrak{b} \mid \mathfrak{c}_{\mathcal{A}}$.

In Theorem 4.5 we also consider the case $k=1$. In Theorem 4.8 we build on the result above and consider a more general case: we calculate the constant terms of all level-raised Eisenstein series $\left.E_{k}\left(\chi_{1}, \chi_{2}\right)\right|_{\mathfrak{m}}$, where $\chi_{1}$ and $\chi_{2}$ are possibly imprimitive, under certain mild conditions. These results are essential in our arithmetic applications [3], 4]. In those works we construct cusp forms by taking the appropriate linear combinations of products of Eisenstein series considered here with certain other auxiliary forms constructed in [11].

In $\$ 5$ we conclude with the following result on the cuspidality of ordinary forms that is applied in our arithmetic applications [3, 4]. Fix a prime $p$ and let $\mathfrak{P}=\operatorname{gcd}\left(p^{\infty}, \mathfrak{n}\right)$ denote the $p$-part of $\mathfrak{n}$. The set $C_{\infty}(\mathfrak{P}, \mathfrak{n})$ defined in (1) may be viewed as the set of " $p$-unramified" cusps.

Theorem 3. Let $p$ be a prime. If $f \in M_{k}(\mathfrak{n})$ is $\mathfrak{p}$-ordinary for each prime $\mathfrak{p} \subset \mathcal{O}_{F}$ dividing $p$, then $f$ is cuspidal if and only if the constant term of $f$ vanishes at each cusp in $C_{\infty}(\mathfrak{P}, \mathfrak{n})$.
2. Notation on Hilbert modular forms. We refer the reader to [2, §2.1] for precise definitions and notations, following Shimura [10], concerning
the space of classical Hilbert modular forms over the totally real field $F$. We recall certain aspects of those definitions.
2.1. Hilbert modular forms. Let $\mathcal{H}$ denote the complex upper halfplane endowed with the usual action of $\mathbf{G L}_{2}^{+}(\mathbb{R})$ via linear fractional transformations, where $\mathbf{G} \mathbf{L}_{2}^{+}$denotes the group of matrices with positive determinant. We fix an ordering of the $n$ embeddings $F \hookrightarrow \mathbb{R}$, which yields an embedding $\mathbf{G L}_{2}^{+}(F) \hookrightarrow \mathbf{G} \mathbf{L}_{2}^{+}(\mathbb{R})^{n}$ and hence an action of $\mathbf{G L} \mathbf{L}_{2}^{+}(F)$ on $\mathcal{H}^{n}$. Here $\mathbf{G} \mathbf{L}_{2}^{+}(F)$ denotes the group of matrices with totally positive determinant.

For each class $\lambda$ in the narrow class group $\mathrm{Cl}^{+}(F)$, we choose a representative fractional ideal $\mathfrak{t}_{\lambda}$. Let $\mathfrak{n} \subset \mathcal{O}_{F}$ be an ideal, and assume that the representative ideals $\mathfrak{t}_{\lambda}$ have been chosen to be relatively prime to $\mathfrak{n}$. Define the groups

$$
\begin{aligned}
& \Gamma_{0, \lambda}(\mathfrak{n})=\left\{\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right) \in \mathbf{G L}_{2}^{+}(F): \alpha, \delta \in \mathcal{O}_{F}, \gamma \in \mathfrak{t}_{\lambda} \mathfrak{d} \mathfrak{n}\right. \\
& \left.\beta \in\left(\mathfrak{t}_{\lambda} \mathfrak{d}\right)^{-1}, \alpha \delta-\beta \gamma \in \mathcal{O}_{F}^{*}\right\} \\
& \Gamma_{1, \lambda}(\mathfrak{n})=\left\{\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right) \in \Gamma_{0, \lambda}(\mathfrak{n}): \delta \equiv 1(\bmod \mathfrak{n})\right\}
\end{aligned}
$$

Here $\mathfrak{d}$ denotes the different of $F$.
Let $k$ be a positive integer. We denote by $M_{k}(\mathfrak{n})$ the space of Hilbert modular forms for $F$ of level $\mathfrak{n}$ and weight $k$. Each element $f \in M_{k}(\mathfrak{n})$ is a tuple $f=\left(f_{\lambda}\right)_{\lambda \in \mathrm{Cl}^{+}(F)}$ of holomorphic functions $f_{\lambda}: \mathcal{H}^{n} \rightarrow \mathbb{C}$ such that $\left.f_{\lambda}\right|_{A, k}=f_{\lambda}$ for all $\lambda \in \mathrm{Cl}^{+}(F)$ and $A \in \Gamma_{1, \lambda}(\mathfrak{n})$. Here the weight $k$ slash action is defined in the usual way:

$$
\begin{aligned}
& \left.f_{\lambda}\right|_{A, k}\left(z_{1}, \ldots, z_{n}\right) \\
& \quad=\mathrm{N}(\operatorname{det} A)^{k / 2} \prod_{i=1}^{n}\left(\gamma_{i} z_{i}+\delta_{i}\right)^{-k} \cdot f_{\lambda}\left(\frac{\alpha_{1} z_{1}+\beta_{1}}{\gamma_{1} z_{1}+\delta_{1}}, \ldots, \frac{\alpha_{n} z_{n}+\beta_{n}}{\gamma_{n} z_{n}+\delta_{n}}\right)
\end{aligned}
$$

where $\alpha_{i}$ denotes the image of $\alpha$ under the $i$ th real embedding of $F$ and similarly for $\beta_{i}, \gamma_{i}, \delta_{i}$.
2.2. Constant terms and cusp forms. Suppose that $\mathcal{A}=(A, \lambda)$ is an ordered pair with

$$
A=\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right) \in \mathbf{G L}_{2}^{+}(F)
$$

and $\lambda \in \mathrm{Cl}^{+}(F)$. We define the fractional ideal

$$
\mathfrak{b}_{\mathcal{A}}=\alpha \mathcal{O}_{F}+\gamma\left(\mathfrak{t}_{\lambda} \mathfrak{d}\right)^{-1}
$$

Given $f=\left(f_{\lambda}\right) \in M_{k}(\mathfrak{n})$ and a pair $\mathcal{A}=(A, \lambda)$ as above, the function $\left.f_{\lambda}\right|_{A, k}$ has a Fourier expansion

$$
\begin{equation*}
\left.f_{\lambda}\right|_{A, k}(z)=a_{\mathcal{A}}(0)+\sum_{\substack{b \in \mathfrak{a} \\ b \gg 0}} a_{\mathcal{A}}(b) e_{F}(b z), \tag{3}
\end{equation*}
$$

where $\mathfrak{a}$ is a lattice in $F$ depending on $\mathcal{A}$, and

$$
e_{F}(b z):=\exp \left(2 \pi i\left(b_{1} z_{1}+\cdots+b_{n} z_{n}\right)\right)
$$

Definition 2.1. The normalized constant term of the form $f$ at $\mathcal{A}$ is

$$
\begin{equation*}
c_{\mathcal{A}}(0, f)=a_{\mathcal{A}}(0) \cdot\left(\mathrm{Nt}_{\lambda}\right)^{-k / 2}\left(\mathrm{Nb}_{\mathcal{A}}\right)^{-k}(\mathrm{~N}(\operatorname{det} A))^{k / 2} \tag{4}
\end{equation*}
$$

As we will see later, the constant terms with this normalization will exhibit nice invariance properties as well as integrality properties. The space $S_{k}(\mathfrak{n}) \subset M_{k}(\mathfrak{n})$ of cusp forms is defined to be subspace of forms $f$ such that $c_{\mathcal{A}}(0, f)=0$ for all pairs $\mathcal{A}$.
2.3. $q$-Expansion. When $A=1$, we drop the subscript $A$ and write simply

$$
c_{\lambda}(0, f)=a_{\lambda}(0)\left(\mathrm{Nt}_{\lambda}\right)^{-k / 2}
$$

Furthermore when $A=1$, the lattice $\mathfrak{a}$ appearing in (3) is the ideal $\mathfrak{t}_{\lambda}$. Any non-zero integral ideal $\mathfrak{m}$ may be written $\mathfrak{m}=(b) \mathfrak{t}_{\lambda}^{-\mathrm{T}}$ with $b \in \mathfrak{t}_{\lambda}$ totally positive for a unique $\lambda \in \mathrm{Cl}^{+}(F)$. We define the normalized Fourier coefficients

$$
c(\mathfrak{m}, f)=a_{\lambda}(b)\left(\mathrm{Nt}_{\lambda}\right)^{-k / 2}
$$

The collection of normalized Fourier coefficients $\left\{c_{\lambda}(0, f), c(\mathfrak{m}, f)\right\}$ is called the $q$-expansion of $f$. Note that these normalized coefficients are denoted with a capital $C$ in 10 .
2.4. Forms over a field $K$. Each tuple $f \in M_{k}(\mathfrak{n})$ is determined by its $q$-expansion, which is the collection of coefficients

$$
c_{\lambda}(0, f) \in \mathbb{C}, \quad \lambda \in \mathrm{Cl}^{+}(F), \quad c(\mathfrak{m}, f) \in \mathbb{C}, \quad \mathfrak{m} \subset \mathcal{O}_{F}, \mathfrak{m} \neq 0
$$

defined in 82.3 . For any subfield $K \subset \mathbf{C}$, define $M_{k}(\mathfrak{n}, K)$ to be the $K$ vector subspace of $M_{k}(\mathfrak{n})$ consisting of modular forms whose $q$-expansion coefficients all lie in $K$. A fundamental result of Shimura [9, Theorem 7] states

$$
\begin{equation*}
M_{k}(\mathfrak{n}, K)=M_{k}(\mathfrak{n}, \mathbb{Q}) \otimes_{\mathbb{Q}} K \tag{5}
\end{equation*}
$$

We define $M_{k}(\mathfrak{n}, K)$ by more generally if $K$ is any field of characteristic 0 . This generalizes in the obvious way to define $S_{k}(\mathfrak{n}, K)$.
2.5. Hecke operators. The space $M_{k}(\mathfrak{n})$ is endowed with the action of a Hecke algebra $\tilde{\mathbf{T}} \subset \operatorname{End}\left(M_{k}(\mathfrak{n})\right)$ generated over $\mathbb{Z}$ by the following operators:

- $T_{\mathfrak{q}}$ for $\mathfrak{q} \nmid \mathfrak{n}$,
- $U_{\mathfrak{q}}$ for $\mathfrak{q} \mid \mathfrak{n}$,
- the "diamond operators" $S(\mathfrak{e})$ for each class $\mathfrak{e} \in G_{\mathfrak{n}}^{+}=$narrow ray class group of $F$ of conductor $\mathfrak{n}$.

We refer to [10, §2] for the definition of these operators. We warn that in loc. cit. both $T_{\mathfrak{q}}$ and $U_{\mathfrak{q}}$ are denoted by $T_{\mathfrak{q}}$.

Let us recall the definition of the diamond operators $S(\mathfrak{e})$. Let $f=$ $\left(f_{\lambda}\right)_{\lambda \in \mathrm{Cl}^{+}(F)} \in M_{k}(\mathfrak{n})$, and let $\mathfrak{e}$ denote an ideal of $\mathcal{O}_{F}$ that is relatively prime to $\mathfrak{n}$, and which hence represents a class in $G_{\mathfrak{n}}^{+}$. For each $\mu \in \mathrm{Cl}^{+}(F)$, let $\lambda \in \mathrm{Cl}^{+}(F)$ denote the class of $\mu \mathfrak{e}^{-2}$. Write $\mathfrak{t}_{\lambda} \mathfrak{t}_{\mu}^{-1} \mathfrak{e}^{2}=(x)$ where $x$ is a totally positive element of $F^{*}$, uniquely determined up to multiplication by a totally positive unit in $\mathcal{O}_{F}^{*}$. Let

$$
A_{\mu}=\left(\begin{array}{cc}
\alpha_{\mu} & \beta_{\mu}  \tag{6}\\
\gamma_{\mu} & \delta_{\mu}
\end{array}\right) \in \mathbf{G} \mathbf{L}_{2}^{+}(F)
$$

be a matrix satisfying the following conditions:

$$
\begin{align*}
& \alpha_{\mu} \in \mathfrak{e}, \quad \beta_{\mu} \in \mathfrak{t}_{\mu}^{-1} \mathfrak{d}^{-1} \mathfrak{e}, \quad \gamma_{\mu} \in \mathfrak{t}_{\lambda} \mathfrak{d n e}, \quad \delta_{\mu} \in \mathfrak{t}_{\lambda} \mathfrak{t}_{\mu}^{-1} \mathfrak{e},  \tag{7}\\
& \operatorname{det} A_{\mu}=x, \quad \delta_{\mu} \equiv x\left(\bmod \mathfrak{t}_{\lambda} \mathfrak{t}_{\mu}^{-1} \mathfrak{e n}\right)
\end{align*}
$$

Then

$$
\begin{equation*}
\left.f\right|_{S(\mathfrak{e})}=\left(g_{\mu}\right)_{\mu \in \mathrm{Cl}^{+}(F)} \quad \text { where } g_{\mu}=\left.f_{\lambda}\right|_{A_{\mu}} \tag{8}
\end{equation*}
$$

2.6. Raising the level. For a Hilbert modular form $f \in M_{k}(\mathfrak{n})$ and an integral ideal $\mathfrak{q}$ of $F$, there is a form

$$
\left.f\right|_{\mathfrak{q}} \in M_{k}(\mathfrak{n q})
$$

characterized by the fact that for non-zero integral ideals $\mathfrak{a}$ we have

$$
c\left(\mathfrak{a},\left.f\right|_{\mathfrak{q}}\right)= \begin{cases}c(\mathfrak{a} / \mathfrak{q}, f) & \text { if } \mathfrak{q} \mid \mathfrak{a}  \tag{9}\\ 0 & \text { if } \mathfrak{q} \nmid \mathfrak{a}\end{cases}
$$

and

$$
\begin{equation*}
c_{\lambda}\left(0,\left.f\right|_{\mathfrak{q}}\right)=c_{\lambda \mathfrak{q}}(0, f) \tag{10}
\end{equation*}
$$

for all $\lambda \in \mathrm{Cl}^{+}(F)$. We recall the definition of $\left.f\right|_{\mathfrak{q}}$. For every $\lambda$ there is a $\mu \in \mathrm{Cl}^{+}(F)$ and a totally positive element $\alpha_{\mu} \in F$ such that

$$
\mathfrak{q}_{\lambda}=\left(\alpha_{\mu}\right) \mathfrak{t}_{\mu} .
$$

Then

$$
\left.\left(\left.f\right|_{\mathfrak{q}}\right)_{\lambda}:=\mathrm{Nq}^{-k / 2} f_{\mu} \left\lvert\, \begin{array}{rr}
\alpha_{\mu} & 0  \tag{11}\\
0 & 1
\end{array}\right.\right)
$$

The fact that $\left.f\right|_{\mathfrak{q}}$ lies in $M_{k}(\mathfrak{n q})$ and satisfies $(9)-(10)$ is proven in [10, Prop. 2.3].

## 3. Cusps

3.1. Admissibility. Recall the fractional ideal

$$
\begin{equation*}
\mathfrak{b}_{\mathcal{A}}=\alpha \mathcal{O}_{F}+\gamma\left(\mathfrak{t}_{\lambda} \mathfrak{d}\right)^{-1} \tag{12}
\end{equation*}
$$

defined in 2.2 associated to a pair $\mathcal{A}=(A, \lambda)$ with $A \in \mathbf{G L}_{2}^{+}(F)$ and $\lambda \in \mathrm{Cl}^{+}(F)$. We now define the integral ideals

$$
\begin{equation*}
\mathfrak{a}_{\mathcal{A}}=\alpha \mathfrak{b}_{\mathcal{A}}^{-1}, \quad \mathfrak{c}_{\mathcal{A}}=\gamma\left(\mathfrak{t}_{\lambda} \mathfrak{d b}_{\mathcal{A}}\right)^{-1} \tag{13}
\end{equation*}
$$

These ideals are relatively prime.
To explain the meaning of these invariants, consider the case $F=\mathbb{Q}$. If $\mathcal{A}$ represents the cusp $\alpha / \gamma \in \mathbb{P}^{1}(\mathbb{Q})$ with $\alpha, \gamma$ relatively prime integers, then $\mathfrak{a}_{\mathcal{A}}=(\alpha)$ and $\mathfrak{c}_{\mathcal{A}}=(\gamma)$. Finally, we define

$$
\begin{equation*}
\mathfrak{m}_{\mathcal{A}}=\operatorname{gcd}\left(\mathfrak{c}_{\mathcal{A}}, \mathfrak{n}\right) \tag{14}
\end{equation*}
$$

Given a form $f \in M_{k}(\mathfrak{n})$, it is clear that the normalized constant term $c_{\mathcal{A}}(0, f)$ defined in 2.2 depends only on $A$ up to left multiplication by an element of $\Gamma_{1, \lambda}(\mathfrak{n})$. Furthermore, if we write

$$
B=\left\{\left(\begin{array}{cc}
\alpha & \beta \\
0 & \delta
\end{array}\right) \in \mathbf{G} \mathbf{L}_{2}^{+}(F)\right\}
$$

it is almost true that $c_{\mathcal{A}}(0, f)$ depends on $A$ up to right multiplication by an element of $B$-there is a sign ambiguity

$$
c_{\left(A A^{\prime}, \lambda\right)}(0, f)=\operatorname{sgn}(\mathrm{N} \alpha)^{k} c_{(A, \lambda)}(0, f), \quad A^{\prime}=\left(\begin{array}{ll}
\alpha & \beta  \tag{15}\\
0 & \delta
\end{array}\right) \in B
$$

To see this, one combines the following equalities, which follow from the definitions (3) and (12), together with the definition (4) of $c_{\mathcal{A}}(0, f)$ :

$$
\begin{aligned}
& a_{A A^{\prime}}(0)=a_{A}(0) \mathrm{N}(\alpha / \delta)^{k / 2} \\
& \mathrm{Nb}_{A A^{\prime}}=\mathrm{Nb} \\
& A \cdot|\mathrm{~N} \alpha| \\
& \operatorname{det}\left(A A^{\prime}\right)=\operatorname{det}(A) \mathrm{N}(\alpha \delta)
\end{aligned}
$$

If $k$ is odd and the class of $A$ in $\Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbf{G} \mathbf{L}_{2}^{+}(F)$ is fixed under right multiplication by an element $A^{\prime} \in B$ with $\mathrm{N} \alpha<0$, then it follows from 15 ) that $c_{\mathcal{A}}(0, f)=0$. Let us determine the pairs $\mathcal{A}$ for which this is the case.

