# Eisenstein Series and the Trace Formula <br> for GL(2) over a Function Field 

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#### Abstract

We write out and prove the trace formula for a convolution operator on the space of cusp forms on GL(2) over the function field $F$ of a smooth projective absolutely irreducible curve over a finite field. The proof - which follows Drinfeld - is complete and all terms in the formula are explicitly computed. The structure of the homogeneous space $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ is studied in section 2 by means of locally free sheaves of $\mathcal{O}_{X}$-modules. Section 3 deals with the regularization and computation of the geometric terms, over conjugacy classes. Section 4 develops the theory of intertwining operators and Eisenstein Series, and the trace formula is proven in section 5 .

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## 1 Introduction and statement of the Trace Formula

### 1.1 Introduction

The (non-invariant) trace formula for GL(2) over a number field was stated and its proof sketched in chapter 15 of the influential book of Jacquet and Langlands [JL70] of 1970. It was used there for comparison of automorphic representations of the multiplicative group of a quaternion algebra, with automorphic representations of GL(2).
Drinfeld used the trace formula for $\mathrm{GL}(2)$ over a function field $F$ to prove Langlands' conjecture for $\mathrm{GL}(2, F)$, and to count in [D81] the number of two
dimensional irreducible representations of the fundamental group of a smooth projective geometrically irreducible curve $X$ over a finite field. To check the statement of the trace formula of [JL70] in the function field case, Drinfeld gave a detailed (but unpublished) proof. It differs from the one sketched in [JL70]. It is this proof of Drinfeld which is given in this paper.
The main reason why this proof is still interesting is the elementary and unconventional treatment of Eisenstein series (see subsections 4.7-4.8 below), and the computation of traces in the spirit of Tate [T68], see subsection 5.2. In both cases it is based on a "baby model" (see Proposition 4.31, Corollary 4.32, Lemma 5.11), which cries out for generalization.
Let us describe the contents of this article.
The trace formula itself is stated in subsection 1.2 with a few comments. More comments, including informal ones, are given in section 3.
Section 2 contains a dictionary between the language of adèles and the language of vector bundles on the smooth projective curve $X$ corresponding to $F$. In particular, the set of rank $n$ vector bundles on $X$ is identified with $\mathrm{GL}(n, F) \backslash \mathrm{GL}(n, \mathbb{A}) / \mathrm{GL}\left(n, O_{\mathbb{A}}\right)$, where $O_{\mathbb{A}} \subset \mathbb{A}$ is the ring of integral adèles. This dictionary goes back to A. Weil [W38], although in an older language. It underlies Drinfeld's Geometric Langlands program [BD].
The terms which appear in the geometric part of the trace formula - orbital integrals and weighted orbital integrals - are estimated and regularized in section 3.

In section 4 intertwining operators, Eisenstein series, and $L$-functions are introduced. The rationality of the intertwining operator $M\left(\mu_{1}, \mu_{2}, t\right)$ and the functional equation $M^{2}=1$ are first proven using local computations: normalization of the intertwining operators by $L$-functions and $\varepsilon$-factors, and the functional equation of the $L$-functions.
In subsections 4.7-4.8 these facts are proven using an alternative, global approach. The ideas might go back to Selberg. But technically the exposition is quite different and more elementary: in the case of function fields the analytic problems disappear.
The trace formula is proven in section 5. The logarithmic derivative of the intertwining operator appears as a result of a computation of the trace of some operator in a power series space, see Lemma 5.11. This computation is probably related to Tate's article [T68].
Here are some questions.

1. Could the methods of subsections 4.7-4.8 and section 5 be extended to prove the functional equation for Eisenstein series, and the trace formula, for an arbitrary reductive group over a function field?
2. Is there a modification of the technique from subsections 4.7-4.8 that would work in the case of number fields, e.g., for $\mathrm{GL}(2, \mathbb{Q})$ ? One could try to replace the space of formal power series used in subsections 4.7-4.8 by some space of holomorphic functions.
3. What is the precise relationship between Lemma 5.11 and Tate's [T68]?
4. What is the relationship between the approach to Eisenstein series of subsections 4.7-4.8, and the classical approaches: that of Selberg-Langlands-Arthur, and that of scattering theory (see [FP72] or [LP76])?
This author's initial motivation to write out Drinfeld's expression and proof of the trace formula for GL(2) over a function field stems from his search for higher rank analogues of Drinfeld's formula [D81]. This led us to count with Deligne [DF13] the number of rank $n(\geq 2)$ local systems with principal unipotent local monodromy at least at two places. There we use the trace formula in the compact quotient case, and the transfer of automorphic representations from a compact form to GL $(n)$. This explains the condition: "at least at two places". The case of [D81] is rank $n=2$, no monodromy. To complete the study of [D81] and of [DF13] in rank two one has to consider the case of principal unipotent local monodromy at a single place. This is done in $[F]$, using the explicit computations of the trace formula for GL(2) over a function field of the present work. This was our initial motivation to write out this formula. Drinfeld's proof in the case of rank two, no ramification, is also given in [F].
Of course there are numerous expositions of the trace formula of [JL70], e.g. [GJ79], geared to explain the lifting application of [JL70], mainly in the number field case. But none computes explicitly (and accurately, cf. [D81]) all the terms which appear in the trace formula. The latter is precisely what is needed for the counting applications of [D81] and [F]. An attempt at a complete exposition of the computations for $\mathrm{GL}(2)$ in the number field case is at [AFOO].
Of course the trace formula of [JL70] was generalized to the higher rank case by Arthur, see e.g. [A05], in the number field case, and by Lafforgue, see e.g. [Lf97], in the function field case. But the important applications of these works did not require explicit evaluation of all the terms which appear in the trace formula, so our results are not included in those of [Lf97], even in the case of GL(2) considered here.
In the number field case, the Remark on p. 112 of [A05] states: "As a matter of fact, it is only in the case of GL(2) that the general coefficients have been evaluated. It would be very interesting to understand them better in other examples, although this does not seem to be necessary for presently conceived applications of the trace formula". Indeed the applications of [D81], [DF13], [F] - counting rather than comparing - are of different nature than those of [JL70], [A05], [Lf97], where most terms can be erased a-priori in the comparison so they need not be computed.
To repeat what is explained above, we also think the approach of subsections 4.7-4.8 and section 5 is original, substantially different from the currently known methods (which are developed in [A05], [Lf97]), interesting and warrants further development.
I am deeply grateful to V. Drinfeld for making available to me his unpublished notes, for teaching me lots of mathematics in the process, and for his permission to publish this paper; to A. Beilinson for telling me at IHES about Drinfeld's notes; to the referee for the very careful reading. The author was a Schonbrunn visiting Professor at the Hebrew University, Jerusalem. Work
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### 1.2 Statement of the Trace Formula