Definition 3.1. A pair $(\mathfrak{m}, \mathfrak{n})$ with $\mathfrak{m} \mid \mathfrak{n}$ is called admissible if there does not exist a pair of units $\epsilon_{1}, \epsilon_{2} \in \mathcal{O}_{F}^{*}$ such that $\epsilon_{1} \equiv 1(\bmod \mathfrak{m})$ and $\epsilon_{2} \equiv 1(\bmod \mathfrak{n} / \mathfrak{m})$ with $N \epsilon_{1}=\mathrm{N} \epsilon_{2}=-1$ and $\epsilon_{1} / \epsilon_{2}$ totally positive.

Definition 3.2. With the level $\mathfrak{n}$ fixed, a pair $\mathcal{A}=(A, \lambda)$ with $A \in$ $\mathbf{G} \mathbf{L}_{2}^{+}(F)$ and $\lambda \in \mathrm{Cl}^{+}(F)$ is called admissible if $\left(\mathfrak{m}_{A}, \mathfrak{n}\right)$ is admissible.

Theorem 3.3. Given a pair $\mathcal{A}=(A, \lambda)$, the class of $A$ in $\Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbf{G} \mathbf{L}_{2}^{+}(F)$ is fixed by right multiplication by an element $A^{\prime} \in B$ with $\mathrm{N} \alpha<0$ if and only if $\mathcal{A}$ is not admissible.

Before proving the theorem, we introduce some notation and prove an important lemma. Given a fractional ideal $\mathfrak{b}$ and an integral ideal $\mathfrak{m}$, we denote by $(\mathfrak{b} / \mathfrak{b m})^{*}$ the subset of elements of $\mathfrak{b} / \mathfrak{b m}$ that generate this quotient as an $\mathcal{O}_{F} / \mathfrak{m}$-module. This is a principal homogeneous space for the group $\left(\mathcal{O}_{F} / \mathfrak{m}\right)^{*}$.

Definition 3.4. For a fractional ideal $\mathfrak{b}$ and an integral ideal $\mathfrak{m}$, we define

$$
\mathcal{R}_{\mathfrak{m}}^{\mathfrak{b}}=(\mathfrak{b} / \mathfrak{b m})^{*} / \mathcal{O}_{F,+}^{*}
$$

the quotient of the set $(\mathfrak{b} / \mathfrak{b m})^{*}$ by the action of multiplication by the group of totally positive units of $F$.

Definition 3.5. Let $\mathcal{P}_{\lambda}(\mathfrak{n})$ be the set of tuples $(\mathfrak{b}, \mathfrak{m}, \alpha, \gamma)$ where $\mathfrak{b}$ is a fractional ideal of $F, \mathfrak{m}$ is an integral ideal dividing $\mathfrak{n}, \alpha \in \mathcal{R}_{\mathfrak{m}}^{\mathfrak{b}}$, and $\gamma \in \mathcal{R}_{\mathfrak{n} / \mathfrak{m}}^{\mathfrak{b j} \mathfrak{t}_{\lambda} \mathfrak{m}}$.

Note that the number of elements in $\mathcal{R}_{\mathfrak{n} / \mathfrak{m}}^{\mathfrak{b d} \mathfrak{t}_{\lambda} \mathfrak{m}}=\left(\mathfrak{b d} \mathfrak{t}_{\lambda} \mathfrak{m} / \mathfrak{b d} \mathfrak{t}_{\lambda} \mathfrak{n}\right)^{*} / \mathcal{O}_{F,+}^{*}$ is independent of $\lambda$, thanks to the bijection with $\left(\mathcal{O}_{F} /(\mathfrak{n} / \mathfrak{m})\right)^{*} / \mathcal{O}_{F,+}^{*}$. Therefore the cardinality of $\mathcal{P}_{\lambda}(\mathfrak{n})$ is also independent of $\lambda$.

The heart of Theorem 3.3 is the following lemma.
Lemma 3.6. Fix $\lambda \in \mathrm{Cl}^{+}(F)$. There is a canonical bijection

$$
\varphi: \Gamma_{1, \lambda}(\mathfrak{n}) \backslash\left(F^{2} \backslash(0,0)\right) \rightarrow \mathcal{P}_{\lambda}(\mathfrak{n})
$$

given by $(\alpha, \gamma) \mapsto(\mathfrak{b}, \mathfrak{m}, \bar{\alpha}, \bar{\gamma})$, where $\mathfrak{b}=\mathfrak{b}_{\mathcal{A}}$ and $\mathfrak{m}=\mathfrak{m}_{\mathcal{A}}$ are defined as in (12) and (14).

Proof. The fact that the map $\varphi$ is well-defined is elementary and left to the reader. Surjectivity is also not difficult: Given a fractional ideal $\mathfrak{b}$ and an integral ideal $\mathfrak{m} \mid \mathfrak{n}$, choose $\gamma \in \mathfrak{t}_{\lambda} \mathfrak{d} \mathfrak{b}$ such that $\operatorname{gcd}\left(\gamma \mathfrak{d}^{-1} \mathfrak{t}_{\lambda}^{-1}, \mathfrak{n}\right)=\mathfrak{m}$. Scaling $\gamma$ by an appropriate element of $\mathcal{O}_{F}$ relatively prime to $\mathfrak{n}$, we can ensure that $\gamma$ lands in any class in $\mathcal{R}_{\mathfrak{n} / \mathfrak{m}}^{\mathfrak{b d} \mathfrak{t}^{\mathfrak{m}}}$ without changing the gcd condition. Next we choose any $\alpha \in \mathcal{O}_{F}$ such that

$$
\begin{equation*}
\alpha \mathcal{O}_{F}+\gamma\left(\mathfrak{d} \mathfrak{t}_{\lambda}\right)^{-1}=\mathfrak{b} \tag{16}
\end{equation*}
$$

Scaling $\alpha$ by an appropriate element of $\mathcal{O}_{F}$ relatively prime to $\gamma\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}\right)^{-1}$, we can ensure that $\alpha$ lands in any class in $\mathcal{R}_{\mathfrak{m}}^{\mathfrak{b}}$ without affecting 16. This proves the desired surjectivity.

For injectivity, suppose that

$$
\begin{equation*}
\varphi(\alpha, \gamma)=\varphi\left(\alpha^{\prime}, \gamma^{\prime}\right)=(\mathfrak{b}, \mathfrak{m}, \bar{\alpha}, \bar{\gamma}) . \tag{17}
\end{equation*}
$$

Let $\mathcal{A}$ correspond to $(\alpha, \gamma)$ and $\mathcal{A}^{\prime}$ to $\left(\alpha^{\prime}, \gamma^{\prime}\right)$ as in the statement of the lemma. From the third component of (17), there exists a totally positive unit $\epsilon \in \mathcal{O}_{F,+}^{*}$ such that $\alpha^{\prime} \equiv \epsilon \alpha(\bmod \sqrt{\mathfrak{b m}})$. Since $\gamma \mathfrak{d}^{-1} \mathfrak{t}_{\lambda}^{-1}+\mathfrak{b n}=\mathfrak{b m}$, we can act by an element of the form $\left(\begin{array}{cc}\epsilon & \beta \\ 0 & 1\end{array}\right)$ on $(\alpha, \gamma)$ with $\beta \in\left(\mathfrak{d t}_{\lambda}\right)^{-1}$ to ensure that $\alpha \equiv \alpha^{\prime}(\bmod \mathfrak{b n})$.

Next let $u \in \mathcal{O}_{F,+}^{*}$ have $\gamma^{\prime} \equiv u \gamma\left(\bmod \mathfrak{b o t}{ }_{\lambda} \mathfrak{n}\right)$. Such a $u$ exists by the fourth component of 177$)$. Note that the pairs $(\alpha, \gamma),\left(\alpha^{\prime}, \gamma^{\prime}\right)$ can be completed to matrices

$$
M=\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right), \quad \operatorname{det} M=1, \quad M^{\prime}=\left(\begin{array}{ll}
\alpha^{\prime} & \beta^{\prime} \\
\gamma^{\prime} & \delta^{\prime}
\end{array}\right), \quad \operatorname{det} M^{\prime}=u
$$

with $\beta, \beta^{\prime} \in\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}\right)^{-1}$ and $\delta, \delta^{\prime} \in \mathfrak{b}^{-1}$. We claim that $\delta$ and $\delta^{\prime}$ can be chosen to satisfy

$$
\delta^{\prime} \equiv \delta u\left(\bmod \mathfrak{b}^{-1} \mathfrak{n}\right) .
$$

Granting the claim, it is straightforward to check that

$$
M^{\prime} M^{-1}\binom{\alpha}{\gamma}=\binom{\alpha^{\prime}}{\gamma^{\prime}} \quad \text { and } \quad M^{\prime} M^{-1} \in \Gamma_{1, \lambda}(\mathfrak{n})
$$

as desired.
It remains to prove the claim. Given any $y \in \mathfrak{b}^{-1} \mathfrak{c}_{\mathcal{A}}$, we can write $y=\gamma x$ with $x \in\left(\mathfrak{d t}_{\lambda} \mathfrak{b}^{2}\right)^{-1}$ and replace $(\beta, \delta)$ by $(\beta+\alpha x, \delta+\gamma x)$. Hence $\delta$ can be replaced by any element in its equivalence class in $\mathfrak{b}^{-1} / \mathfrak{b}^{-1} \mathfrak{c}_{\mathcal{A}}$. Since $\operatorname{gcd}\left(\mathfrak{c}_{\mathcal{A}}, \mathfrak{n}\right)=\mathfrak{m}$, to prove the claim it therefore suffices to show that $\delta^{\prime} \equiv$ $u \delta\left(\bmod \mathfrak{b}^{-1} \mathfrak{m}\right)$. If $\alpha=0$, this is clear since $\mathfrak{m}=1$. Otherwise, multiplication by $\alpha$ induces an isomorphism $\mathfrak{b}^{-1} / \mathfrak{b}^{-1} \mathfrak{m} \rightarrow \mathfrak{a}_{\mathcal{A}} / \mathfrak{a}_{\mathcal{A}} \mathfrak{m}$, so we must show that $\alpha \delta^{\prime} \equiv \alpha \delta u\left(\bmod \mathfrak{a}_{\mathcal{A}} \mathfrak{m}\right)$. Now $\mathfrak{m}$ divides $\mathfrak{c}_{\mathcal{A}}$, which is coprime to $\mathfrak{a}_{\mathcal{A}}$, so by the Chinese Remainder Theorem it suffices to separately show that $\alpha \delta^{\prime} \equiv$ $\alpha \delta u\left(\bmod \mathfrak{a}_{\mathcal{A}}\right)$ and $\alpha \delta^{\prime} \equiv \alpha \delta u(\bmod \mathfrak{m})$. The first of these is trivial since $\alpha \delta^{\prime}, \alpha \delta u \in \mathfrak{a}_{\mathcal{A}}$. To prove $\alpha \delta^{\prime} \equiv \alpha \delta u(\bmod \mathfrak{m})$ we note

$$
\alpha \delta^{\prime}-\alpha \delta u=\beta^{\prime} \gamma^{\prime}+\left(\alpha-\alpha^{\prime}\right) \delta^{\prime}-\beta \gamma u
$$

with $\gamma=\beta^{\prime} \gamma^{\prime} \in \mathfrak{c}_{\mathcal{A}^{\prime}} \subset \mathfrak{m}, \beta \gamma u \in \mathfrak{c}_{\mathcal{A}} \subset \mathfrak{m}$, and $\left(\alpha-\alpha^{\prime}\right) \delta \in \mathfrak{n} \subset \mathfrak{m}$. This concludes the proof of Lemma 3.6.

Proof of Theorem 3.3. Suppose that $\mathcal{A}=(A, \lambda)$ is not admissible with $A=\binom{\alpha}{\gamma}$, and let $\mathfrak{m}=\mathfrak{m}_{\mathcal{A}}$. Let $\epsilon=\epsilon_{1} \in \mathcal{O}_{F}^{*}$ in the definition of admissibility, so $N \epsilon=-1, \epsilon \equiv 1(\bmod \mathfrak{m})$, and $\epsilon \equiv \epsilon_{1} \epsilon_{2}(\bmod \mathfrak{n} / \mathfrak{m})$ with $\epsilon_{1} \epsilon_{2} \in \mathcal{O}_{F}^{*}$ a totally positive unit. Right multiplication by $M=\left(\begin{array}{cc}\epsilon & 0 \\ 0 & \epsilon^{-1}\end{array}\right)$ sends $A$ to $\binom{\alpha \epsilon *}{\gamma \epsilon *}$. It is immediate from the definition that $\varphi(\alpha, \gamma)=\varphi(\alpha \epsilon, \gamma \epsilon)$, hence Lemma 3.6 implies that there exists $A^{\prime} \in \Gamma_{1, \lambda}(\mathfrak{n})$ such that $A^{\prime} A$ and $A M$
have the same first column. Therefore there exists $N=\left(\begin{array}{ll}1 & * \\ 0 & *\end{array}\right) \in B$ such that $A^{\prime} A=A M N$. It follows that the class of $A$ in $\Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbf{G L}_{2}^{+}(F)$ is fixed by right multiplication by $M N \in B$, and the lower right entry of $M N$ has negative norm.

To prove the converse, suppose that $A=\left(\begin{array}{cc}\alpha & \beta \\ \gamma & \delta\end{array}\right) \in \mathbf{G L}_{2}^{+}(F)$ and the class of $A$ in $\Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbf{G L}_{2}^{+}(F)$ is fixed by right multiplication by $M=\left(\begin{array}{cc}x & * \\ 0 & z\end{array}\right) \in B$, where $\mathrm{N} z<0$. By Lemma 3.6, we have $\varphi(\alpha, \gamma)=\varphi(\alpha x, \gamma x)$. From the first component of this equation, we see that $x \in \mathcal{O}_{F}^{*}$. From the third and fourth components we see that there exist totally positive units $u_{1}, u_{2} \in \mathcal{O}_{F,+}^{*}$ such that

$$
x \equiv u_{1}(\bmod \mathfrak{m}), \quad x \equiv u_{2}(\bmod \mathfrak{n} / \mathfrak{m})
$$

where $\mathfrak{m}=\mathfrak{m}_{\mathcal{A}}$. Note that $\mathrm{N} x<0$ (and hence $\mathrm{N} x=-1$ ) since $\mathrm{N} z<0$ and $x z$ is totally positive. Then letting $\epsilon_{1}=x / u_{1}$ and $\epsilon_{2}=x / u_{2}$ shows that $(\mathfrak{m}, \mathfrak{n})$ is not admissible.

Corollary 3.7. Let $k$ be odd, and let $f \in M_{k}(\mathfrak{n})$. The constant term $c_{\mathcal{A}}(0, f)$ vanishes if $\mathcal{A}$ is not admissible.

Proof. By Theorem 3.3, if $\mathcal{A}=(A, \lambda)$ is not admissible then the class of $A$ in $\Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbf{G L}_{2}^{+}(F)$ is fixed by right multiplication by an element $A^{\prime} \in B$ with $\mathrm{N} \alpha<0$. Then $c_{\mathcal{A}}(0, f)=c_{\left(A A^{\prime}, \lambda\right)}(0, f)$; but by 15 we also have $c_{\mathcal{A}}(0, f)=-c_{\left(A A^{\prime}, \lambda\right)}(0, f)$. The result follows.
3.2. Definition of cusps. For $\Gamma \subset \Gamma_{0, \lambda}(\mathfrak{n})$ any congruence subgroup, the associated set of cusps is by definition the finite set

$$
\operatorname{cusps}(\Gamma):=\Gamma \backslash \mathbf{G} \mathbf{L}_{2}^{+}(F) /\left\{\left(\begin{array}{cc}
\alpha & \beta  \tag{18}\\
0 & \delta
\end{array}\right) \in \mathbf{G L}_{2}^{+}(F)\right\} \leftrightarrow \Gamma \backslash \mathbb{P}^{1}(F)
$$

The bijection in 18 is $\left(\begin{array}{cc}\alpha & \beta \\ \gamma & \delta\end{array}\right) \mapsto(\alpha: \gamma)$. We define

$$
\begin{equation*}
\operatorname{cusps}(\mathfrak{n})=\bigsqcup_{\lambda} \operatorname{cusps}\left(\Gamma_{1, \lambda}(\mathfrak{n})\right) \tag{19}
\end{equation*}
$$

An ordered pair $\mathcal{A}=(A, \lambda)$ with $A \in \mathbf{G L}_{2}^{+}(F)$ and $\lambda \in \mathrm{Cl}^{+}(F)$ gives rise to an element of $\operatorname{cusps}(\mathfrak{n})$ by considering the image of $A$ in $\operatorname{cusps}\left(\Gamma_{1, \lambda}(\mathfrak{n})\right)$ in the $\lambda$-component of the disjoint union $\sqrt[19]{ }$. The cusp represented by $\mathcal{A}$ will be denoted $[\mathcal{A}] \in \operatorname{cusps}(\mathfrak{n})$.

Definition 3.8. For each $\mathfrak{m} \mid \mathfrak{n}$, we define $\mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}$ to be the set of cusps $[\mathcal{A}] \in \operatorname{cusps}(\mathfrak{n})$ such that $\mathfrak{m}_{\mathcal{A}}=\mathfrak{m}$.

The set of admissible cusps is defined by

$$
\operatorname{cusps}^{*}(\mathfrak{n})=\bigsqcup_{\substack{\mathfrak{m} \mid \mathfrak{n} \\(\mathfrak{m}, \mathfrak{n}) \text { admissible }}} \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}=\{[\mathcal{A}]: \mathcal{A} \text { is admissible }\}
$$

There is a canonical action of the diamond operators on cusps( $\mathfrak{n}$ ) that is compatible with its action on modular forms. Given an integral ideal $\mathfrak{e}$ coprime to $\mathfrak{n}$ and a cusp $\mathcal{A}=(A, \mu)$, we define $\lambda$ and $A_{\mu}$ as in 7 and define

$$
S(\mathfrak{e}) \mathcal{A}=\mathcal{A}^{\prime}=\left(A^{\prime}, \lambda\right), \quad \text { where } \quad A^{\prime}=A_{\mu} A
$$

Proposition 3.9. Each set $\mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}$ is invariant under the action of $G_{\mathfrak{n}}^{+}$via the diamond operators.

Proof. With notation as above, one checks directly from the definitions that $\mathfrak{b}_{\mathcal{A}^{\prime}}=\mathfrak{b}_{\mathcal{A}} \mathfrak{m}$. Furthermore one calculates that

$$
\begin{aligned}
\mathfrak{m}_{\mathcal{A}^{\prime}} \mathfrak{b}_{\mathcal{A}^{\prime}} & =\operatorname{gcd}\left(\left(\gamma_{\mu} \alpha+\delta_{\mu} \gamma\right)\left(\mathfrak{d}_{\lambda}\right)^{-1}, \mathfrak{n} \mathfrak{b}_{\mathcal{A}^{\prime}}\right)=\operatorname{gcd}\left(\delta_{\mu} \gamma\left(\mathfrak{d}_{\lambda}\right)^{-1}, \mathfrak{n} \mathfrak{b}_{\mathcal{A}^{\prime}}\right) \\
& \subset \operatorname{gcd}\left(\gamma\left(\mathfrak{d} \mathfrak{t}_{\mu}\right)^{-1} \mathfrak{m}, \mathfrak{n} \mathfrak{b}_{\mathcal{A}^{\prime}}\right)=\operatorname{gcd}\left(\mathfrak{c}_{\mathcal{A}} \mathfrak{b}_{\mathcal{A}} \mathfrak{m}, \mathfrak{n b}_{\mathcal{A}^{\prime}}\right)=\mathfrak{m}_{\mathcal{A}^{\prime}} \mathfrak{b}_{\mathcal{A}^{\prime}}
\end{aligned}
$$

Therefore $\mathfrak{m}_{\mathcal{A}^{\prime}} \subset \mathfrak{m}_{\mathcal{A}}$. However, this is a group action; replacing $\mathfrak{m}$ by an ideal whose image in $G_{\mathfrak{n}}^{+}$is inverse to $\mathfrak{m}$, and switching the roles of $\mathcal{A}$ and $\mathcal{A}^{\prime}$, we find the reverse inclusion $\mathfrak{m}_{\mathcal{A}} \subset \mathfrak{m}_{\mathcal{A}^{\prime}}$. Hence $\mathfrak{m}_{\mathcal{A}}=\mathfrak{m}_{\mathcal{A}^{\prime}}$, and the result follows.
3.3. Enumeration of cusps. We now enumerate $\operatorname{cusps}(\mathfrak{n})$ and more specifically each subset $\mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}$. Recall the set of 4-tuples $\mathcal{P}_{\lambda}(\mathfrak{n})$ defined in Definition 3.5. For each ideal $\mathfrak{m} \mid \mathfrak{n}$, let $\mathcal{P}_{\lambda}(\mathfrak{m}, \mathfrak{n}) \subset \mathcal{P}_{\lambda}(\mathfrak{n})$ denote the set of tuples whose second coordinate is $\mathfrak{m}$. There is a natural action of $F^{*}$ on $\mathcal{P}_{\lambda}(\mathfrak{n})$ that preserves each $\mathcal{P}_{\lambda}(\mathfrak{m}, \mathfrak{n})$, given by

$$
x \cdot(\mathfrak{b}, \mathfrak{m}, \alpha, \gamma)=(\mathfrak{b}(x), \mathfrak{m}, \alpha x, \gamma x) .
$$

The following is an immediate corollary of Lemma 3.6 .
Corollary 3.10. For each class $\lambda \in G_{1}^{+}$, there is a canonical bijection

$$
\Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbb{P}^{1}(F) \rightarrow P_{\lambda}(\mathfrak{n}) / F^{*}
$$

There are canonical bijections

$$
\begin{equation*}
\mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}} \rightarrow \bigsqcup_{\lambda \in \mathrm{Cl}^{+}(F)} P_{\lambda}(\mathfrak{m}, \mathfrak{n}) / F^{*}, \quad \operatorname{cusps}(\mathfrak{n}) \rightarrow \bigsqcup_{\lambda \in \mathrm{Cl}^{+}(F)} P_{\lambda}(\mathfrak{n}) / F^{*} \tag{20}
\end{equation*}
$$

Definition 3.11. Let $Q_{\mathfrak{m}, \mathfrak{n}}$ denote the quotient of $G_{\mathfrak{m}}^{+} \times G_{\mathfrak{n} / \mathfrak{m}}^{+}$by the subgroup generated by diagonally embedded principal ideals $(x)$, where $x \in \mathcal{O}_{F}$ is congruent to 1 modulo $\mathfrak{n}$.

Corollary 3.12. We have $\# \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}=\# Q_{\mathfrak{m}, \mathfrak{n}}$, hence

$$
\# \operatorname{cusps}(\mathfrak{n})=\sum_{\mathfrak{m} \mid \mathfrak{n}} \# Q_{\mathfrak{m}, \mathfrak{n}}
$$

Proof. From 20) we deduce that $\# \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}=h^{+} \#\left(P_{\lambda}(\mathfrak{m}, \mathfrak{n}) / F^{*}\right)$, where $h^{+}=\# \mathrm{Cl}^{+}(F)$. There is a surjective map

$$
\begin{equation*}
P_{\lambda}(\mathfrak{m}, \mathfrak{n}) / F^{*} \rightarrow \mathrm{Cl}(F), \quad(\mathfrak{b}, \alpha, \gamma) \mapsto[\mathfrak{b}] . \tag{21}
\end{equation*}
$$

For an ideal $\mathfrak{a}$, we put $\mathcal{R}_{\mathfrak{a}}=\mathcal{R}_{\mathfrak{a}}^{\mathcal{O}_{F}}$. If $\mathcal{U}$ denotes the image of $\mathcal{O}_{F}^{*}$ mapped diagonally to $\mathcal{R}_{\mathfrak{m}} \times \mathcal{R}_{\mathfrak{n} / \mathfrak{m}}$, then the fiber over a point in (21) is a principal homogeneous space for $\left(\mathcal{R}_{\mathfrak{m}} \times \mathcal{R}_{\mathfrak{n} / \mathfrak{m}}\right) / \mathcal{U}$. Therefore

$$
\begin{equation*}
\# \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}=\left(h^{+} h\right) \cdot \#\left(\left(\mathcal{R}_{\mathfrak{m}} \times \mathcal{R}_{\mathfrak{n} / \mathfrak{m}}\right) / \mathcal{U}\right), \tag{22}
\end{equation*}
$$

where $h=\# \mathrm{Cl}(F)$. Meanwhile, there is an exact sequence

$$
\begin{equation*}
1 \rightarrow\left(\mathcal{R}_{\mathfrak{m}} \times \mathcal{R}_{\mathfrak{n} / \mathfrak{m}}\right) / \mathcal{U} \rightarrow Q_{\mathfrak{m}, \mathfrak{n}} \rightarrow \mathrm{Cl}^{+}(F) \times \mathrm{Cl}(F) \rightarrow 1 \tag{23}
\end{equation*}
$$

where the second non-trivial arrow is $([\mathfrak{a}],[\mathfrak{b}]) \mapsto([\mathfrak{a} / \mathfrak{b}],[\mathfrak{a}])$. From (23) we deduce that $\# Q_{\mathfrak{m}, \mathfrak{n}}$ also equals the right side of 22$)$.
3.4. Constant term map. If $k$ is even we define

$$
C_{k}=\bigoplus_{\operatorname{cusps}(\mathfrak{n})} \mathbb{C} .
$$

Then (15) implies that we have a well-defined constant term map

$$
\begin{equation*}
\operatorname{con}_{k}: M_{k}(\mathfrak{n}) \rightarrow C_{k}, \quad f \mapsto\left(c_{\mathcal{A}}(0, f)\right)_{[\mathcal{A}]} . \tag{24}
\end{equation*}
$$

For $k$ odd we must deal with the sign ambiguity in (15). For $k$ an odd integer, let

$$
\tilde{C}_{k}=\bigoplus_{\lambda \in \mathrm{Cl}^{+}(F)} \bigoplus_{\substack{\mathcal{A} \in T_{1, \lambda},(\mathfrak{n}) \backslash \mathbf{G L}_{2}^{+}(F) \\ \mathcal{A} \text { admissible }}} \mathbb{C} .
$$

Endow $\tilde{C}_{k}$ with a right action of the upper triangular Borel $B$ by

$$
[A] \cdot A^{\prime} \mapsto\left[A A^{\prime}\right] \operatorname{sgn}(\mathrm{N} \delta), \quad \text { where } \quad A^{\prime}=\left(\begin{array}{cc}
* & * \\
0 & \delta
\end{array}\right) .
$$

We let

$$
C_{k}=H_{0}\left(B, \tilde{C}_{k}\right)=\tilde{C}_{k} /\left\langle c \cdot b-c: c \in \tilde{C}_{k}, b \in B\right\rangle .
$$

Of course, Theorem 3.3 implies that

$$
C_{k} \cong \bigoplus_{\operatorname{cusps}^{*}(\mathfrak{n})} \mathbb{C}
$$

However, fixing such an isomorphism requires making a non-canonical choice, which we would like to avoid. For positive odd $k$ we again have a canonical constant term map

$$
\begin{equation*}
\operatorname{con}_{k}: M_{k}(\mathfrak{n}) \rightarrow C_{k} \tag{25}
\end{equation*}
$$

that sends a modular form $f=\left(f_{\lambda}\right)$ to the tuple of normalized constant terms $c_{\mathcal{A}}(0, f)$. The discussion above implies that this map is well-defined.
3.5. Cusps above $\infty$ and 0 . We introduce the suggestive notation

$$
C_{\infty}(\mathfrak{n})=\mathfrak{Q}_{\mathfrak{n}, \mathfrak{n}}=\left\{[\mathcal{A}] \in \operatorname{cusps}(\mathfrak{n}): \mathfrak{n} \mid \mathfrak{c}_{\mathcal{A}}\right\} .
$$

This is the smallest set of cusps containing the cusps $\infty \in \Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbb{P}^{1}(F)$ for each $\lambda \in \mathrm{Cl}^{+}(F)$ that is stable under the action of the diamond operators $S(\mathfrak{e})$. It follows from equation (23) with $\mathfrak{m}=\mathfrak{n}$ that $\# C_{\infty}(\mathfrak{n})=h^{+} h_{\mathfrak{n}}$, where $h_{\mathfrak{n}}=\# G_{\mathfrak{n}}$, the size of the wide ray class group of conductor $\mathfrak{n}$. Similarly, let

$$
C_{0}(\mathfrak{n})=\mathfrak{Q}_{1, \mathfrak{n}}=\left\{[\mathcal{A}] \in \operatorname{cusps}(\mathfrak{n}): \operatorname{gcd}\left(\mathfrak{n}, \mathfrak{c}_{\mathcal{A}}\right)=1\right\}
$$

This is the smallest $G_{\mathfrak{n}}^{+}$-invariant set of cusps containing $0 \in \Gamma_{1, \lambda}(\mathfrak{n}) \backslash \mathbb{P}^{1}(F)$ for each $\lambda \in \mathrm{Cl}^{+}(F)$.

If $\mathfrak{n}=1$, then $\operatorname{cusps}(\mathfrak{n})=C_{\infty}(\mathfrak{n})=C_{0}(\mathfrak{n})$. If $\mathfrak{n}$ is prime, then $\operatorname{cusps}(\mathfrak{n})$ is the disjoint union of $C_{\infty}(\mathfrak{n})$ and $C_{0}(\mathfrak{n})$. In general there will be more cusps.

If $\mathfrak{b}$ is a divisor of $\mathfrak{n}$, then for $*=0, \infty$ we define $C_{*}(\mathfrak{b}, \mathfrak{n})$ to be all elements in cusps( $\mathfrak{n}$ ) whose image under the canonical map cusps $(\mathfrak{n}) \rightarrow \operatorname{cusps}(\mathfrak{b})$ lies in $C_{*}(\mathfrak{b})$, i.e.

$$
C_{\infty}(\mathfrak{b}, \mathfrak{n})=\bigsqcup_{\mathfrak{b}|\mathfrak{m}| \mathfrak{n}} \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}, \quad C_{0}(\mathfrak{b}, \mathfrak{n})=\bigsqcup_{\substack{\mathfrak{m} \mid \mathfrak{n} \\(\mathfrak{b}, \mathfrak{m})=1}} \mathfrak{Q}_{\mathfrak{m}, \mathfrak{n}}
$$

3.6. Lemma on level raising and cusps. The following remark will be used in later computations.

Lemma 3.13. Let $f \in M_{k}(\mathfrak{n})$ and let $\mathfrak{q} \nmid \mathfrak{n}$ be a prime ideal. Let $\mathcal{A}=$ $(A, \lambda)$ represent a cusp $[\mathcal{A}] \in C_{\infty}(\mathfrak{q m}, \mathfrak{q n})$ for some $\mathfrak{m} \mid \mathfrak{n}$. There exists a pair $\mathcal{A}^{\prime}=\left(A^{\prime}, \mu\right)$ such that $\left[\mathcal{A}^{\prime}\right] \in C_{\infty}(\mathfrak{m}, \mathfrak{n})$ and

$$
\begin{equation*}
c_{\mathcal{A}}\left(0,\left.f\right|_{\mathfrak{q}}\right)=c_{\mathcal{A}^{\prime}}(0, f) \tag{26}
\end{equation*}
$$

Proof. Write $\mathcal{A}=(A, \lambda)$. By definition,

$$
\left.\left(\left(\left.f\right|_{\mathfrak{q}}\right)_{\lambda}\right)\right|_{A}=\left.\mathrm{Nq}^{-k / 2} f_{\mu}\right|_{A^{\prime}}, \quad \text { where } \quad A^{\prime}=\left(\begin{array}{cc}
\alpha_{\mu} \alpha & \alpha_{\mu} \beta \\
\gamma & \delta
\end{array}\right)
$$

and $\mathfrak{q t}_{\lambda}=\left(\alpha_{\mu}\right) \mathfrak{t}_{\mu}$. Let $\mathcal{A}^{\prime}=\left(A^{\prime}, \mu\right)$. Then $\left[\mathcal{A}^{\prime}\right] \in C_{\infty}(\mathfrak{m}, \mathfrak{n})$ and $c_{\mathcal{A}}\left(0,\left.f\right|_{\mathfrak{q}}\right)$ $=c_{\mathcal{A}^{\prime}}(0, f)$ are direct calculations using (11).
4. Eisenstein series. We recall the well-known Eisenstein series using [2. §2.2] and [1] as convenient references. Let $\eta, \chi$ be narrow ray class characters of modulus $\mathfrak{a}, \mathfrak{b}$, respectively, such that $\eta \chi$ has $\operatorname{sign}(k, \ldots, k)$. With the exception of the case $F=\mathbb{Q}, k=2, \eta^{0}=\chi^{0}=1$, when there are convergence issues, there is an Eisenstein series $E_{k}(\eta, \chi) \in M_{k}(\mathfrak{a b})$. Here $\eta, \chi$ are not assumed to be primitive characters, and $\eta^{0}, \chi^{0}$ denote the primitive characters associated to $\eta, \chi$. The Eisenstein series $E_{k}(\eta, \chi)$ has $q$-expansion
coefficients given by

$$
c\left(\mathfrak{m}, E_{k}(\eta, \chi)\right)=\sum_{\mathfrak{r} \mid \mathfrak{m}} \eta(\mathfrak{m} / \mathfrak{r}) \chi(\mathfrak{r}) \mathrm{Nr}^{k-1}
$$