Let us write the trace formula for GL(2) over a function field $F$ of a smooth projective geometrically connected curve $X$ over a finite field $\mathbb{F}_{q}$, and a test function $f$ in $C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$ (subscript $c$ for "compactly supported", superscript $\infty$ for "locally constant", $\mathbb{A}$ denotes the ring of adèles of $F$ ).
Let $r_{0}$ be the representation of $\mathrm{GL}(2, \mathbb{A})$ by right translation on the space $A_{0, \alpha}$ of cusp forms on $\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$, and $r_{0}(f)=\int f(g) r_{0}(g) d g$ $(g \in \mathrm{GL}(2, \mathbb{A}))$ the convolution operator; $d g=\otimes_{v} d g_{v}$ is a Haar measure. Here $\alpha$ is a fixed idèle of degree 1 , whose components are almost all equal to 1 .
A cusp form is a function $\phi: \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) \rightarrow E$ ( $E$ is a fixed algebraically closed subfield of $\mathbb{C}$ ) which is invariant on the right by some open compact subgroup of $\mathrm{GL}(2, \mathbb{A})$, and $\int_{N(F) \backslash N(\mathbb{A})} \phi(n x) d n=0$ for all $x$ in GL $(2, \mathbb{A})$. Here $N$ denotes the unipotent upper triangular subgroup of GL(2). We also write $A$ for the diagonal subgroup, and $A^{\prime}=A-Z$ where $Z$ is the center of GL(2). By a well known result of G. Harder, when $F$ is a function field (but not a number field) a cusp form is compactly supported modulo $Z(\mathbb{A})$.

Theorem 1.1. For any $f \in C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$ we have $\operatorname{tr} r_{0}(f)=\sum_{1 \leq i \leq 8} S_{i}(f)$. The terms are:

$$
\begin{gathered}
S_{1}(f)=\left|\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right| \sum_{\gamma \in \alpha^{\mathbb{Z}} \cdot F^{\times}} f(\gamma) . \\
S_{2}(f)=\sum_{F_{2}} S_{2, F_{2}}(f), \\
S_{2, F_{2}}(f)=\left|\operatorname{Aut}_{F} F_{2}\right|^{-1} \sum_{\gamma \in \alpha^{\mathbb{Z}}\left(F_{2}-F\right)} \int_{\mathrm{GL}(2, \mathbb{A}) / \alpha^{\mathbb{Z}} \cdot F_{2}^{\times}} f\left(x \gamma x^{-1}\right) d x .
\end{gathered}
$$

Here $F_{2}$ ranges over the set of isomorphism classes of quadratic extensions of the field $F$. For each $F_{2}$ we fix an embedding $F_{2} \hookrightarrow M(2, F)$ into the ring of $2 \times 2$ matrices over $F$.

$$
S_{3}(f)=\sum_{\gamma \in \alpha^{Z} \cdot A^{\prime}(F)} \int_{A(\mathbb{A}) \backslash \operatorname{GL}(2, \mathbb{A})} f\left(x^{-1} \gamma x\right) v(x) d x
$$

Any $x \in \mathrm{GL}(2, \mathbb{A})$ can be written in the form ank, $a \in A(\mathbb{A}), k \in \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$, $n=\left(\begin{array}{ll}1 & b \\ 0 & 1\end{array}\right), b$ is determined uniquely by $x$ up to $b \mapsto u b+w, u \in O_{\mathbb{A}}^{\times}, w \in O_{\mathbb{A}}$. $\operatorname{Put} v(x)=\sum_{v} \log _{q}\left(\max \left(1,\left|b_{v}\right|_{v}\right)\right)$.

$$
\begin{gathered}
S_{4}(f)=\sum_{a \in F^{\times} \alpha^{\mathbb{Z}}} \tilde{\theta}_{a, f}(1), \quad \tilde{\theta}_{a, f}(t)=\frac{1}{2}\left(\theta_{a, f}(t)+\theta_{a, f}\left(t^{-1}\right)\right), \\
\theta_{a, f}(t)=\int_{F^{\times} \alpha^{\mathbb{Z}} N(F) \backslash \operatorname{GL}(2, \mathbb{A})} f\left(x^{-1}\left(\begin{array}{cc}
a & a \\
0 & a
\end{array}\right) x\right) t^{\mathrm{ht}^{+}(x)} d x, \\
\text { DOCUMENTA MATHEMATICA } 19 \text { (2014) } 1-62
\end{gathered}
$$

$\mathrm{ht}^{+}: \mathrm{GL}(2, \mathbb{A}) \rightarrow \mathbb{Z}$ is defined by $\left.\mathrm{ht}^{+}\left(\begin{array}{cc}a & c \\ 0 & b\end{array}\right) k\right)=\operatorname{deg} a-\operatorname{deg} b\left(k \in \mathrm{GL}\left(2, O_{\mathbb{A}}\right)\right.$; $\left.a, b \in \mathbb{A}^{\times} ; c \in \mathbb{A}\right)$.