For $k>1$, we have

$$
c_{\lambda}\left(0, E_{k}(\eta, \chi)\right)= \begin{cases}2^{-n} \eta^{-1}(\lambda) L\left(\chi \eta^{-1}, 1-k\right) & \text { if } \mathfrak{a}=1 \\ 0 & \text { if } \mathfrak{a} \neq 1\end{cases}
$$

For $k=1, \mathfrak{a}=1$, and $\eta=1$, we note that

$$
c_{\lambda}\left(0, E_{1}(1, \chi)\right)=2^{-n} \cdot \begin{cases}L(\chi, 0) & \text { if } \mathfrak{b} \neq 1, \\ L(\chi, 0)+\chi^{-1}\left(\mathfrak{t}_{\lambda}\right) L\left(\chi^{-1}, 0\right) & \text { if } \mathfrak{b}=1\end{cases}
$$

for all $\lambda \in \mathrm{Cl}^{+}(F)$.
Given a fixed level $\mathfrak{n}$, the Eisenstein subspace $E_{k}(\mathfrak{n}) \subset M_{k}(\mathfrak{n})$ is defined to be the subspace spanned by the Eisenstein series $\left.E_{k}(\eta, \chi)\right|_{\mathfrak{q}}$ where $\eta, \chi$ are primitive narrow ray class characters of conductor $\mathfrak{a}, \mathfrak{b}$, respectively, such that $\eta \chi$ has $\operatorname{sign}(k, \ldots, k)$ and $\mathfrak{a b q} \mid \mathfrak{n}$. An elementary argument using Hecke operators shows that for $k>1$, these Eisenstein series are linearly independent (see [1, Prop. 3.8]). For $k=1$, we have $\left.E_{1}(\eta, \chi)\right|_{\mathfrak{q}}=\left.E_{1}(\chi, \eta)\right|_{\mathfrak{q}}$, and these equations generate the space of relations among the Eisenstein series.

The Eisenstein subspace is a complement to the space $S_{k}(\mathfrak{n})$ of cusp forms, i.e.

$$
\begin{equation*}
M_{k}(\mathfrak{n})=E_{k}(\mathfrak{n}) \oplus S_{k}(\mathfrak{n}) . \tag{27}
\end{equation*}
$$

Furthermore, for $k \geq 2$ (excluding the case $F=\mathbb{Q}, k=2$ ) the restriction of the map $\operatorname{con}_{k}$ defined in (24) and (25) to the subspace $E_{k}(\mathfrak{n})$ is an isomorphism:

$$
\begin{equation*}
\operatorname{con}_{k, E}: E_{k}(\mathfrak{n}) \xrightarrow{\sim} C_{k} . \tag{28}
\end{equation*}
$$

The results (27) and (28) are proven in [12, Prop. 1.5] for weight $k=2$, and we now sketch a proof in the general case. Firstly, one can show that (excluding the case $F=\mathbb{Q}, k=2$ ) there is an equality of dimensions in (28):

$$
\operatorname{dim}_{\mathbb{C}} E_{k}(\mathfrak{n})= \begin{cases}\# \operatorname{cusps}(\mathfrak{n})=\sum_{\mathfrak{m} \mid \mathfrak{n}} \# Q_{\mathfrak{m}, \mathfrak{n}} & \text { if } k \geq 2 \text { is even, }  \tag{29}\\ \# \operatorname{cusps}^{*}(\mathfrak{n})=\sum_{\mathfrak{m} \mid \mathfrak{n}}^{*} \# Q_{\mathfrak{m}, \mathfrak{n}} & \text { if } k \geq 3 \text { is odd }\end{cases}
$$

To see this for $k$ even, note that $\# Q_{\mathfrak{m}, \mathfrak{n}}$ is the number of pairs $\left(\chi_{1}, \chi_{2}\right)$ where $\chi_{1}, \chi_{2}$ are ray class characters of modulus $\mathfrak{m}, \mathfrak{n} / \mathfrak{m}$, respectively, such that $\chi_{1} \chi_{2}$ is totally even. To such a pair we can associate the Eisenstein series $\left.E_{k}\left(\chi_{1}^{0}, \chi_{2}^{0}\right)\right|_{\mathfrak{m} / \operatorname{cond}\left(\chi_{1}\right)}$. Here $\chi_{1}^{0}$ and $\chi_{2}^{0}$ denote the primitive avatars of $\chi_{1}$ and $\chi_{2}$, respectively. These Eisenstein series form the defining basis for $E_{k}(\mathfrak{n})$.

For $k \geq 3$ odd, it is not hard to show that there exists a pair of characters $\chi_{1}, \chi_{2}$ of modulus $\mathfrak{m}, \mathfrak{n} / \mathfrak{m}$, respectively, with $\chi_{1} \chi_{2}$ totally odd if and only if $(\mathfrak{m}, \mathfrak{n})$ is an admissible pair. In this case, $\# Q_{\mathfrak{m}, \mathfrak{n}}$ is equal to the size of the set of such pairs $\left(\chi_{1}, \chi_{2}\right)$. The argument then continues as in the case for $k$ even, and we find that the dimension of $E_{k}(\mathfrak{n})$ is $\sum_{\mathfrak{m} \mid \mathfrak{n}}^{*} \# Q_{\mathfrak{m}, \mathfrak{n}}$.

With (29) in hand, both (27) and (28) for $k \geq 2$ follow from the fact that

$$
E_{k}(\mathfrak{n}) \cap S_{k}(\mathfrak{n})=\{0\}
$$

which is usually proven using the Petersson inner product (see for instance [1, Prop. 3.9] or [12, p. 423]).
4.1. Evaluation of constant terms of Eisenstein series. Let $\chi_{1}$ and $\chi_{2}$ be ray class characters of modulus $\mathfrak{a}$ and $\mathfrak{b}$ and signature $q_{1}$ and $q_{2}$, respectively. Put $\mathfrak{n}=\mathfrak{a b}$. Let $k$ be a positive integer and assume that $q_{1}+q_{2} \equiv(k, \ldots, k)(\bmod 2)$. In this section we compute the constant terms of Eisenstein series $E_{k}\left(\chi_{1}, \chi_{2}\right)$ at various cusps $\mathcal{A}=(A, \lambda)$, where $A=\binom{\alpha *}{\gamma *} \in \mathbf{G L}_{2}^{+}(F)$.

We write $\mathfrak{a}_{0}=\operatorname{cond}\left(\chi_{1}\right), \mathfrak{b}_{0}=\operatorname{cond}\left(\chi_{2}\right)$ and let $\mathfrak{a}_{1}=\mathfrak{a} / \mathfrak{a}_{0}, \mathfrak{b}_{1}=\mathfrak{b} / \mathfrak{b}_{0}$. Without loss of generality, we assume that $\operatorname{gcd}\left(\mathfrak{a}_{0}, \mathfrak{a}_{1}\right)=1$ and that $\mathfrak{a}_{1}$ is square-free, since increasing the modulus at a prime already dividing the conductor or increasing the power of a prime already dividing the modulus does not affect the character or associated Eisenstein series. We make the same assumptions about $\mathfrak{b}_{1}$.

Definition 4.1. The Gauss sum associated to a primitive character $\chi$ of conductor $\mathfrak{b}$ and $\operatorname{sign} r \in(\mathbb{Z} / 2 \mathbb{Z})^{n}$ is given by

$$
\tau(\chi)=\sum_{x \in \mathfrak{b}^{-1} \mathfrak{d}^{-1} / \mathfrak{d}^{-1}} \operatorname{sgn}(x)^{r} \chi(x \mathfrak{b d}) e_{F}(x)
$$

For a general character $\chi$ we define $\tau(\chi)=\tau\left(\chi^{0}\right)$ where $\chi^{0}$ is the associated primitive character.

Recall the invariants $\mathfrak{a}_{\mathcal{A}}, \mathfrak{b}_{\mathcal{A}}, \mathfrak{c}_{\mathcal{A}}$ defined in $\$ 3.1$.
Definition 4.2. Assume that $[\mathcal{A}] \in C_{\infty}(\mathfrak{b}, \mathfrak{n})$, i.e. $\mathfrak{b} \mid \mathfrak{c}_{\mathcal{A}}$. Write $\mathfrak{b}_{0}=$ $\operatorname{cond}\left(\chi_{2}\right)$. Let $\chi$ denote the primitive character associated to $\chi_{1} \chi_{2}^{-1}$, and write $\mathfrak{f}=\operatorname{cond}(\chi)$. Define

$$
\begin{align*}
& \text { (30) } \quad P_{\mathcal{A}}\left(\chi_{1}, \chi_{2}, k\right)=  \tag{30}\\
& \frac{1}{2^{n}} \frac{\tau\left(\chi_{1} \chi_{2}^{-1}\right)}{\tau\left(\chi_{2}^{-1}\right)}\left(\frac{\mathrm{N} \mathfrak{b}_{0}}{\mathrm{Nf}}\right)^{k} \operatorname{sgn}(-\gamma)^{q_{1}} \operatorname{sgn}(\alpha)^{q_{2}} \chi_{1}\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}_{0}\right)\left(\chi_{2}^{0}\right)^{-1}\left(\mathfrak{a}_{\mathcal{A}}\right) L\left(\chi^{-1}, 1-k\right)
\end{align*}
$$

where $\chi_{2}^{0}$ denotes the primitive avatar of $\chi_{2}$. Here and throughout, we adopt the convention that $\chi_{1}(\mathfrak{m})=0$ if $\operatorname{gcd}(\mathfrak{m}, \mathfrak{a}) \neq 1$, and similarly for any ray class character. We also use the convention that if $a=0$ and $\mathfrak{m}$ is a fractional
ideal, then

$$
\operatorname{sgn}(a)^{q_{1}} \chi_{1}(a \mathfrak{m})= \begin{cases}0 & \text { if } \mathfrak{a} \neq 1 \\ \chi_{1}(\mathfrak{m}) & \text { if } \mathfrak{a}=1\end{cases}
$$

For example, when $\gamma=0$ in 30 the expression $\operatorname{sgn}(-\gamma)^{q_{1}} \chi_{1}\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}\right)$ should be interpreted as 0 if $\mathfrak{a} \neq 1$ and as $\chi_{1}^{-1}\left(\mathfrak{t}_{\lambda} \mathfrak{d b}_{A} \mathfrak{b}\right)$ if $\mathfrak{a}=1$. The analogous convention holds for the term $\operatorname{sgn}(\alpha)^{q_{2}} \chi_{2}^{-1}\left(\mathfrak{a}_{\mathcal{A}}\right)=\operatorname{sgn}(\alpha)^{q_{2}} \chi_{2}^{-1}\left(\alpha \mathfrak{b}_{\mathcal{A}}^{-1}\right)$.

For finite sets $S$ and $T$ of finite places of $F$, define

$$
P_{\mathcal{A}}\left(\chi_{1}, \chi_{2}, k, S, T\right)=P_{\mathcal{A}}\left(\chi_{1}, \chi_{2}, k\right) \prod_{\mathfrak{q} \in S}\left(1-\chi^{-1}(\mathfrak{q}) \mathrm{Nq}^{k-1}\right) \prod_{\mathfrak{q} \in T}\left(1-\chi(\mathfrak{q}) \mathrm{Nq}^{-k}\right)
$$

Further, for an ideal $\mathfrak{a} \mid \mathfrak{n}$ define

$$
\delta_{0, \mathcal{A}}(\mathfrak{a})=\left\{\begin{array}{ll}
0 & \text { if }[\mathcal{A}] \notin C_{0}(\mathfrak{a}, \mathfrak{n}), \\
1 & \text { if }[\mathcal{A}] \in C_{0}(\mathfrak{a}, \mathfrak{n}),
\end{array} \quad \delta_{\infty, \mathcal{A}}(\mathfrak{a})= \begin{cases}0 & \text { if }[\mathcal{A}] \notin C_{\infty}(\mathfrak{a}, \mathfrak{n}) \\
1 & \text { if }[\mathcal{A}] \in C_{\infty}(\mathfrak{a}, \mathfrak{n})\end{cases}\right.
$$

REMARK 4.3. Suppose that $\mathfrak{a}_{0}=\operatorname{cond}\left(\chi_{1}\right)$ and $\mathfrak{b}_{0}=\operatorname{cond}\left(\chi_{2}\right)$ are coprime. Then

$$
\tau\left(\chi_{1} \chi_{2}^{-1}\right)=\chi_{1}\left(\mathfrak{b}_{0}\right) \chi_{2}^{-1}\left(\mathfrak{a}_{0}\right) \tau\left(\chi_{1}\right) \tau\left(\chi_{2}^{-1}\right)
$$

and hence

$$
\begin{align*}
& P_{\mathcal{A}}\left(\chi_{1}, \chi_{2}, k\right)  \tag{31}\\
& \quad=\frac{\tau\left(\chi_{1}\right)}{2^{n}\left(\mathrm{Na}_{0}^{k}\right)} \operatorname{sgn}(-\gamma)^{q_{1}} \operatorname{sgn}(\alpha)^{q_{2}} \chi_{1}\left(\mathfrak{c}_{\mathcal{A}}\right) \chi_{2}^{-1}\left(\mathfrak{a}_{\mathcal{A}} / \mathfrak{a}_{0}\right) L\left(\chi^{-1}, 1-k\right)
\end{align*}
$$

We require one more piece of notation.
Definition 4.4. Let $\mathfrak{m} \mid \mathfrak{n}$. We write $J_{\mathfrak{m}}=J_{\mathfrak{m}}(\mathcal{A})$ for the set of prime divisors $\mathfrak{q} \mid \mathfrak{m}$ such that $[\mathcal{A}] \in C_{0}(\mathfrak{q}, \mathfrak{n})$, and $J_{\mathfrak{m}}^{c}$ for the set of prime divisors $\mathfrak{q} \mid \mathfrak{m}$ such that $[\mathcal{A}] \in C_{\infty}(\mathfrak{q}, \mathfrak{n})$.

Theorem 4.5. Assume that $\chi_{2}$ is primitive of conductor $\mathfrak{b}$. The normalized constant terms of $E_{k}\left(\chi_{1}, \chi_{2}\right)$ at $\mathcal{A}$ are given as follows:
(1) Let $k>1$. The normalized constant term of $E_{k}\left(\chi_{1}, \chi_{2}\right)$ at $\mathcal{A}$ equals

$$
\begin{equation*}
\delta_{\infty, \mathcal{A}}(\mathfrak{b}) P_{\mathcal{A}}\left(\chi_{1}, \chi_{2}, k, \emptyset, T_{\mathfrak{n}, \mathfrak{f}}\right) \tag{32}
\end{equation*}
$$

where $T_{\mathfrak{n}, \mathfrak{f}}$ is the set of primes dividing $\mathfrak{n}$ but not $\mathfrak{f}$.
(2) Let $k=1$. Suppose that $\mathfrak{a}$ and $\mathfrak{b}$ are coprime. The normalized constant term of $E_{1}\left(\chi_{1}, \chi_{2}\right)$ at a cusp $\mathcal{A}$ equals

$$
\begin{align*}
& \delta_{0, \mathcal{A}}(\mathfrak{a}) \delta_{\infty, \mathcal{A}}(\mathfrak{b}) P_{\mathcal{A}}\left(\chi_{1}, \chi_{2}, 1, \emptyset, T_{\mathfrak{n}, \mathfrak{f}}\right)  \tag{33}\\
&+\delta_{\infty, \mathcal{A}}\left(\mathfrak{a}_{0}\right) \delta_{0, \mathcal{A}}(\mathfrak{b}) P_{\mathcal{A}}\left(\chi_{2}, \chi_{1}, 1, J_{\mathfrak{a}_{1}}^{c}, \emptyset\right) \prod_{\mathfrak{q} \in J_{\mathfrak{a}_{1}}}\left(1-\mathrm{Nq}^{-1}\right)
\end{align*}
$$

REmARK 4.6. Note that $E_{1}\left(\chi_{1}, \chi_{2}\right)=E_{1}\left(\chi_{2}, \chi_{1}\right)$. The theorem only assumes that $\chi_{2}$ is primitive; if we assume also that $\chi_{1}$ is primitive, then in
the setting of part (2) the sets $T_{\mathfrak{n}, \mathfrak{f}}, J_{\mathfrak{a}_{1}}$ and $J_{\mathfrak{a}_{1}}^{c}$ are empty and (33) becomes symmetric with respect to $\chi_{1}$ and $\chi_{2}$.