$$
S_{5}(f)=\frac{-1}{4 \pi i} \sum_{\mu_{1}, \mu_{2}} \oint_{|z|=1} \operatorname{tr} I\left(\mu_{1} \nu_{z}, \mu_{2} \nu_{z^{-1}}, f\right) \frac{m^{\prime}\left(\mu_{1} / \mu_{2}, z\right)}{m\left(\mu_{1} / \mu_{2}, z\right)} 2 z d z
$$

Here $m(\mu, z)=L(\mu, z) / L(\mu, z / q)$. The $\mu_{1}, \mu_{2}$ range over the set of characters of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}, \nu_{z}(x)=z^{\operatorname{deg}(x)}$. Also $I\left(\mu_{1}, \mu_{2}\right)$ is the space of right locally constant functions $\phi$ on $\mathrm{GL}(2, \mathbb{A})$ with

$$
\phi\left(\left(\begin{array}{ll}
a & c \\
0 & b
\end{array}\right) x\right)=|a / b|^{1 / 2} \mu_{1}(a) \mu_{2}(b) \phi(x) \quad\left(x \in \mathrm{GL}(2, \mathbb{A}) ; a, b \in \mathbb{A}^{\times} ; c \in \mathbb{A}\right) .
$$

It is a $\mathrm{GL}(2, \mathbb{A})$-module by right translation, and $\operatorname{tr} I\left(\mu_{1} \nu_{z}, \mu_{2} \nu_{z^{-1}}, f\right)$ is the trace of the indicated convolution operator.
$S_{6}(f)=\frac{-1}{4 \pi i} \sum_{\mu_{1}, \mu_{2}} \oint_{|z|=1} \operatorname{tr}\left[I\left(\mu_{1} \nu_{z}, \mu_{2} \nu_{z^{-1}}, f\right) \cdot R\left(\mu_{1}, \mu_{2}, z\right)^{-1} \frac{d}{d z} R\left(\mu_{1}, \mu_{2}, z\right)\right] d z$.
Notations are as in $S_{5}(f)$, and $R\left(\mu_{1}, \mu_{2}, z\right): I\left(\mu_{1} \nu_{z}, \mu_{2} \nu_{z^{-1}}\right) \rightarrow I\left(\mu_{2} \nu_{z^{-1}}, \mu_{1} \nu_{z}\right)$ is an operator, rational in $z$, defined as a product $\otimes_{v} R\left(\mu_{1 v}, \mu_{2 v}, z_{v}\right), z_{v}=$ $z^{\operatorname{deg}(v)}$. The product is well defined as the local operator maps the function in the source whose restriction to $\mathrm{GL}\left(2, O_{v}\right)$ is 1 to such function in the target. Further, $R\left(\mu_{1 v}, \mu_{2 v}, z\right)$ is defined to be

$$
\left[L\left(\mu_{1 v} / \mu_{2 v}, z^{2} / q_{v}\right) / L\left(\mu_{1 v} / \mu_{2 v}, z^{2}\right)\right] M\left(\mu_{1 v}, \mu_{2 v}, z\right)
$$

The operator $M\left(\mu_{1 v}, \mu_{2 v}, z\right)=M\left(\mu_{1 v} \nu_{z}, \mu_{2 v} \nu_{z^{-1}}\right)$ is defined first by an integral

$$
\phi \mapsto \int \phi\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)\left(\begin{array}{cc}
1 & y \\
0 & 1
\end{array}\right) x\right) d y \quad \text { if } \quad\left|\left(\mu_{1 v} / \mu_{2 v}\right)\left(\boldsymbol{\pi}_{v}\right) z^{2}\right|<1
$$

then by analytic continuation, as it is a rational function in $z$. The operators $I\left(\mu_{1} \nu_{z}, \mu_{2} \nu_{z^{-1}}, f\right)$ and $R\left(\mu_{1}, \mu_{2}, z\right)$ are considered as operators on

$$
\begin{gathered}
I_{0}\left(\mu_{1}, \mu_{2}\right)=\left\{\phi \in C^{\infty}\left(\mathrm{GL}\left(2, O_{\mathbb{A}}\right)\right) ; \phi\left(\left(\begin{array}{ll}
a & c \\
0 & b
\end{array}\right) x\right)=\mu_{1}(a) \mu_{2}(b) \phi(x) ;\right. \\
\left.x \in \mathrm{GL}\left(2, O_{\mathbb{A}}\right), a, b \in O_{\mathbb{A}}^{\times} ; c \in O_{\mathbb{A}}\right\} . \\
S_{7}(f)=\frac{1}{4} \sum_{\mu} \operatorname{tr} I(\mu, \mu, f), \quad S_{8}(f)=-\sum_{\mu} \int_{\mathrm{GL}(2, \mathbb{A})} f(x) \mu(\operatorname{det} x) d x .
\end{gathered}
$$

Both sums range over all characters $\mu$ of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{2 \mathbb{Z}}$. The sum of $S_{8}$ is over all automorphic one dimensional representations ( $\mu \circ \mathrm{det}$ ) of $\alpha^{\mathbb{Z}} \backslash \mathrm{GL}(2, \mathbb{A})$. The integral there represents the trace of the convolution operator associated with $f$.

The terms $S_{1}(f)$ and $S_{2}(f)$ are finite by Proposition 3.5, 3.6, 3.9. The argument used in the proof of Proposition 3.9 shows that for any $\gamma \in \alpha^{\mathbb{Z}}(A(F)-Z(F))$ the function $x \mapsto f\left(x^{-1} \gamma x\right)$ on $A(\mathbb{A}) \backslash \mathrm{GL}(2, \mathbb{A})$ has compact support, hence the integral in $S_{3}(f)$ converges.