REmARK 4.7. The term $\delta_{0, \mathcal{A}}(\mathfrak{a})$ is unnecessary in 47) below since $P_{\mathcal{A}}(\chi, \psi, 1)$ already vanishes if $[\mathcal{A}] \notin C_{0}(\mathfrak{a}, \mathfrak{n})$. In this case the assumption $\operatorname{gcd}(\mathfrak{a}, \mathfrak{b})=1$ implies that $\operatorname{gcd}\left(\mathfrak{a}, \mathfrak{c}_{\mathcal{A}}\right)$ is non-trivial, whence $\chi_{1}\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}_{0}\right)=0$. We include this factor simply as a reminder that this portion of the constant term is supported on $C_{0}(\mathfrak{a}, \mathfrak{n}) \cap C_{\infty}\left(\mathfrak{b}_{0}, \mathfrak{n}\right)$.

Proof of Theorem 4.5. Recall from [10, Prop. 3.4] the definition of the Eisenstein series $E_{k}\left(\chi_{1}, \chi_{2}\right)$ (see also [2, Section 2.2]). Let

$$
U=\left\{u \in \mathcal{O}_{F}^{*}: \mathrm{N} u^{k}=1, u \equiv 1(\bmod \mathfrak{n})\right\} .
$$

For $k \geq 1$, we have $E_{k}\left(\chi_{1}, \chi_{2}\right)=\left(f_{\lambda}(z)\right)_{\lambda \in \mathrm{Cl}^{+}(F)}$, where $f_{\lambda}(z)$ is defined via Hecke's trick as follows. For $z \in \mathcal{H}$ and $s \in \mathbb{C}$ with $\operatorname{Re}(2 s+k)>2$, define

$$
\begin{equation*}
f_{\lambda}(z, s)=C_{\lambda} \tau\left(\chi_{2}\right) \sum_{\mathfrak{r} \in \mathrm{Cl}(F)} \mathrm{Nr}^{k} g_{\lambda}(z, s) \tag{34}
\end{equation*}
$$

where

$$
C_{\lambda}=\frac{\sqrt{d_{F}} \Gamma(k)^{n} \mathrm{~N}\left(\mathfrak{t}_{\lambda}\right)^{-k / 2}}{\left[\mathcal{O}_{F}^{*}: U\right] \mathrm{N}(\mathfrak{d}) \mathrm{N}(\mathfrak{b})(-2 \pi i)^{k n}}
$$

and

$$
\begin{align*}
g_{\lambda}(z, s)= & \sum_{a, b} \frac{\operatorname{sgn}(a)^{q_{1}} \chi_{1}\left(a \mathfrak{r}^{-1}\right) \operatorname{sgn}(-b)^{q_{2}} \chi_{2}^{-1}\left(-b \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right)}{(a z+b)^{k}|a z+b|^{2 s}}  \tag{35}\\
= & \sum_{\left(a_{0}, b_{0}\right)} \operatorname{sgn}\left(a_{0}\right)^{q_{1}} \chi_{1}\left(a_{0} \mathfrak{r}^{-1}\right) \operatorname{sgn}\left(-b_{0}\right)^{q_{2}} \chi_{2}^{-1}\left(-b_{0} \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right)  \tag{36}\\
& \times \sum_{(a, b) \equiv\left(a_{0}, b_{0}\right)} \frac{1}{(a z+b)^{k}|a z+b|^{2 s}} .
\end{align*}
$$

The sum in (34) runs over representatives $\mathfrak{r}$ for the wide class group $\mathrm{Cl}(F)$. The sum in (35) runs over representatives $(a, b)$ for the non-zero elements of the product $\mathfrak{r} \times \mathfrak{d}^{-1} \mathfrak{b}^{-1} \mathfrak{t}_{\lambda}^{-1} \mathfrak{r}$ modulo the diagonal action of $U$. In the equality (36) the $\left(a_{0}, b_{0}\right)$ sum runs through $(\mathfrak{r} / \mathfrak{r a b}) \times\left(\mathfrak{d}^{-1} \mathfrak{b}^{-1} \mathfrak{t}_{\lambda}^{-1} \mathfrak{r} / \mathfrak{d}^{-1} \mathfrak{t}_{\lambda}^{-1} \mathfrak{r a}\right)$, while $(a, b)$ ranges over the non-zero elements of $\mathfrak{r} \times \mathfrak{d}^{-1} \mathfrak{b}^{-1} \mathfrak{t}_{\lambda}^{-1} \mathfrak{r}$ modulo the diagonal action of $U$ such that $a \equiv a_{0}(\bmod \mathfrak{r a b})$ and $b \equiv b_{0}\left(\bmod \mathfrak{d}^{-1} \mathfrak{t}_{\lambda}^{-1} \mathfrak{r a}\right)$.

Here we use

$$
\begin{array}{ll}
\chi_{1}(a)=\operatorname{sgn}(a)^{q_{1}} & \text { for } a \equiv 1(\bmod \mathfrak{a}) \\
\chi_{2}(b)=\operatorname{sgn}(b)^{q_{2}} & \text { for } b \equiv 1(\bmod \mathfrak{b})
\end{array}
$$

We remark that in the definition we already use the assumption that $\chi_{2}$ is primitive, applying [10, (3.11)]). The function $f_{\lambda}(z, s)$ can be analytically continued in the variable $s$ to the entire complex plane, and we set $f_{\lambda}(z)$ $=f_{\lambda}(z, 0)$.

We choose representatives of the cusp $[\mathcal{A}]=[(A, \lambda)]$ with $A=\left(\begin{array}{ll}\alpha & \beta \\ \gamma & \delta\end{array}\right)$, as follows. Let $\mathfrak{g}=\operatorname{gcd}\left(\mathfrak{b}, \mathfrak{c}_{\mathcal{A}}\right)$. The cusp $[\mathcal{A}]$ only depends on $(\alpha, \gamma)$, so for convenience we are free to choose $\beta, \delta$ such that

$$
\begin{equation*}
\operatorname{det}(A)=1, \quad \beta \in\left(\mathfrak{t}_{\lambda} \mathfrak{d} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}, \quad \delta \in \mathfrak{b}\left(\mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \tag{37}
\end{equation*}
$$

Such $\beta, \delta$ exist because (using notation from (13) ) we have

$$
\begin{equation*}
\alpha \mathfrak{b}\left(\mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}+\gamma\left(\mathfrak{t}_{\lambda} \mathfrak{d} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}=\mathfrak{a}_{\mathcal{A}} \mathfrak{b} \mathfrak{g}^{-1}+\mathfrak{c}_{\mathcal{A}} \mathfrak{g}^{-1}=\mathcal{O}_{F} \tag{38}
\end{equation*}
$$

The last equality holds as $\mathfrak{c}_{\mathcal{A}} \mathfrak{g}^{-1}$ is coprime to both $\mathfrak{b g}{ }^{-1}$ and $\mathfrak{a}_{\mathcal{A}}$. The map $(a, b) \mapsto(u, v)=(a, b) A$ induces a bijection

$$
\begin{equation*}
\mathfrak{r} \times \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}\right)^{-1} \rightarrow \mathfrak{r} \mathfrak{b}_{\mathcal{A}} \mathfrak{g} \mathfrak{b}^{-1} \times \mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \tag{39}
\end{equation*}
$$

To see this, note that

$$
(u, v)=(a \alpha+b \gamma, a \beta+b \delta) \quad \text { and } \quad(a, b)=(u \delta-v \gamma,-u \beta+v \alpha)
$$

For $(a, b) \in \mathfrak{r} \times \mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}\right)^{-1}$ we have

$$
u=a \alpha+b \gamma \in \mathfrak{r b} \mathcal{A}_{\mathcal{A}} \mathfrak{b}^{-1}\left(\alpha \mathfrak{b}\left(\mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}+\gamma\left(\mathfrak{t}_{\lambda} \mathfrak{d b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}\right)=\mathfrak{r b}_{\mathcal{A}} \mathfrak{g b} \mathfrak{b}^{-1}
$$

by (38). Meanwhile

$$
v=a \beta+b \delta \in \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}+\mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}\right)^{-1} \mathfrak{b}\left(\mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}=\mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}
$$

For the other direction, note that for $(u, v) \in \mathfrak{r b}_{\mathcal{A}} \mathfrak{g b}^{-1} \times \mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}$ we have

$$
\begin{aligned}
a & =u \delta-v \gamma \in \mathfrak{r b}_{\mathcal{A}} \mathfrak{g b}^{-1} \mathfrak{b}\left(\mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}+\mathfrak{r g}^{-1} \gamma\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)^{-1}=\mathfrak{r}+\mathfrak{r}\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{g}\right)=\mathfrak{r} \\
b & =-u \beta+v \alpha \in \mathfrak{r b}_{\mathcal{A}} \mathfrak{g b}^{-1}\left(\mathfrak{d t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}+\mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \alpha=\mathfrak{r}\left(\mathfrak{d t}_{\lambda} \mathfrak{b}\right)^{-1}
\end{aligned}
$$

This completes the proof that 39 is a bijection. This bijection restricts to a bijection

$$
\mathfrak{r a b} \times \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda}\right)^{-1} \mathfrak{a} \rightarrow \mathfrak{r b}_{\mathcal{A}} \mathfrak{g a} \times \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \mathfrak{n}
$$

The function $\left.g_{\lambda}(z, s)\right|_{A}$ can be written as

$$
\begin{align*}
&\left.g_{\lambda}(z, s)\right|_{A}=\sum_{u_{0}, v_{0}} \operatorname{sgn}\left(u_{0} \delta-v_{0} \gamma\right)^{q_{1}} \chi_{1}\left(\left(u_{0} \delta-v_{0} \gamma\right) \mathfrak{r}^{-1}\right)  \tag{40}\\
& \times \operatorname{sgn}\left(u_{0} \beta\right.\left.-v_{0} \alpha\right)^{q_{2}} \chi_{2}^{-1}\left(\left(u_{0} \beta-v_{0} \alpha\right) \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right) \\
& \times \sum_{(u, v) \equiv\left(u_{0}, v_{0}\right)} \frac{1}{(u z+v)^{k}|u z+v|^{2 s}}
\end{align*}
$$

Here $u_{0}$ and $v_{0}$ run through complete sets of representatives of

$$
\mathfrak{r} \mathfrak{b}_{\mathcal{A}} \mathfrak{g} \mathfrak{b}^{-1} / \mathfrak{r} \mathfrak{b}_{\mathcal{A}} \mathfrak{g a} \quad \text { and } \quad \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} / \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \mathfrak{n}
$$

respectively, and the pair $(u, v)$ runs through representatives for

$$
\left(\left(\mathfrak{r b} \mathfrak{A}_{\mathcal{A}} \mathfrak{g} \mathfrak{b}^{-1} \times \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1}\right) \backslash\{(0,0)\}\right) / U
$$

such that $(u, v) \equiv\left(u_{0}, v_{0}\right)$. We now recall notation from [10, §3]. Up to a constant factor the sum

$$
\sum_{(u, v) \equiv\left(u_{0}, v_{0}\right)} \frac{1}{(u z+v)^{k}|u z+v|^{2 s}}
$$

equals the series $E_{k, U}\left(z, u_{0}, v_{0}, \mathfrak{r b}_{\mathcal{A}} \mathfrak{g a}, \mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \mathfrak{n}\right)$ defined in [10, (3.1)], with $r$ in loc. cit. set to 0 .

We are now ready to prove (1), though we remark that much of what is said below also applies to (2). We have $k \geq 2$. By [10, (3.7)] the constant term of $\left.g_{\lambda}(z, s)\right|_{A}$ at $s=0$ is the value at $s=0$ of

$$
\begin{aligned}
\sum_{u_{0}, v_{0}} \operatorname{sgn}\left(u_{0} \delta-v_{0} \gamma\right)^{q_{1}} \chi_{1} & \left(\left(u_{0} \delta-v_{0} \gamma\right) \mathfrak{r}^{-1}\right) \\
& \times \operatorname{sgn}\left(u_{0} \beta-v_{0} \alpha\right)^{q_{2}} \chi_{2}^{-1}\left(\left(u_{0} \beta-v_{0} \alpha\right) \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right) \\
& \times(-1)^{k n} \delta\left(u_{0}, \mathfrak{r b}_{\mathcal{A}} \mathfrak{g a}\right) \sum_{v \equiv v_{0}} \operatorname{sgn}(\mathrm{~N} v)^{k}|\mathrm{~N} v|^{-k-2 s}
\end{aligned}
$$

Here $\delta\left(u_{0}, \mathfrak{r b}_{\mathcal{A}} \mathfrak{g a}\right)$ is 0 if $u_{0} \notin \mathfrak{r b}_{\mathcal{A}} \mathfrak{g a}$ and is 1 otherwise. Therefore only the summands where the class of $u_{0}$ is trivial contribute to the constant term of $\left.g_{\lambda}(z, s)\right|_{A}$ at $s=0$.

If $u_{0} \in \mathfrak{r b}_{\mathcal{A}} \mathfrak{g a}$, then using the relation $b_{0} \gamma=u_{0}-a_{0} \alpha$ we deduce that $b_{0} \gamma \in \mathfrak{r b}_{\mathcal{A}}$. On the other hand $b_{0} \gamma \in \mathfrak{r b}_{\mathcal{A}} \mathfrak{c}_{\mathcal{A}} \mathfrak{b}^{-1}$. Therefore

From this we obtain $b_{0} \in \mathfrak{r}\left(\mathfrak{t}_{\lambda} \mathfrak{d g}\right)^{-1}$. Now consider the case that $\mathcal{A}$ does not represent a cusp in $C_{\infty}(\mathfrak{b}, \mathfrak{n})$. This is equivalent to $\mathfrak{b} \nmid \mathfrak{c}_{\mathcal{A}}$, which in turn is equivalent to $\mathfrak{g} \neq \mathfrak{b}$. Hence $b_{0} \mathfrak{b d t} \mathfrak{r}^{-1}$ is an integral ideal not coprime to $\mathfrak{b}$. Therefore $\chi_{2}^{-1}\left(-b_{0} \mathfrak{b d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right)=0$ and thus the terms where the class of $u_{0}$ is trivial also vanish, so the constant term of $\left.g_{\lambda}(z, s)\right|_{A}$ at $s=0$ is 0 . Note that here we have used $b_{0}=-u_{0} \beta+v_{0} \alpha$.