By Proposition 3.11 the function $\theta_{a, f}(t)$ is rational and may have at $t=1$ a pole of order at most 1 , for each $a \in \mathbb{A}^{\times}$. Hence $\tilde{\theta}_{a, f}(t)$ is regular at $t=1$. From Proposition 3.5 it follows that the sums in $S_{3}(f)$ and $S_{4}(f)$ are finite, so these terms are well defined.
For any $f=\otimes f_{v}$ in $C_{c}^{\infty}(\operatorname{GL}(2, \mathbb{A}))$, the operator $I\left(\mu_{1}, \mu_{2}, f\right)$ is zero unless $\mu_{i}$ are unramified at each $v$ where $f_{v}$ is $\mathrm{GL}\left(2, O_{v}\right)$ biinvariant. This implies that the sums in $S_{i}(f)(5 \leq i \leq 8)$ are finite, for a given $f$. To see that $S_{5}(f)$ and $S_{6}(f)$ are well defined, note that the rational functions $m(\mu, t), R\left(\mu_{1}, \mu_{2}, t\right)$, $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ are regular on $|t|=1$ for all characters $\mu, \mu_{1}, \mu_{2}$ of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$. For $m(\mu, t)$ this follows from Proposition 4.11, for $R$ and $R^{-1}$ from Corollary 4.28 .

The distributions [linear forms on $\left.C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))\right] f \mapsto \operatorname{tr} r_{0}(f), S_{i}(f)(i=$ $1,2,5,7,8)$ are invariant, namely take the same value at $f$ and $f^{h}(x)=$ $f\left(h^{-1} x h\right), h \in \operatorname{GL}(2, \mathbb{A})$. For $i=3,4,6$ we have $S_{i}\left(f^{h}\right)=S_{i}(f)$ if $h \in$ $\mathrm{GL}\left(2, O_{\mathbb{A}}\right)$, but $S_{i}$ is not invariant.
If $f \in C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$ takes values in $\mathbb{Q}$ then $\operatorname{tr} r_{0}(f) \in \mathbb{Q}$, since the representation $r_{0}$ is defined over $\mathbb{Q}$. For $i=1,2,3,4,8$ it is clear that $S_{i}(f) \in \mathbb{Q}$. For $i=7$ the integrand contains the factor $\mu(a b)|a / b|^{1 / 2}$ which involves $\sqrt{q}$. However the sum includes with $\mu$ also $\mu \varepsilon, \varepsilon(\alpha)=-1$, and so the sum of the terms indexed by $\mu$ and $\mu \varepsilon$ can be written as an integral over the domain where $|a / b|$ is in $q^{2 \mathbb{Z}}$.
To see that $S_{5}(f)$ is rational, we put $a\left(\mu_{1}, \mu_{2}\right)=\frac{1}{2 \pi i} \oint_{|t|=1} f\left(\mu_{1}, \mu_{2}, t\right) d t$ where

$$
f\left(\mu_{1}, \mu_{2}, t\right)=\operatorname{tr} I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}, f\right) \cdot \frac{d}{d t} \ln m\left(\mu_{1} / \mu_{2}, t^{2}\right)
$$

and claim that for any $\sigma \in \operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q})$ one has $\sigma\left(a\left(\mu_{1}, \mu_{2}\right)\right)=a\left({ }^{\sigma} \mu_{1},{ }^{\sigma} \mu_{2}\right)$. Note that $\operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q})$ acts on the group of characters on $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$ as they are all $\mathbb{Q}$-valued. Now $a\left(\mu_{1}, \mu_{2}\right)$ is the sum of the residues of $f\left(\mu_{1}, \mu_{2}, t\right)$ at the points of the unit disc. We have that $\sigma\left(f\left(\mu_{1}, \mu_{2}, t\right)\right)=f\left({ }^{\sigma} \mu_{1},{ }^{\sigma} \mu_{2}, \varepsilon(\sigma) \cdot{ }^{\sigma} t\right)$ with $\varepsilon(\sigma)=\sigma(\sqrt{q}) / \sqrt{q}$. However, if $f\left(\mu_{1}, \mu_{2}, t\right)$ has a pole at $t=t_{0}$ and $\left|t_{0}\right|<1$, then by Proposition 4.11, $\left|\sigma\left(t_{0}\right)\right|<1$ for any $\sigma \in \operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q})$. Hence $S_{5}(f) \in \mathbb{Q}$. To see that $S_{6}(f) \in \mathbb{Q}$ one proceeds similarly, using the results of Corollary 4.28 on the poles of $R\left(\mu_{1}, \mu_{2}, t\right)$ and $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$.