Next we turn to the case that $\mathcal{A}$ does represent a cusp in $C_{\infty}(\mathfrak{b}, \mathfrak{n})$, so $\mathfrak{g}=\mathfrak{b}$. As observed above we are only interested in the term with $u_{0} \in \mathfrak{r b} \mathcal{A}_{\mathcal{A}} \mathfrak{a}$. We choose $u_{0}=0$ to represent the trivial coset in $\mathfrak{r b}_{\mathcal{A}} \mathfrak{g b}^{-1} / \mathfrak{r b}_{\mathcal{A}} \mathfrak{g a}$. Therefore the constant term of $\left.g_{\lambda}(z, s)\right|_{A}$ at $s=0$ is the value at $s=0$ of

$$
\begin{aligned}
& \sum_{v_{0}} \operatorname{sgn}\left(\mathrm{~N}\left(-v_{0}\right)\right)^{k} \operatorname{sgn}(\gamma)^{q_{1}} \operatorname{sgn}(\alpha)^{q_{2}} \\
& \quad \times \chi_{1}\left(-v_{0} \gamma \mathfrak{r}^{-1}\right) \chi_{2}^{-1}\left(-v_{0} \alpha \mathfrak{b} \mathfrak{d}_{\lambda} \mathfrak{r}^{-1}\right)(-1)^{k n} \sum_{v \equiv v_{0}} \frac{\operatorname{sgn}(\mathrm{~N} v)^{k}}{|\mathrm{~N} v|^{k+2 s}} \\
& =\sum_{v} \operatorname{sgn}(\gamma)^{q_{1}} \operatorname{sgn}(\alpha)^{q_{2}} \chi_{1}\left(-v \gamma \mathfrak{r}^{-1}\right) \chi_{2}^{-1}\left(-v \alpha \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right) \mathrm{N}\left(v \mathcal{O}_{F}\right)^{-k-2 s}
\end{aligned}
$$

Hence the constant term of $\left.f_{\lambda}(z, s)\right|_{A}$ at $s=0$ is the value of the following
sum at $s=0$ :

$$
\begin{aligned}
& C_{\lambda} \tau\left(\chi_{2}\right) \sum_{\mathfrak{r} \in \mathrm{Cl}(F)} \mathrm{Nr}^{k} \sum_{v} \operatorname{sgn}(\gamma)^{q_{1}} \operatorname{sgn}(\alpha)^{q_{2}} \\
& \quad \times \chi_{1}\left(-v \gamma \mathfrak{r}^{-1}\right) \chi_{2}^{-1}\left(-v \alpha \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right) \mathrm{N}\left(v \mathcal{O}_{F}\right)^{-k-2 s} \\
& = \\
& C_{\lambda} \tau\left(\chi_{2}\right) \operatorname{sgn}(\gamma)^{q_{1}} \operatorname{sgn}(\alpha)^{q_{2}}(-1)^{-k n-2 n s} \\
& \quad \times \chi_{1}\left(\gamma\left(\mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)^{-1}\right) \chi_{2}^{-1}\left(\alpha\left(\mathfrak{b}_{\mathcal{A}}\right)^{-1}\right) \mathrm{N}\left(\mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)^{k+2 s} \\
& \quad \times \sum_{\mathfrak{r}} \sum_{v} \chi_{1}\left(v \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{r}^{-1}\right) \chi_{2}^{-1}\left(v \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{r}^{-1}\right) \mathrm{N}\left(v \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)^{-k-2 s} \mathrm{~N}\left(\mathfrak{r}^{-1}\right)^{-k}
\end{aligned}
$$

Here we have made the change of variable $v \mapsto-v$. The value of this at $s=0$ is

$$
\begin{aligned}
& C_{\lambda} \tau\left(\chi_{2}\right) \operatorname{sgn}(\gamma)^{q_{1}} \operatorname{sgn}(\alpha)^{q_{2}}(-1)^{-k n} \chi_{1}\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}\right) \chi_{2}^{-1}\left(\mathfrak{a}_{\mathcal{A}}\right) \mathrm{N}\left(\mathfrak{b} \mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)^{k} \\
& \times\left[\mathcal{O}_{F}^{*}: U\right] L(\chi, k) \prod_{\mathfrak{q} \in T_{\mathfrak{n}, \mathfrak{f}}}\left(1-\chi(\mathfrak{q}) \mathrm{Nq}^{-k}\right)
\end{aligned}
$$

where $\chi$ is the primitive character associated to $\chi_{1} \chi_{2}^{-1}$, and $\mathfrak{q}$ runs through the set $T_{\mathfrak{n}, \mathfrak{f}}$ of all primes dividing $\mathfrak{n}$ but not $\mathfrak{f}=\operatorname{cond}(\chi)$. Next we use the functional equation

$$
L(\chi, k)=\frac{d(F)^{1 / 2-k} N \mathfrak{f}^{1-k}(2 \pi i)^{k n}}{2^{n} \Gamma(k)^{n} \tau\left(\chi^{-1}\right)} L\left(\chi^{-1}, 1-k\right)
$$

together with the relations

$$
\tau\left(\chi_{2}\right) \tau\left(\chi_{2}^{-1}\right)=\operatorname{sgn}(-1)^{q_{2}} \mathrm{Nb}, \quad \tau(\chi) \tau\left(\chi^{-1}\right)=\operatorname{sgn}(-1)^{q_{1}+q_{2}} \mathrm{Nf}
$$

We find that the unnormalized constant term $a_{\lambda, A}(0)$ of $\left.f_{\lambda}(z, 0)\right|_{A}$ is

$$
\begin{aligned}
& \frac{\tau\left(\chi_{1} \chi_{2}^{-1}\right)}{\tau\left(\chi_{2}^{-1}\right)}\left(\frac{\mathrm{Nbb}_{\mathcal{A}}}{\mathrm{Nf}}\right)^{k} \mathrm{Nt}_{\lambda}^{k / 2} \operatorname{sgn}(-\gamma)^{q_{1}} \operatorname{sgn}(\alpha)^{q_{2}} \chi_{1}\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}\right) \chi_{2}^{-1}\left(\mathfrak{a}_{\mathcal{A}}\right) \\
& \times \frac{L\left(\chi^{-1}, 1-k\right)}{2^{n}} \prod_{\mathfrak{q}}\left(1-\chi(\mathfrak{q}) \mathrm{Nq}^{-k}\right)
\end{aligned}
$$

The normalized constant term is $\left(\mathrm{Nb}_{\mathcal{A}}\right)^{-k} \mathrm{Nt}_{\lambda}^{-k / 2}$ multiplied by this, yielding statement (1) of the theorem.

Everything up to this point also applies when $k=1$. However, when $k=1$, the formula in [10, (3.7)] shows that there is an additional term which arises from the constant term in the $q$-expansion of

$$
\sum_{(u, v) \equiv\left(u_{0}, v_{0}\right)} \frac{1}{(u z+v)|u z+v|^{2 s}}
$$

at $s=0$; its value is the following sum at $s=0$ :

$$
\frac{(-2 \pi i)^{k n} \mathrm{~N}\left(\mathfrak{b}^{-1}\right)}{2^{n} \sqrt{d(F)}} \sum_{u \equiv u_{0}} \frac{\operatorname{sgn}(\mathrm{~N}(u))}{|\mathrm{N} u|^{2 s}} .
$$

Therefore the second term in the constant term of $\left.g_{\lambda}(z, s)\right|_{A}$ at $s=0$ is the value of the following at $s=0$ :

$$
\begin{equation*}
\sum_{u_{0}, v_{0}} \operatorname{sgn}\left(u_{0} \delta-v_{0} \gamma\right)^{q_{1}} \chi_{1}\left(\left(u_{0} \delta-v_{0} \gamma\right) \mathfrak{r}^{-1}\right) \tag{41}
\end{equation*}
$$

$$
\times \operatorname{sgn}\left(u_{0} \beta-v_{0} \alpha\right)^{q_{2}} \chi_{2}^{-1}\left(\left(u_{0} \beta-v_{0} \alpha\right) \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right) \frac{(-2 \pi i)^{k n} \mathrm{~N}\left(\mathfrak{b}^{-1}\right)}{2^{n} \sqrt{d(F)}} \sum_{u \equiv u_{0}} \frac{\operatorname{sgn}(\mathrm{~N}(u))}{|\mathrm{N} u|^{2 s}}
$$

$$
=\frac{(-2 \pi i)^{k n} \mathrm{~N}\left(\mathfrak{b}^{-1}\right)}{2^{n} \sqrt{d(F)}} \sum_{u} \frac{\operatorname{sgn}(\mathrm{~N}(u))}{|\mathrm{N} u|^{2 s}} \sum_{v_{0}} \operatorname{sgn}\left(u \delta-v_{0} \gamma\right)^{q_{1}} \chi_{1}\left(\left(u \delta-v_{0} \gamma\right) \mathfrak{r}^{-1}\right)
$$

$$
\times \operatorname{sgn}\left(u \beta-v_{0} \alpha\right)^{q_{2}} \chi_{2}^{-1}\left(\left(u \beta-v_{0} \alpha\right) \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right)
$$

Here the second sum on the right side of (41) runs through all $v_{0}$ in a set of representatives for

$$
\mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} / \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \mathfrak{a b}
$$

By the definition of $\mathfrak{g}$, we have $\gamma \in \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}$, and hence $v_{0} \in \mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \mathfrak{a}$ implies $v_{0} \gamma \in \mathfrak{r a}$. Therefore the last sum above (i.e. the expression appearing in (41) after the $\times$ symbol) can be written as a double sum

$$
\begin{align*}
& \sum_{v_{0} \in \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b} \mathcal{A}^{\mathfrak{A}}\right)^{-1} / \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \mathfrak{a}} \operatorname{sgn}\left(u \delta-v_{0} \gamma\right)^{q_{1}} \chi_{1}\left(\left(u \delta-v_{0} \gamma\right) \mathfrak{r}^{-1}\right)  \tag{42}\\
& \times\left(\sum_{\substack{v_{0}^{\prime} \in \mathfrak{r}\left(\mathfrak{d t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} / \mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A} \mathfrak{A}}\right)^{-1} \mathfrak{a b b} \\
v_{0}^{\prime} \equiv v_{0}\left(\bmod \mathfrak{r}\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}} \mathfrak{g}\right)^{-1} \mathfrak{a}\right)}} \operatorname{sgn}\left(u \beta-v_{0}^{\prime} \alpha\right)^{q_{2}} \chi_{2}^{-1}\left(\left(u \beta-v_{0}^{\prime} \alpha\right) \mathfrak{b} \mathfrak{d}_{\lambda} \mathfrak{r}^{-1}\right)\right) . \tag{43}
\end{align*}
$$

Recall that the "finite part" of the character $\chi_{2}$ is the character

$$
\chi_{2, f}:\left(\mathcal{O}_{F} / \mathfrak{b}\right)^{*} \rightarrow \mathbb{C}^{*}, \quad \chi_{2, f}(\alpha)=\operatorname{sgn}(\alpha)^{q_{2}} \chi_{2}((\alpha))
$$

We extend $\chi_{2, f}$ to a function of $\mathcal{O} / \mathfrak{b}$ by dictating $\chi_{2, f}(\alpha)=0$ if $\operatorname{gcd}(\alpha, \mathfrak{b})$ $\neq 1$. Up to multiplication by a non-zero scalar, the expression 43) in large parenthesis is the sum of $\chi_{2, f}$ over a coset of the ideal in $\mathcal{O} / \mathfrak{b}$ generated by $\mathfrak{a}_{\mathcal{A}} \mathfrak{b g}^{-1}$. Since $\chi_{2}$ is primitive of conductor $\mathfrak{b}$, it is elementary that such a sum vanishes unless $\mathfrak{a}_{\mathcal{A}} \mathfrak{b g} \mathfrak{g}^{-1}$ is divisible by $\mathfrak{b}$, i.e. unless $\mathfrak{g}=\mathcal{O}_{F}$. In other words, if $[\mathcal{A}] \notin C_{0}(\mathfrak{b}, \mathfrak{n})$ then the sum 43$)$, is 0 and if $[\mathcal{A}] \in C_{0}(\mathfrak{b}, \mathfrak{n})$ then 43$)$ equals

$$
\mathrm{Nb} \cdot \operatorname{sgn}(u \beta)^{q_{2}} \chi_{2}^{-1}\left(u \beta \mathfrak{b} \mathfrak{d t}_{\lambda} \mathfrak{r}^{-1}\right)
$$

As we now show, a similar argument implies that the sum (41) is zero unless we also have $[\mathcal{A}] \in C_{\infty}\left(\mathfrak{a}_{0}, \mathfrak{n}\right)$. Since $\mathfrak{g}=1$, the sum

$$
\begin{equation*}
\sum_{v_{0}} \operatorname{sgn}\left(u \delta-v_{0} \gamma\right)^{q_{1}} \chi_{1}\left(\left(u \delta-v_{0} \gamma\right) \mathfrak{r}^{-1}\right) \tag{44}
\end{equation*}
$$

appearing in $(42)$ is the sum of $\chi_{1, f}$ over a coset of the ideal in $\mathcal{O} / \mathfrak{a}$ generated by $\gamma\left(\mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)^{-1}=\mathfrak{c}_{\mathcal{A}}$. This vanishes unless $\mathfrak{a}_{0}=\operatorname{cond}\left(\chi_{1}\right)$ divides $\mathfrak{c}_{\mathcal{A}}$. Hence (44) vanishes unless $\mathfrak{a}_{0} \mid \mathfrak{c}_{A}$, i.e. unless $[\mathcal{A}] \in C_{\infty}\left(\mathfrak{a}_{0}, \mathfrak{n}\right)$. Furthermore, when $[\mathcal{A}] \in C_{\infty}\left(\mathfrak{a}_{0}, \mathfrak{n}\right)$ the value of (44) can be easily calculated directly. Let $\mathfrak{a}_{2}=$ $\prod_{\mathfrak{q} \in J_{\mathfrak{a}_{1}}^{c}} \mathfrak{q}$ and $\mathfrak{a}_{3}=\mathfrak{a}_{1} / \mathfrak{a}_{2}=\prod_{\mathfrak{q} \in J_{\mathfrak{a}_{1}}} \mathfrak{q}$. Then the value of 44$)$ is

$$
N \mathfrak{a}_{3} \prod_{\mathfrak{q} \mid \mathfrak{a}_{3}}\left(1-\mathrm{Nq}^{-1}\right) \operatorname{sgn}\left(u_{0} \delta\right)^{q_{1}} \chi_{1}^{*}\left(u_{0} \delta \mathfrak{r}^{-1}\right)
$$

where $\chi_{1}^{*}$ is the character $\chi_{1}$ with modulus $\mathfrak{a}_{0} \mathfrak{a}_{2}$.
Combining these calculations, we find that for $[\mathcal{A}] \in C_{0}(\mathfrak{b}, \mathfrak{n}) \cap C_{\infty}\left(\mathfrak{a}_{0}, \mathfrak{n}\right)$, the second part of the constant term of $\left.f_{\lambda}(z, s)\right|_{A}$ at $s=0$ is the value at $s=0$ of the following:

$$
\begin{gathered}
C_{\lambda} \tau\left(\chi_{2}\right) \frac{(-2 \pi i)^{n}}{2^{n} \mathrm{~N}\left(\mathfrak{r}\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)^{-1} \mathfrak{a}_{3}\right) \sqrt{d(F)}} \cdot \mathrm{Na}_{3} \prod_{\mathfrak{q} \mid \mathfrak{a}_{3}}\left(1-\mathrm{Nq}^{-1}\right) \sum_{\mathfrak{r} \in \mathrm{Cl}(F)} \mathrm{Nr} \\
\times \sum_{u_{0}} \operatorname{sgn}\left(u_{0} \delta\right)^{q_{1}} \operatorname{sgn}\left(u_{0} \beta\right)^{q_{2}} \chi_{1}^{*}\left(u_{0} \delta \mathfrak{r}^{-1}\right) \chi_{2}^{-1}\left(u_{0} \beta \mathfrak{b} \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{r}^{-1}\right) \sum_{u \equiv u_{0}} \frac{\operatorname{sgn}(\mathrm{~N} u)}{|\mathrm{N} u|^{2 s}} \\
=C_{\lambda} \tau\left(\chi_{2}\right) \frac{(-2 \pi i)^{n}}{2^{n} \mathrm{~N}\left(\left(\mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)^{-1}\right) \sqrt{d(F)}} \\
\times \prod_{\mathfrak{q} \mid \mathfrak{a}_{3}}\left(1-\mathrm{Nq}^{-1}\right) \operatorname{sgn}(\delta)^{q_{1}} \operatorname{sgn}(\beta)^{q_{2}} \chi_{1}^{*}\left(\delta \mathfrak{b}^{-1} \mathfrak{b}_{\mathcal{A}}\right) \chi_{2}^{-1}\left(\beta \mathfrak{d}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right) \\
\times \sum_{\mathfrak{r}} \sum_{u} \chi_{1}^{*} \chi_{2}^{-1}\left(u \mathfrak{b}\left(\mathfrak{b}_{\mathcal{A}}\right)^{-1} \mathfrak{r}^{-1}\right) \frac{1}{|\mathrm{~N} u|^{2 s}}
\end{gathered}
$$

The value at $s=0$ is

$$
\begin{align*}
& \frac{\tau\left(\chi_{2}\right) \mathrm{Nt}_{\lambda}^{1 / 2} \mathrm{Nb}_{\mathcal{A}}}{2^{n} \mathrm{~N} \mathfrak{b}} \operatorname{sgn}(\delta)^{q_{1}} \operatorname{sgn}(\beta)^{q_{2}} \chi_{1}^{*}\left(\delta \mathfrak{b}^{-1} \mathfrak{b}_{\mathcal{A}}\right) \chi_{2}^{-1}\left(\beta \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)  \tag{45}\\
& \times L(\chi, 0) \prod_{\mathfrak{q} \mid \mathfrak{a}_{2}}(1-\chi(\mathfrak{q})) \prod_{\mathfrak{q} \mid \mathfrak{a}_{3}}\left(1-\mathrm{Nq}^{-1}\right)
\end{align*}
$$

Since $\mathfrak{a}_{0} \mathfrak{a}_{2} \mid \mathfrak{c}_{\mathcal{A}}$, we have $\beta \gamma \in \mathfrak{a}_{0} \mathfrak{a}_{2}$ and hence $\alpha \delta \equiv 1\left(\bmod \mathfrak{a}_{0} \mathfrak{a}_{2}\right)$, whence

$$
\operatorname{sgn}(\delta)^{q_{1}} \chi_{1}^{*}\left(\delta \mathfrak{b}^{-1} \mathfrak{b}_{\mathcal{A}}\right)=\operatorname{sgn}(\alpha)^{q_{1}}\left(\chi_{1}^{*}\right)^{-1}\left(\mathfrak{a}_{\mathcal{A}} / \mathfrak{b}\right)=\operatorname{sgn}(\alpha)^{q_{1}} \chi_{1}^{-1}\left(\mathfrak{a}_{\mathcal{A}} / \mathfrak{b}\right)
$$

where the last equality follows since $\operatorname{gcd}\left(\mathfrak{a}_{\mathcal{A}}, \mathfrak{a}_{2}\right)=1$. Similarly, $\alpha \delta \in \mathfrak{b}$ implies $-\beta \gamma \equiv 1(\bmod \mathfrak{b})$, hence

$$
\operatorname{sgn}(\beta)^{q_{2}} \chi_{2}^{-1}\left(\beta \mathfrak{d} \mathfrak{t}_{\lambda} \mathfrak{b}_{\mathcal{A}}\right)=\operatorname{sgn}(-\gamma)^{q_{2}} \chi_{2}\left(\mathfrak{c}_{\mathcal{A}}\right)
$$

Therefore, noting equation (31), after scaling by the normalization factor
$\left(\mathrm{Nb}_{\mathcal{A}}\right)^{-1} \mathrm{Nt}_{\lambda}^{-1 / 2}$ for constant terms, the value in 45 is equal to

$$
P_{\mathcal{A}}\left(\chi_{2}, \chi_{1}, 1\right) \prod_{\mathfrak{q} \mid \mathfrak{a}_{2}}(1-\chi(\mathfrak{q})) \prod_{\mathfrak{q} \mid \mathfrak{a}_{3}}\left(1-\mathrm{Nq}^{-1}\right)
$$

The first term calculated above (for $k \geq 1$ ) is non-zero only when the cusp $[\mathcal{A}]$ belongs to $C_{\infty}(\mathfrak{b}, \mathfrak{n}) \cap C_{0}(\mathfrak{a}, \mathfrak{n})$. The second term is non-zero only when $[\mathcal{A}]$ belongs to $C_{0}(\mathfrak{b}, \mathfrak{n}) \cap C_{\infty}\left(\mathfrak{a}_{0}, \mathfrak{n}\right)$. This finishes the proof of 4.5 .
4.2. Constant terms for raised level and imprimitive characters. In our arithmetic application [3], we require the constant terms of the levelraised Eisenstein series $\left.E_{k}(\chi, \psi)\right|_{\mathfrak{m}}$ for auxiliary square-free ideals $\mathfrak{m}$, with $\chi$ and $\psi$ possibly imprimitive. This level raising is related to the $T$-smoothing operation of Deligne-Ribet [5].