## 2 Locally free sheaves of $\mathcal{O}_{X}$-Modules

### 2.1 Stable Bundles

Let $X$ be a smooth geometrically connected projective curve over $\mathbb{F}_{q}$ (we take minimal $q$ ). Denote by $\mathcal{O}_{X}$ the structure sheaf of $X$. Denote by Bun $n$ the set of isomorphism classes of rank $n$ locally free sheaves of $\mathcal{O}_{X}$-modules. By a (vector) bundle we mean here simply a locally free sheaf. In particular, $\mathrm{Bun}_{1}=\operatorname{Pic} X$. The Picard group Pic $X$ of invertible, or rank 1, locally free sheaves $\mathcal{L}$ of $\mathcal{O}_{X^{-}}$ modules, is naturally isomorphic to the group of classes $\bar{D}$ of (Weil) divisors $D=\sum_{v} n_{v} v\left(n_{v} \in \mathbb{Z}, v \in|X|\right)$. Here $|X|$ is the set of closed points of $X$,
and the divisors $D, D^{\prime}$ lie in the same class (are linearly equivalent) if their difference is the (principal) divisor $(f)=\sum_{v} \operatorname{ord}_{v}(f) v$ where $f$ is a nonzero rational function on $X$ and $\operatorname{ord}_{v}(f)$ is the order of $f$ at $v \in|X|\left(\operatorname{ord}_{v}(f)>0\right.$ if $v$ is a zero, $\operatorname{ord}_{v}(f)<0$ if $v$ is a pole, $\operatorname{ord}_{v}(f)=0$ otherwise). If $\mathcal{L}, \mathcal{M} \in \operatorname{Pic} X$ correspond to the divisors $D, D^{\prime}$ then $\mathcal{L} \otimes \mathcal{M}$ corresponds to $D+D^{\prime}$.
There is a degree map deg on Pic $X$ : $\operatorname{deg}\left(\sum_{v} n_{v} v\right)=\sum_{v} n_{v} \operatorname{deg}(v)$ defines $\operatorname{deg}(\mathcal{L})=\operatorname{deg}(D)$, where $\operatorname{deg}(v)=\left[k_{v}: \mathbb{F}_{q}\right]$. Here $k_{v}$ is the residue field of the function field $F=\mathbb{F}_{q}(X)$ of $X$ over $\mathbb{F}_{q}$ at $v$; assume $\mathbb{F}_{q}$ is algebraically closed in $F$. We write $F_{v}$ for the completion of $F$ at $v, O_{v}$ for its ring of integers. The cardinality of the residue field $k_{v}=\mathbb{F}_{q_{v}}$ at $v$ is denoted by $q_{v}$, thus $q_{v}=q^{\operatorname{deg}(v)}$. We also write $\operatorname{deg}(\bar{D})$ for $\operatorname{deg}(D)$, as the degree of a principal divisor is 0 ; recall that $\bar{D}$ denotes the class of $D$.
Denote by $\chi(\mathcal{L})=\operatorname{dim}_{\mathbb{F}_{q}} H^{0}(X, \mathcal{L})-\operatorname{dim}_{\mathbb{F}_{q}} H^{1}(X, \mathcal{L})$ the Euler-Poincaré characteristic of $\mathcal{L} \in \operatorname{Pic} X$. Here $H^{i}(X, \mathcal{L})$ are finite dimensional vector spaces over $\mathbb{F}_{q}$. Then $\chi\left(\mathcal{O}_{X}\right)=1-g$ where $g=\operatorname{dim}_{\mathbb{F}_{q}} H^{1}\left(X, \mathcal{O}_{X}\right)$ is named the genus of $X$. The Riemann-Roch theorem asserts that $\chi(\mathcal{L})-\operatorname{deg}(\mathcal{L})=\chi\left(\mathcal{O}_{X}\right)$ is independent of $\mathcal{L} \in \operatorname{Pic} X$.
Define the degree of a locally free sheaf $\mathcal{E}$ of $\mathcal{O}_{X}$-modules of rank $n$ to be $\operatorname{deg} \mathcal{E}=\chi(\mathcal{E})-n \chi\left(\mathcal{O}_{X}\right)$. The determinant of $\mathcal{E}$ is $\operatorname{det} \mathcal{E}=\bigwedge^{n} \mathcal{E} \in \operatorname{Pic} X$. We have $\operatorname{deg} \mathcal{E}=\operatorname{deg} \operatorname{det} \mathcal{E}$. This gives an alternative definition of the degree. A proof of this equality is as follows. If $\mathcal{E}$ is a line bundle, then there is nothing to prove. In the general case, use the fact that $\operatorname{both} \operatorname{deg} \mathcal{E}$ and $\operatorname{deg} \operatorname{det} \mathcal{E}$ are additive (if $\mathcal{E}^{\prime} \subset \mathcal{E}$ is a subbundle, then $\operatorname{deg} \mathcal{E}=\operatorname{deg} \mathcal{E}^{\prime}+\operatorname{deg}\left(\mathcal{E} / \mathcal{E}^{\prime}\right)$ and similarly for $\operatorname{deg} \operatorname{det} \mathcal{E})$, and that each vector bundle has a flag, $\mathcal{E}_{i}$, such that $\mathcal{E}_{i} / \mathcal{E}_{i-1}$ are line bundles.
The height of a rank two locally free sheaf $\mathcal{E}$ of $\mathcal{O}_{X}$-modules is the integer $\operatorname{ht}(\mathcal{E})=\max _{\mathcal{L}}(2 \operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{E}), \mathcal{L}$ ranges over all invertible subsheaves of $\mathcal{E}$.

Proposition 2.1. We have $-2 g \leq \operatorname{ht}(\mathcal{E})<\infty$.
Proof. Let $\mathcal{L}$ be an invertible subsheaf of $\mathcal{E}$. From the Riemann-Roch theorem $\chi(\mathcal{L})=\operatorname{deg} \mathcal{L}+1-g$ we obtain $\operatorname{dim}_{\mathbb{F}_{q}} H^{0}(X, \mathcal{L}) \geq \operatorname{deg} \mathcal{L}+1-g$, whence $\operatorname{deg} \mathcal{L} \leq \operatorname{dim}_{\mathbb{F}_{q}} H^{0}(X, \mathcal{L})+g-1 \leq \operatorname{dim}_{\mathbb{F}_{q}} H^{0}(X, \mathcal{E})+g-1$, so $\operatorname{ht}(\mathcal{E})$ is finite.
Let $\mathcal{L}$ be an invertible subsheaf of $\mathcal{E}$ of maximal degree. Let $\mathcal{M}$ be an invertible sheaf with $\operatorname{deg} \mathcal{M}=\operatorname{deg} \mathcal{L}+1$. Then $\operatorname{Hom}(\mathcal{M}, \mathcal{E})=0$. Also, by RiemannRoch for the rank $2 \operatorname{sheaf} \mathcal{E}, \operatorname{dim}_{\mathbb{F}_{q}} \operatorname{Hom}(\mathcal{M}, \mathcal{E})=\operatorname{dim}_{\mathbb{F}_{q}} H^{0}\left(X, \mathcal{M}^{-1} \mathcal{E}\right) \geq$ $\operatorname{deg}\left(\mathcal{M}^{-1} \mathcal{E}\right)+2-2 g=\operatorname{deg} \mathcal{E}-2 \operatorname{deg} \mathcal{M}+2-2 g=\operatorname{deg} \mathcal{E}-2 \operatorname{deg} \mathcal{L}-2 g$, so $2 \operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{E} \geq-2 g$.

A rank two locally free sheaf $\mathcal{E}$ of $\mathcal{O}_{X}$-modules is called stable if $\operatorname{ht}(\mathcal{E})<0$ and semistable if $\operatorname{ht}(\mathcal{E}) \leq 0$. In general, the slope $\mu(\mathcal{E})$ of a locally free sheaf $\mathcal{E}$ over an algebraic curve is defined to be $\operatorname{deg} \mathcal{E} / \mathrm{rk} \mathcal{E}$, and $\mathcal{E}$ is called stable if $\mu(\mathcal{F})<\mu(\mathcal{E})$ for all proper nonzero subbundles $\mathcal{F}$ of $\mathcal{E}$ (semistable if $\leq$ ). A locally free sheaf $\mathcal{E}$ of rank two is called almost stable if $\operatorname{ht}(\mathcal{E})<2 g-1$, and very unstable if $\operatorname{ht}(\mathcal{E}) \geq 2 g-1$. If $g=0$, every $\mathcal{E}$ is very unstable.

Remark 1. A very unstable vector bundle $\mathcal{E}$ of rank 2 splits into the direct sum of two line bundles. We give here a relatively elementary treatment. An extension can be found in the work of Harder and Narasimhan. If $\mathcal{E}$ is very unstable, $\mathcal{L}$ is an invertible subsheaf of $\mathcal{E}$ of maximal degree, and $\mathcal{M}=\mathcal{E} / \mathcal{L}$, then $\mathcal{M}$ is invertible and $\operatorname{Ext}(\mathcal{M}, \mathcal{L})=H^{1}\left(X, \mathcal{M}^{-1} \mathcal{L}\right)$ is 0 since

$$
\operatorname{deg} \mathcal{M}^{-1} \mathcal{L}=\operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{M}=2 \operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{E}=\operatorname{ht} \mathcal{E} \geq 2 g-1
$$

Indeed, by Serre duality $H^{1}\left(X, \mathcal{M}^{-1} \mathcal{L}\right)=H^{0}\left(X, \mathcal{L}^{-1} \mathcal{M} \omega\right)$ where $\omega$ denotes the canonical bundle. But $\operatorname{deg} \mathcal{L}^{-1} \mathcal{M} \omega \leq 2 g-2-(2 g-1)<0$, and $H^{0}(X, \mathcal{F})=0$ for an invertible sheaf $\mathcal{F}$ with negative degree.

Proposition 2.2. The number of isomorphism classes of almost stable rank two locally free sheaves $\mathcal{E}$ of $\mathcal{O}_{X}$-modules with a fixed degree is finite.

Proof. The height of an almost stable sheaf lies in $[-2 g, 2 g-2]$. Hence it suffices to show the finiteness for $\mathcal{E}$ with a fixed degree $n$ and height $h$. Every such sheaf lies in an exact sequence $0 \rightarrow \mathcal{L} \rightarrow \mathcal{E} \rightarrow \mathcal{M} \rightarrow 0$, where $\mathcal{L}$ and $\mathcal{M}$ are invertible sheaves and $2 \operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{E}=h$. Then $\operatorname{deg} \mathcal{L}=(n+h) / 2$, $\operatorname{deg} \mathcal{M}=(n-h) / 2$. Since the degrees of $\mathcal{L}$ and $\mathcal{M}$ are fixed, there are only finitely many possibilities for $\mathcal{L}$ and $\mathcal{M}$ (set of cardinality of the $\mathbb{F}_{q}$-points on the abelian variety $\left.\operatorname{Pic}^{0}(X)\right)$. With $\mathcal{L}$ and $\mathcal{M}$ fixed there are only finitely many choices for $\mathcal{E}$ as $\operatorname{Ext}(\mathcal{L}, \mathcal{M})$ is finite.

The group $\operatorname{Pic} X$ acts on $\mathrm{Bun}_{2}:\left(\mathcal{L} \in \operatorname{Pic} X, \mathcal{E} \in \mathrm{Bun}_{2}\right) \mapsto \mathcal{L} \otimes \mathcal{E}$. As

$$
\operatorname{deg}(\mathcal{L} \otimes \mathcal{E})=2 \operatorname{deg}(\mathcal{L})+\operatorname{deg}(\mathcal{E})
$$

the set of almost stable sheaves is invariant under this action. In a Pic $X$-orbit we may choose $\mathcal{E}$ to have $\operatorname{deg}(\mathcal{E})$ in $\{0,1\}$. Hence we deduce

Corollary 2.3. The number of $\operatorname{Pic} X$-orbits on the set of isomorphism classes of almost stable rank two locally free sheaves of $\mathcal{O}_{X}$-modules is finite.

### 2.2 Bundles and lattices

Let $\mathcal{E}$ be a rank $n$ locally free sheaf of $\mathcal{O}_{X}$-modules. Denote by $\mathcal{E}_{\eta}$ the fiber (= stalk) of $\mathcal{E}$ over the generic point $\eta$ of $X$. Let $\mathcal{E}_{(v)}$ be the stalk of $\mathcal{E}$ at the closed point $v \in|X|$. Let $O_{(v)}$ be the local ring of $X$ at $v$. Then $\mathcal{E}_{\eta}$ is an $n$-dimensional vector space over $F$, and $\mathcal{E}_{(v)}$ is an $O_{(v)}$-lattice in $\mathcal{E}_{\eta}$, namely a rank $n$ free $O_{(v)}$-submodule of $\mathcal{E}_{\eta}$.
A set $M$ of $O_{(v)}$-lattices $M_{(v)}$ in a finite dimensional vector space $V$ over $F, v$ ranges over the set $|X|$ of closed points in $X$, is called adelic if there exists a basis $\left\{e_{1}, \ldots, e_{n}\right\}$ in $V$ such that $M_{(v)}=O_{(v)} e_{1}+\cdots+O_{(v)} e_{n}$ for almost all $v$ in $|X|$. "Almost all" means "with at most finitely many exceptions". If $M$ is adelic then it is adelic with respect to any basis $\left\{e_{1}, \ldots, e_{n}\right\}$ of $V$.
The set of stalks $\left\{\mathcal{E}_{(v)} ; v \in|X|\right\}$ of a locally free sheaf $\mathcal{E}$ of $\mathcal{O}_{X}$-modules is adelic. Conversely, an adelic set of lattices $M=\left\{M_{(v)} ; v \in|X|\right\}$ in a finite
dimensional vector space $V$ over $F$ is the set of stalks of the locally free sheaf $\mathcal{E}$ of $\mathcal{O}_{X}$-modules defined by