The following notation will be in effect throughout this section. Let $\chi$ and $\psi$ be characters of modulus $\mathfrak{a}$ and $\mathfrak{b}$ and signatures $q_{1}$ and $q_{2}$, respectively. Let $k$ be a positive integer such that $q_{1}+q_{2} \equiv(k, \ldots, k)(\bmod 2)$. We denote the conductors of $\chi$ and $\psi$ by $\mathfrak{a}_{0}$ and $\mathfrak{b}_{0}$, respectively, and put $\mathfrak{a}_{1}=\mathfrak{a} / \mathfrak{a}_{0}$ and $\mathfrak{b}_{1}=\mathfrak{b} / \mathfrak{b}_{0}$. Assume $\operatorname{gcd}\left(\mathfrak{b}_{1}, \mathfrak{a}\right)=1$. Let $\mathfrak{n}=\mathfrak{a b l}$ for a square-free integral ideal $\mathfrak{l}$ with $\operatorname{gcd}(\mathfrak{a b}, \mathfrak{l})=1$. We assume that $\mathfrak{a}_{1}$ is square-free and coprime to $\mathfrak{a}_{0}$, and similarly for $\left(\mathfrak{b}_{1}, \mathfrak{b}_{0}\right)$.

Let $\mathcal{A}=(A, \lambda)$ with

$$
A=\left(\begin{array}{ll}
\alpha & * \\
\gamma & *
\end{array}\right) \in \mathbf{G L}_{2}^{+}(F), \quad \lambda \in \mathrm{Cl}^{+}(F)
$$

TheOrem 4.8. Let $\mathfrak{m}$ be a divisor of $\mathfrak{l}$.

- If $k \geq 2$, then the normalized constant term of $\left.E_{k}(\chi, \psi)\right|_{\mathfrak{m}}$ at $\mathcal{A}$ is

$$
\begin{align*}
\delta_{\infty, \mathcal{A}}\left(\mathfrak{b}_{0}\right) P_{\mathcal{A}}(\chi, \psi & \left., k, J_{\mathfrak{b}_{1}}^{c}, J_{\mathfrak{a}_{1}}\right)  \tag{46}\\
& \times \prod_{\mathfrak{q} \in J_{\mathfrak{b}_{1}}}\left(1-\mathrm{Nq}^{-1}\right) \prod_{\mathfrak{q} \in J_{\mathfrak{m}}}\left(\psi(\mathfrak{q}) \mathrm{Nq}^{k}\right)^{-1} \prod_{\mathfrak{q} \in J_{\mathfrak{m}}^{c}} \chi^{-1}(\mathfrak{q}) .
\end{align*}
$$

- If $k=1$, further assume that $\operatorname{gcd}(\mathfrak{a}, \mathfrak{b})=1$. Then the normalized constant term of $\left.E_{1}(\chi, \psi)\right|_{\mathfrak{m}}$ at $\mathcal{A}$ is

$$
\begin{align*}
\delta_{0, \mathcal{A}}(\mathfrak{a}) \delta_{\infty, \mathcal{A}}\left(\mathfrak{b}_{0}\right) P_{\mathcal{A}}(\chi, & \left.\psi, 1, J_{\mathfrak{b}_{1}}^{c}, J_{\mathfrak{a}_{1}}\right)  \tag{47}\\
& \times \prod_{\mathfrak{q} \in J_{\mathfrak{b}_{1}}}\left(1-\mathrm{Nq}^{-1}\right) \prod_{\mathfrak{q} \in J_{\mathfrak{m}}}(\psi(\mathfrak{q}) \mathrm{Nq})^{-1} \prod_{\mathfrak{q} \in J_{\mathfrak{w}}^{c}} \chi^{-1}(\mathfrak{q}) \\
+\delta_{\infty, \mathcal{A}}\left(\mathfrak{a}_{0}\right) \delta_{0, \mathcal{A}}(\mathfrak{b}) P_{\mathcal{A}}( & \left.\psi, \chi, 1, J_{\mathfrak{a}_{1}}^{c}, J_{\mathfrak{b}_{1}}\right) \\
& \times \prod_{\mathfrak{q} \in J_{\mathfrak{a}_{1}}}\left(1-\mathrm{Nq}^{-1}\right) \prod_{\mathfrak{q} \in J_{\mathfrak{m}}}(\chi(\mathfrak{q}) \mathrm{Nq})^{-1} \prod_{\mathfrak{q} \in J_{\mathfrak{m}}^{c}} \psi^{-1}(\mathfrak{q})
\end{align*}
$$

Proof. We give the proof for $k \geq 2$. The argument for $k=1$ is identical and left to the reader.

First we assume that $\mathfrak{b}=\mathfrak{b}_{0}$, i.e. $\mathfrak{b}_{1}=1$, and calculate the constant term of $\left.E_{k}(\chi, \psi)\right|_{\mathfrak{m}}$. Let $\mathfrak{m}=\mathfrak{q}_{1} \cdots \mathfrak{q}_{j}$ and use induction on $j$. The base case $j=0$ follows directly from (32). Write $\mu$ for the Möbius function defined by $\mu(\mathfrak{t})=(-1)^{k}$ where $k$ is the number of prime factors of $\mathfrak{t}$ if $\mathfrak{t}$ is square-free, and $\mu(\mathfrak{t})=0$ otherwise. For $j>0$ we use the expression

$$
\begin{equation*}
E_{k}\left(\chi_{\mathfrak{m}}, \psi\right)=\left.\sum_{\mathfrak{t} \mid \mathfrak{m}} \mu(\mathfrak{t}) \chi(\mathfrak{t}) E_{k}(\chi, \psi)\right|_{\mathfrak{t}} \tag{48}
\end{equation*}
$$

Here $\chi_{\mathfrak{m}}$ means the character $\chi$ viewed with modulus $\mathfrak{a m}$. If $J_{\mathfrak{m}}=\left\{\mathfrak{q}_{1}, \ldots, \mathfrak{q}_{j}\right\}$, then $\delta_{0, \mathcal{A}}(\mathfrak{a m})=\delta_{0, \mathcal{A}}(\mathfrak{a})$, so by 32 the normalized constant term of $E_{k}\left(\chi_{\mathfrak{m}}, \psi\right)$ at $\mathcal{A}$ is

$$
\begin{equation*}
\delta_{0, \mathcal{A}}(\mathfrak{a}) \delta_{\infty, \mathcal{A}}\left(\mathfrak{b}_{0}\right) P_{\mathcal{A}}\left(\chi, \psi, k, \emptyset, J_{\mathfrak{a}_{1}}\right) \prod_{\mathfrak{q} \mid \mathfrak{m}}\left(1-\chi \psi^{-1}\left(\mathfrak{q}_{i}\right) \mathrm{Nq}_{i}^{-k}\right) \tag{49}
\end{equation*}
$$

The induction hypothesis gives the normalized constant term at $\mathcal{A}$ of each term on the right side of (48) except for $\left.E_{k}(\chi, \psi)\right|_{\mathfrak{m}}$. Therefore one can use (48) and 49 to solve for the normalized constant term of $\left.E_{k}(\chi, \psi)\right|_{\mathfrak{m}}$ at $\mathcal{A}$. One obtains

$$
\delta_{0, \mathcal{A}}(\mathfrak{a}) \delta_{\infty, \mathcal{A}}\left(\mathfrak{b}_{0}\right) P_{\mathcal{A}}\left(\chi, \psi, k, \emptyset, J_{\mathfrak{a}_{1}}\right) \prod_{\mathfrak{q} \mid \mathfrak{m}}\left(\chi \psi^{-1}\left(\mathfrak{q}_{i}\right) \mathrm{Nq}_{i}^{k}\right)^{-1}
$$

as desired. Now suppose $J_{\mathfrak{m}} \neq\left\{\mathfrak{q}_{1}, \ldots, \mathfrak{q}_{j}\right\}$, which is equivalent to $[\mathcal{A}] \notin$ $C_{0}(\mathfrak{m}, \mathfrak{n})$. This in turn implies that $[\mathcal{A}] \in C_{\infty}\left(\mathfrak{q}_{i}, \mathfrak{n}\right)$ for some $i$, whence $\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}_{0}, \mathfrak{a m}\right) \neq 1$ and $\chi_{\mathfrak{m}}\left(\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}_{0}\right)=0$. Therefore $(32)$ implies that the constant term of $E_{k}\left(\chi_{\mathfrak{m}}, \psi\right)$ at $\mathcal{A}$ is 0 . Without loss of generality, assume that $\mathfrak{q}_{j} \notin J_{\mathfrak{m}}$. For every subset $I \subset\left\{\mathfrak{q}_{1}, \ldots, \mathfrak{q}_{j-1}\right\}$, put $\mathfrak{q}_{I}=\prod_{\mathfrak{q}_{i} \in I} \mathfrak{q}_{i}$. If $I \neq\left\{\mathfrak{q}_{1}, \ldots, \mathfrak{q}_{j-1}\right\}$, then we can apply the induction hypothesis to the forms $\left.E_{k}(\chi, \psi)\right|_{\mathfrak{t}}$ for both $\mathfrak{t}=\mathfrak{q}_{I}$ and $\mathfrak{t}=\mathfrak{q}_{I} \mathfrak{q}_{j}$ on the right side of (48) to see that the contributions made by their constant terms at $\mathcal{A}$ cancel. It follows that the constant term of $\left.E_{k}(\chi, \psi)\right|_{\mathfrak{m}}$ at $\mathcal{A}$ equals $\chi^{-1}\left(\mathfrak{q}_{j}\right)$ times that of $\left.E_{k}(\chi, \psi)\right|_{\mathfrak{m} / \mathfrak{q}_{j}}$, and we are done by the induction hypothesis.

Next we relax the condition that $\mathfrak{b}_{1}=1$. We use the expression

$$
\begin{equation*}
\left.E_{k}(\chi, \psi)\right|_{\mathfrak{m}}=\left.\sum_{\mathfrak{t} \mid \mathfrak{b}_{1}} \mu(\mathfrak{t}) \psi(\mathfrak{t})(\mathrm{Nt})^{k-1} E_{k}\left(\chi, \psi^{0}\right)\right|_{\mathfrak{t} \mathfrak{m}} \tag{50}
\end{equation*}
$$

where $\psi^{0}$ is the primitive character associated with $\psi$. The case already completed for $\psi$ primitive gives the constant terms of the forms on the right
of (50). The result then follows from the formula

$$
\begin{aligned}
& \sum_{\mathfrak{t} \mid \mathfrak{b}_{1}} \mu(\mathfrak{t}) \psi(\mathfrak{t})(\mathrm{Nt})^{k-1} \prod_{\mathfrak{q} \in J_{\mathfrak{t}}}\left(\psi(\mathfrak{q}) \mathrm{Nq}^{k}\right)^{-1} \prod_{\mathfrak{q} \in J_{\mathfrak{t}}^{c}} \chi^{-1}(\mathfrak{q}) \\
&=\prod_{\mathfrak{q} \in J_{\mathfrak{b}_{1}}}\left(1-\mathrm{Nq}^{-1}\right) \prod_{\mathfrak{q} \in J_{\mathfrak{b}_{1}}^{c}}\left(1-\chi^{-1} \psi(\mathfrak{q}) \mathrm{Nq}^{k-1}\right)
\end{aligned}
$$

5. Ordinary forms. Let $\mathfrak{p}$ be a prime ideal of $\mathcal{O}_{F}$ dividing a prime number $p$. Following Hida, we define the ordinary operator

$$
e_{\mathfrak{p}}^{\mathrm{ord}}=\lim _{n \rightarrow \infty} U_{\mathfrak{p}}^{n!}
$$

Let $\mathfrak{P}=\operatorname{gcd}\left(p^{\infty}, \mathfrak{n}\right)$ be the $p$-part of $\mathfrak{n}$. We define

$$
e_{\mathfrak{P}}^{\text {ord }}=\prod_{\mathfrak{p} \mid \mathfrak{P}} e_{\mathfrak{p}}^{\text {ord }}
$$

Let $E$ be a finite extension of $\mathbb{Q}_{p}$. The space of $\mathfrak{P}$-ordinary forms is defined by

$$
M_{k}(\mathfrak{n}, E)^{\mathfrak{P}-\text { ord }}=e_{\mathfrak{P}}^{\operatorname{ord}} M_{k}(\mathfrak{n}, E)
$$

This is the largest subspace on which the operator $U_{\mathfrak{p}}$ acts invertibly for each $\mathfrak{p} \mid \mathfrak{P}$.

TheOrem 5.1. A form $f \in M_{k}(\mathfrak{n}, E)^{\mathfrak{P} \text {-ord }}$ is cuspidal if and only if its constant terms at all cusps in $C_{\infty}(\mathfrak{P}, \mathfrak{n})$ are zero. If $f \in M_{k}(\mathfrak{n}, E)$ has constant terms zero at all cusps in $C_{\infty}(\mathfrak{P}, \mathfrak{n})$, then $e_{\mathfrak{P}}^{\text {ord }}(f)$ is cuspidal.

We provide two proofs of Theorem 5.1. The first proof is longer, but its method could have other applications, so we include full details. We begin with the following elementary lemma from linear algebra.

Lemma 5.2. Let $V$ be a finite-dimensional vector space over a field and let $B=\left\{v_{1}, \ldots, v_{n}\right\}$ be a basis. Let $S$ be a possibly infinite set of commuting endomorphisms of $V$ satisfying the following properties:

- After re-ordering, the matrix for each $T \in S$ with respect to the basis $B$ is in Jordan canonical form.
- Every Jordan block of size greater than 1 has associated eigenvalue 0.

Let $B^{\prime} \subset B$ be the set of basis vectors that are actual (non-generalized) eigenvectors for every $T \in S$. Suppose the elements of $B^{\prime}$ are distinguished by their $S$-eigenvalues, i.e. for $v_{i} \neq v_{j}$ in $B^{\prime}$ there exists $T \in S$ such that the $T$-eigenvalues of $v_{i}$ and $v_{j}$ are distinct. Finally, let $W \subset V$ be a subspace that is preserved by each $T$. Then $W$ is non-zero if and only if it contains some $v_{i} \in B^{\prime}$.

Proof. Suppose $v=\sum a_{i} v_{i} \in W$ with the $a_{i}$ not all zero. We first show that we can find another non-zero $v^{\prime} \in W$ such that its expression as a
linear combination of elements in $B$ only contains elements of $B^{\prime}$. For this, suppose that $v_{i} \in B \backslash B^{\prime}$ occurs in $v$ with a non-zero coefficient $a_{i}$. Let $T \in S$ be such that $v_{i}$ is not an eigenvector for $T$. Then there is a unique $n \geq 1$ such that $T^{n}\left(v_{i}\right) \in B$ is an eigenvector for $T$. We replace $v$ by $T^{n}(v)$. This is another element of $W$; its expression as a linear combination of the $v_{i}$ has at most as many elements of $B \backslash B^{\prime}$ as $v$. And the term $a_{i} v_{i}$ has been replaced by $a_{i} T^{n}\left(v_{i}\right)$-this uses the fact that the $T$-eigenvalue of $T^{n}\left(v_{i}\right)$ is 0 . Note in particular that since $T^{n}\left(v_{i}\right) \in B$ occurs with a non-zero coefficient, $T^{n}(v) \neq 0$. If $T^{n}\left(v_{i}\right) \in B^{\prime}$, we have reduced the number of elements of $B \backslash B^{\prime}$ in our linear combination. If not, there is some other $T^{\prime} \in S$ that we can apply a certain number of times, say $m$, to replace $T^{n}\left(v_{i}\right)$ by its associated $T^{\prime}$-eigenvector. Continuing in this way, we get a sequence of non-zero vectors $v \rightarrow T^{n}(v) \rightarrow\left(T^{\prime}\right)^{m} T^{n}(v) \rightarrow \cdots$ and a corresponding sequence of terms occurring in the expression of these vectors in terms of $B$ :

$$
a_{i} v_{i} \rightarrow a_{i} T^{n}\left(v_{i}\right) \rightarrow a_{i}\left(T^{\prime}\right)^{m} T^{n}(v) \rightarrow \cdots
$$

Since this latter sequence clearly cannot contain a cycle, and $B$ is finite, it must terminate. This occurs when the corresponding element of $B$ actually lies in $B^{\prime}$. We have therefore created a new non-zero element of $W$ whose expression in the basis $B$ contains fewer elements of $B \backslash B^{\prime}$. Continuing this procedure yields a non-zero element of $W$ that is in the span of $B^{\prime}$.