$$
H^{0}(\mathrm{U}, \mathcal{E})=\left\{s \in V ; \forall v \in \mathrm{U}, s \in M_{(v)}\right\}
$$

for any open subset U of $X$. Obtained is an equivalence of the category of finite rank locally free sheaves of $\mathcal{O}_{X}$-modules, with the category of finite dimensional vector spaces over $F$ with adelic sets of $O_{(v)}$-lattices.
Let $O_{v}$ be the completion of $O_{(v)}$. The completion of $F$ at $v$ is denoted $F_{v}$. Let $V$ be a finite dimensional vector space over $F$. Put $V_{v}=V \otimes_{F} F_{v}$. There is a natural bijection between the set of $\mathcal{O}_{(v)}$-lattices in $V$, and $O_{v}$-lattices in $V_{v}$ : an $O_{(v)}$-lattice $M \subset V$ corresponds to the lattice $M \otimes_{\mathcal{O}_{(v)}} \mathcal{O}_{v}$ in $V_{v}$; an $O_{v}$-lattice $N \subset V_{v}$ corresponds to the $O_{(v)}$-lattice $N \cap V$.
The category $\mathcal{C}$ whose objects are finite dimensional $F$-vector spaces $V$ with adelic sets $\left\{M_{v} ; v \in|X|\right\}$ of $O_{v}$-lattices $M_{v}$ in $V_{v}$ is equivalent to the category of finite rank locally free sheaves of $\mathcal{O}_{X}$-modules $\mathcal{E}$, by $\mathcal{E} \mapsto\left(\mathcal{E}_{\eta},\left\{\mathcal{E}_{v}\right\}\right)$, where $\mathcal{E}_{\eta}$ is the generic fiber of $\mathcal{E}$ and $\mathcal{E}_{v}$ is the completion of the stalk of $\mathcal{E}$ at the closed point $v \in|X|$.
Let $R_{n}$ be the set of isomorphism classes of pairs $(\mathcal{E}, i)$ where $\mathcal{E}$ is a rank $n$ locally free sheaf of $\mathcal{O}_{X}$-modules, and $i$ is an isomorphism from the generic fiber of $\mathcal{E}$ to $F^{n}$. The pairs $(\mathcal{E}, i)$ and $\left(\mathcal{E}^{\prime}, i^{\prime}\right)$ are isomorphic if there is an isomorphism $\mathcal{E} \xrightarrow{\sim} \mathcal{E}^{\prime}$ which induces a commutative diagram when restricted to the generic fiber with sides $i$ and $i^{\prime}$ and the identity $F^{n} \rightarrow F^{n}$. The group $\mathrm{GL}(n, F)$ acts on $R_{n}$ by $g:(\mathcal{E}, i) \mapsto(\mathcal{E}, g \circ i)$. Then $\mathrm{GL}(n, F) \backslash R_{n}=\operatorname{Bun}_{n}$ is the set of isomorphism classes of rank $n$ locally free sheaves of $\mathcal{O}_{X}$-modules.
The set $R_{n}$ is the set of adelic collections of $O_{v}$-lattices $M_{v} \subset F_{v}^{n}, v \in|X|$. The group $\mathrm{GL}\left(n, F_{v}\right)$ acts transitively on the set of $O_{v}$-lattices in $F_{v}^{n}$. The stabilizer of the standard lattice $O_{v}^{n}$ in $F_{v}^{n}$ is GL $\left(n, O_{v}\right)$. Thus the set of $O_{v}$-lattices in $F_{v}^{n}$ is $\operatorname{GL}\left(n, F_{v}\right) / \mathrm{GL}\left(n, O_{v}\right)$, and $R_{n}$ is $\operatorname{GL}(n, \mathbb{A}) / \mathrm{GL}\left(n, O_{\mathbb{A}}\right)$, where $\mathbb{A}$ is the ring of adèles in $F$ and $O_{\mathbb{A}}=\prod_{v \in|X|} O_{v}$. Thus

$$
\operatorname{Bun}_{n}=\mathrm{GL}(n, F) \backslash \operatorname{GL}(n, \mathbb{A}) / \mathrm{GL}\left(n, O_{\mathbb{A}}\right)
$$

The elements of $\operatorname{GL}(n, \mathbb{A}) / \mathrm{GL}\left(n, O_{\mathbb{A}}\right)$ are called matrix divisors, and the elements of $\mathrm{GL}(n, F) \backslash \operatorname{GL}(n, \mathbb{A}) / \mathrm{GL}\left(n, O_{\mathbb{A}}\right)$ classes of matrix divisors. For $n=1$, the identification of $\mathrm{GL}(n, F) \backslash \mathrm{GL}(n, \mathbb{A}) / \mathrm{GL}\left(n, O_{\mathbb{A}}\right)$ with $\operatorname{Bun}_{n}$ is the identification of classes of divisors with invertible sheaves.
The group $\operatorname{GL}(n, \mathbb{A})$ can be identified with the set of triples

$$
\left(\mathcal{E}, i_{\eta}: \mathcal{E}_{\eta} \xrightarrow{\sim} F^{n},\left(i_{v}: \mathcal{E}_{v} \xrightarrow{\sim} O_{v}^{n}\right)\right) .
$$