Now let $v=\sum a_{i} v_{i} \in W$ with $v_{i} \in B^{\prime}$ be such an element. If more than one $v_{i}$ occurs in this linear combination with non-zero coefficient, say $v_{1}$ and $v_{2}$, then by assumption we can find $T \in S$ such that the associated eigenvalues $\lambda_{1}(T)$ and $\lambda_{2}(T)$ are distinct. We can replace $v$ by $T(v)-\lambda_{1}(T) v$. This annihilates the $v_{1}$ term, but is non-zero because it does not annihilate the $v_{2}$ term. Furthermore it has fewer non-zero coefficients than $v$. Continuing in this way, we can repeatedly decrease the number of terms in the expression of $v$ until we find that there is some $v_{i} \in B^{\prime} \cap W$.

Proof 1 of Theorem 5.1. The second statement of the theorem follows from the first since $e_{\mathfrak{P}}^{\text {ord }}(f)$ preserves the space $E\left[C_{\infty}(\mathfrak{P}, \mathfrak{n})\right]$. To prove the first statement, let $f \in M_{k}(\mathfrak{n}, E)^{\mathfrak{P} \text { - ord }}$ be a form whose constant terms at all cusps in $C_{\infty}(\mathfrak{P}, \mathfrak{n})$ are zero. Then $f$ is a sum of a cusp form and a linear combination of Eisenstein series. The cusp form does not affect any constant terms; we can therefore assume that $f$ is a linear combination of Eisenstein series, and we must show that $f=0$. The Eisenstein subspace has the following convenient basis, for which each of the Hecke operators is in Jordan canonical form:

$$
\begin{equation*}
B=\left\{\left.E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}}\right)\right|_{\mathfrak{c}}\right\} \tag{51}
\end{equation*}
$$

where

- $\eta, \psi$ are primitive characters of conductor $\mathfrak{a}, \mathfrak{b}$ respectively,
- $\mathfrak{r}, \mathfrak{s}$ are square-free products of primes such that $\operatorname{gcd}(\mathfrak{a}, \mathfrak{r})=\operatorname{gcd}(\mathfrak{b}, \mathfrak{s})=1$,
- $\mathfrak{a b r s}$ is divisible by all primes dividing $\mathfrak{n}$,
- $\mathfrak{c}$ is only divisible by primes dividing $\operatorname{gcd}(\mathfrak{a r}, \mathfrak{b s})$.
- $\mathfrak{a b r s c}$ divides $\mathfrak{n}$.

Since this is a lot of notation, it behooves us to demonstrate this with an example. Suppose that $\eta, \psi$ are primitive of conductor $\mathfrak{a}, \mathfrak{b}$ with associated Eisenstein series $E_{k}(\eta, \psi) \in M_{k}(\mathfrak{a b})$. Let $\mathfrak{p}$ be a prime not dividing $\mathfrak{a b}$. For $n \geq 1$, the generalized eigenspace of $M_{k}\left(\mathfrak{a b p}^{n}\right)$ corresponding to $E_{k}(\eta, \psi)$-i.e. the subspace on which all the Hecke operators away from $\mathfrak{p}$ act via the eigenvalues of $E_{k}(\eta, \psi)$-has two or three Jordan blocks for the action of $U_{\mathfrak{p}}:(1)$ the form $E_{k}\left(\eta_{\mathfrak{p}}, \psi\right)$ with $U_{\mathfrak{p}}$-eigenvalue $\psi(\mathfrak{p}) \mathrm{Np}^{k-1}$, (2) the form $E_{k}\left(\eta, \psi_{\mathfrak{p}}\right)$ with $U_{\mathfrak{p}}$-eigenvalue $\eta(\mathfrak{p})$, and (3) if $n \geq 2$, a Jordan block with $U_{\mathfrak{p}}$-eigenvalue 0 and basis $\left.E_{k}\left(\eta_{\mathfrak{p}}, \psi_{\mathfrak{p}}\right)\right|_{\mathfrak{p}^{i}}$ as $i=0, \ldots, n-2$. Here $U_{\mathfrak{p}}\left(\left.E_{k}\left(\eta_{\mathfrak{p}}, \psi_{\mathfrak{p}}\right)\right|_{\mathfrak{p}^{i}}\right)=\left.E_{k}\left(\eta_{\mathfrak{p}}, \psi_{\mathfrak{p}}\right)\right|_{\mathfrak{p}^{i-1}}$ for $i \geq 1$, and $U_{\mathfrak{p}}\left(E_{k}\left(\eta_{\mathfrak{p}}, \psi_{\mathfrak{p}}\right)\right)=0$. The basis (51) is the generalization of this case to the general setting.

Now, the space $E_{k}(\mathfrak{n}, E)^{\mathfrak{P} \text { - ord }}$ is the subspace of $E_{k}(\mathfrak{n}, E)$ generated by the subset $B_{\mathfrak{P}} \subset B$ consisting of the $\left.E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}}\right)\right|_{\mathfrak{c}}$ such that $\mathfrak{a r}$ is coprime to $\mathfrak{P}$. We apply Lemma 5.2 where the set $S$ of endomorphisms is the set of Hecke operators indexed by the primes not dividing $\mathfrak{P}$ : the $T_{\mathfrak{q}}$ for $\mathfrak{q} \nmid \mathfrak{n}$ and the $U_{\mathfrak{q}}$ for $\mathfrak{q} \mid \mathfrak{n} / \mathfrak{P}$. The subspace $W \subset E_{k}(\mathfrak{n}, E)^{\mathfrak{P} \text { - ord }}$ is taken to be the subspace of elements whose constant terms at all cusps in $C_{\infty}(\mathfrak{P}, \mathfrak{n})$ vanish. This subspace is fixed by the Hecke operators away from $\mathfrak{P}$. We need to prove that $W=\{0\}$, and the lemma implies that it suffices to show that no eigenvector in $B_{\mathfrak{P}}$ lies in $W$. The subset $B_{\mathfrak{P}}^{\prime} \subset B_{\mathfrak{P}}$ of eigenvectors is the set of $\left.E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}}\right)\right|_{\mathfrak{c}}$ such that $\mathfrak{a r}$ is coprime to $\mathfrak{P}$ and $\mathfrak{c}=1$. It remains to prove that for such a form $E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}}\right)$, there exists a cusp $[\mathcal{A}] \in C_{\infty}(\mathfrak{P}, \mathfrak{n})$ such that the constant term at $\mathcal{A}$ is non-zero.

For this, we first note that it suffices to show this at the minimal level at which $E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}}\right)$ appears, namely $\mathfrak{n}^{\prime}=\mathfrak{a b r s}$. Indeed, if we let $\mathfrak{P}^{\prime}$ be the $p$-part of $\mathfrak{n}^{\prime}$, then the canonical map $C_{\infty}(\mathfrak{P}, \mathfrak{n}) \rightarrow C_{\infty}\left(\mathfrak{P}^{\prime}, \mathfrak{n}^{\prime}\right)$ is surjective. Therefore we assume that $\mathfrak{n}=\mathfrak{a b v s}$.

Write $\mathfrak{s}=\mathfrak{s}_{1} \mathfrak{s}_{2}$, where $\mathfrak{s}_{1}=\operatorname{gcd}(\mathfrak{s}, \mathfrak{a r})$. Note that $\mathfrak{P} \mid \mathfrak{b s}_{2}$. Let $\mathcal{A}$ be a cusp in $C_{\infty}(\mathfrak{b s}, \mathfrak{n}) \subset C_{\infty}(\mathfrak{P}, \mathfrak{n})$ such that $\mathfrak{c}_{\mathcal{A}} / \mathfrak{b}_{1}$ is coprime to $\mathfrak{a r}$. We use the expression

$$
E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}}\right)=\left.\sum_{\mathfrak{m} \mid \mathfrak{s}_{1}} \mu(\mathfrak{m}) \psi(\mathfrak{m}) \mathrm{Nm}^{k-1} E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}_{2}}\right)\right|_{\mathfrak{m}}
$$

to show that the constant term of $E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}}\right)$ at $\mathcal{A}$ is non-zero. By Lemma 3.13, the constant term of $\left.E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}_{2}}\right)\right|_{\mathfrak{m}}$ at $\mathcal{A}$ equals the constant term of $E_{k}\left(\eta_{\mathfrak{r}}, \psi_{\mathfrak{s}_{2}}\right)$ at some other $\mathcal{A}^{\prime}$, where $\left[\mathcal{A}^{\prime}\right] \in C_{\infty}(\mathfrak{b s} / \mathfrak{m}, \mathfrak{n} / \mathfrak{m})$. If $\mathfrak{m} \neq \mathfrak{s}_{1}$, then by definition
$\mathfrak{b s} / \mathfrak{m}$ is not coprime to $\mathfrak{a r}$. Thus $P_{\mathcal{A}^{\prime}}(\eta, \psi, k)=0$ because of the $\eta\left(\mathfrak{c}_{\mathcal{A}^{\prime}}\right)$ factor, and hence by Theorem 4.8 the constant term is 0 . On the other hand, if $\mathfrak{m}=\mathfrak{s}_{1}$ then Theorem 4.8 shows that this constant term is non-zero as long as $\mathfrak{c}_{\mathcal{A}^{\prime}} / \mathfrak{b}$ is coprime to $\mathfrak{a r}$. This holds because $\mathfrak{c}_{\mathcal{A}} / \mathfrak{b} \mathfrak{s}_{1}$ is coprime with $\mathfrak{a r}$. The result follows.

Proof 2 of Theorem 5.1. Our second proof of the first statement is a direct computation using the action of the Hecke operator $U_{\mathfrak{p}}$ for each $\mathfrak{p} \mid \mathfrak{P}$. First we note that since we are on the ordinary subspace, each $U_{\mathfrak{p}}$ acts semisimply (see [6, p. 382]). Furthermore the operator $U_{\mathfrak{p}}$ preserves $C_{\infty}(\mathfrak{P}, \mathfrak{n})$. Therefore it suffices to consider the case where $f$ is a $U_{\mathfrak{p}}$-eigenvector for each $\mathfrak{p} \mid \mathfrak{P}$.

Let us recall the explicit definition of the operator $U_{\mathfrak{p}}$. For each $\mu \in$ $\mathrm{Cl}^{+}(F)$, let $\lambda \in \mathrm{Cl}^{+}(F)$ denote the class of $\mu \mathfrak{p}^{-1}$. Write $\mathfrak{t}_{\lambda} \mathfrak{t}_{\mu}^{-1} \mathfrak{p}=(x)$ where $x$ is a totally positive element of $F^{*}$. Given $\beta \in F$ define $m_{\beta}=\left(\begin{array}{ll}1 & \beta \\ 0 & x\end{array}\right)$. Then $\left.f\right|_{U_{\mathfrak{p}}}=\left(g_{\mu}\right)_{\mu \in \mathrm{Cl}^{+}(F)}$ where

$$
\begin{equation*}
g_{\mu}=\left.N \mathfrak{p}^{(k-2) / 2} \sum_{\beta \in \mathfrak{t}_{\mu}^{-1} \mathfrak{d}^{-1} / \mathfrak{t}_{\mu}^{-1} \mathfrak{d}^{-1} \mathfrak{p}}\left(f_{\lambda}\right)\right|_{m_{\beta}} \tag{52}
\end{equation*}
$$

Let $\mathcal{A}=(A, \mu)$ represent a cusp. If $\mathfrak{p}^{r} \mid \mathfrak{c}_{\mathcal{A}}$ with $r>0$, then one readily checks that $\mathfrak{p}^{r+1} \mid \mathfrak{c}_{\mathcal{A}^{\prime}}$, where $\mathcal{A}^{\prime}=\left(m_{\beta} A, \lambda\right)$ is an associated cusp appearing in (52). Therefore, since $f$ is a $U_{\mathfrak{p}}$-eigenform with non-zero eigenvalue for each $\mathfrak{p} \mid \mathfrak{P}$, and its constant terms vanish on $C_{\infty}(\mathfrak{P}, \mathfrak{n})$, we see by applying $U_{\mathfrak{p}}$ repeatedly that the constant terms of $f$ vanish on $C_{\infty}\left(\mathfrak{P}_{0}, \mathfrak{n}\right)$, where $\mathfrak{P}_{0}$ is the product of the distinct primes dividing $\mathfrak{P}$.

Next, to show that $f$ has vanishing constant terms at all cusps in cusps( $\mathfrak{n}$ ) $=C_{\infty}(1, \mathfrak{n})$, we show that we can remove the primes in $\mathfrak{P}_{0}$ one-by-one. Therefore, let $\mathfrak{P}_{1} \mid \mathfrak{P}_{0}$, and let $\mathfrak{p} \mid \mathfrak{P}_{1}$. We will show that the cuspidality of $f$ on $C_{\infty}\left(\mathfrak{P}_{1}, \mathfrak{n}\right)$ implies its cuspidality on $C_{\infty}\left(\mathfrak{P}_{1} / \mathfrak{p}, \mathfrak{n}\right)$. Sequentially removing all the primes $\mathfrak{p} \mid \mathfrak{P}_{0}$ in this fashion will then give the desired result.

For this, we use the expression (52) once again. We also introduce the notation $f(\mathcal{A})$ to mean the normalized constant term of $f$ at the cusp $\mathcal{A}$. If $\mathcal{A} \in C_{\infty}\left(\mathfrak{P}_{1} / \mathfrak{p}, \mathfrak{n}\right)$ but $\mathcal{A} \notin C_{\infty}\left(\mathfrak{P}_{1}, \mathfrak{n}\right)$, then one can check directly from the definitions that there is a unique $\beta \in \mathfrak{t}_{\mu}^{-1} \mathfrak{d}^{-1} / \mathfrak{t}_{\mu}^{-1} \mathfrak{d}^{-1} \mathfrak{p}$ such that the associated cusp $\mathcal{A}^{\prime}=\left(m_{\beta} A, \lambda\right)$ does not lie in $C_{\infty}\left(\mathfrak{P}_{1}, \mathfrak{n}\right)$ either; for all the other $\beta$, the associated cusp does lie in $C_{\infty}\left(\mathfrak{P}_{1}, \mathfrak{n}\right)$. The cuspidality of $f$ on $C_{\infty}\left(\mathfrak{P}_{1}, \mathfrak{n}\right)$ therefore implies that

$$
a_{\mathfrak{p}} f(\mathcal{A})=\left.f\right|_{U_{\mathfrak{p}}}(\mathcal{A})=\mathrm{Np}^{k-1} f\left(\mathcal{A}^{\prime}\right)
$$

where $a_{\mathfrak{p}}$ denotes the $U_{\mathfrak{p}}$-eigenvalue of $f$. Note that the constant $\mathrm{N} \mathfrak{p}^{k-1}$ arises from tracing through our normalization factors on constant terms. Now, the set $C_{\infty}\left(\mathfrak{P}_{1} / \mathfrak{p}, \mathfrak{n}\right) \backslash C_{\infty}\left(\mathfrak{P}_{1}, \mathfrak{n}\right)$ is finite, so continually repeating this process,
we see that the sequence

$$
\begin{equation*}
\mathcal{A} \rightarrow \mathcal{A}^{\prime} \rightarrow \cdots \tag{53}
\end{equation*}
$$

must eventually arrive at a repetition. At this point we obtain an equation of the form

$$
a_{\mathfrak{p}}^{r} f\left(\mathcal{A}^{\prime \prime}\right)=\mathrm{Np}^{(k-1) r} f\left(\mathcal{A}^{\prime \prime}\right)
$$

for some positive integer $r$ and some cusp $\mathcal{A}^{\prime \prime}$. As the Hecke eigenvalue $a_{\mathfrak{p}}$ is a $p$-adic unit and $k>1$, we have $a_{\mathfrak{p}}^{r} \neq \mathrm{Np}^{(k-1) r}$ for any positive integer $r$. We obtain $f\left(\mathcal{A}^{\prime \prime}\right)=0$, and hence the same is true for every other cusp appearing in the sequence 53 ; in particular $f(\mathcal{A})=0$ as desired.

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