Given a rank $n$ locally free sheaf $\mathcal{E}$, an isomorphism $i_{\eta}: \mathcal{E}_{\eta} \xrightarrow{\sim} F^{n}$, and for each closed point $v$ in $|X|$ an isomorphism $i_{v}: \mathcal{E}_{v} \xrightarrow{\sim} O_{v}^{n}$ of the completion $\mathcal{E}_{v}$ of the stalk $\mathcal{E}_{(v)}$ at $v$ with $O_{v}^{n}$, let us define the corresponding $g=\left(g_{v}\right)$ in $\operatorname{GL}(n, \mathbb{A})$. Each $g_{v}$ has to be an automorphism $F_{v}^{n} \rightarrow F_{v}^{n}$, with $g_{v}\left(O_{v}^{n}\right)=O_{v}^{n}$ for almost
all $v$. Construct $g_{v}$ as the composition $i_{v} \circ i_{\eta}^{-1}$ :

$$
F_{v}^{n}=F^{n} \otimes_{F} F_{v} \stackrel{i_{\eta}}{\sim} \mathcal{E}_{\eta} \otimes_{F} F_{v}=\mathcal{E}_{F_{v}}=\mathcal{E}_{v} \otimes_{O_{v}} F_{v} \stackrel{i_{v}}{\rightarrow} O_{v}^{n} \otimes_{O_{v}} F_{v}=F_{v}^{n}
$$

Note that since $\mathcal{E}$ is locally free, for almost all $v$ the map $g_{v}=i_{v} \circ i_{\eta}^{-1}$ takes $O_{v}^{n} \subset F_{v}^{n}$ to $\mathcal{E}_{v} \subset \mathcal{E}_{\eta} \otimes_{F} F_{v}$ via $i_{\eta}^{-1}$, and then to $O_{v}^{n}$ via $i_{v}$. To show that the map $\left\{\left(\mathcal{E}, i_{\eta},\left(i_{v}\right)\right)\right\} \rightarrow \mathrm{GL}(n, \mathbb{A})$ is bijective one shows that GL( $\left.n, \mathbb{A}\right)$ acts on the set of triples, simply transitively. Viewing the trivial locally free sheaf as $O_{\mathbb{A}}^{n}$ (space of columns), $g\left(\mathcal{E}, i_{\eta},\left(i_{v}\right)\right)$ is defined to be $\left(g \mathcal{E}, i_{\eta},\left(i_{v} \circ g_{v}^{-1}\right)\right)$, where $i_{v} \circ g_{v}^{-1}$ maps the stalk $g_{v} \mathcal{E}_{v}$ of $g \mathcal{E}$ at $v$ to $O_{v}^{n}$. The set of pairs $\left\{\left(\mathcal{E}, i_{\eta}\right)\right\}$ then corresponds to $\mathrm{GL}(n, \mathbb{A}) / \mathrm{GL}\left(n, O_{\mathbb{A}}\right)$, the set of pairs $\left\{\left(\mathcal{E},\left(i_{v}\right)\right)\right\}$ to $\operatorname{GL}(n, F) \backslash \mathrm{GL}(n, \mathbb{A})$, and the set $\{\mathcal{E}\}$ to $\operatorname{GL}(n, F) \backslash \operatorname{GL}(n, \mathbb{A}) / \operatorname{GL}\left(n, O_{\mathbb{A}}\right)$.
To an idèle $a=\left(\boldsymbol{\pi}_{v}^{-n_{v}} u_{v} ; v \in|X|\right)$, where $\boldsymbol{\pi}_{v}$ denotes a generator of the maximal ideal in the ring $O_{v}$ of integers in $F_{v}, u_{v} \in O_{v}^{\times}$and $n_{v} \in \mathbb{Z}$, we associate the divisor $D=\sum_{v} n_{v} v$, and the degree

$$
\operatorname{deg}(a)=\operatorname{deg}(D)=\sum_{v} n_{v} \operatorname{deg}(v), \quad \operatorname{deg}(v)=\left[\mathbb{F}_{v}: \mathbb{F}_{q}\right]
$$

where $\mathbb{F}_{v}$ is the residue field of $F$ at $v$, a finite field of $q_{v}=q^{\operatorname{deg}(v)}$ elements. For $g \in \mathrm{GL}(2, \mathbb{A})$ write $\operatorname{deg} g$ for $\operatorname{deg} \operatorname{det} g$. Recall that $O_{\mathbb{A}}=\prod_{v} O_{v}(v \in|X|)$. For $t \in \mathbb{C}^{\times}$we write

$$
\nu_{t}(a)=t^{-\operatorname{deg}(a)}=\prod_{v} t_{v}^{-n_{v}}
$$

where $t_{v}=t^{\operatorname{deg}(v)}$. Then $\nu_{q^{-1}}(a)=\prod_{v} q_{v}^{n_{v}}=|a|$ is equal to $\nu(a)=q^{\operatorname{deg}(a)}$. Also $\nu_{t}\left(\boldsymbol{\pi}_{v}\right)=t_{v}, \nu_{q^{-1}}\left(\boldsymbol{\pi}_{v}\right)=\left|\boldsymbol{\pi}_{v}\right|$.
Let $\mathcal{L}$ and $\mathcal{M}$ be invertible sheaves. Fix isomorphisms $i_{\mathcal{L}}, i_{\mathcal{M}}$ of their generic fibers with $F$. Each of $\left(\mathcal{L}, i_{\mathcal{L}}\right)$ and $\left(\mathcal{M}, i_{\mathcal{M}}\right)$ defines an element of $\mathbb{A}^{\times} / O_{\mathbb{A}}^{\times}$, namely a divisor on $X$. Choose representatives $a, b$ in $\mathbb{A}^{\times}$, for example $\sum_{v} n_{v} v$ is represented by $\left(\boldsymbol{\pi}_{v}^{-n_{v}}\right)$. Given an exact sequence $0 \rightarrow \mathcal{L} \rightarrow \mathcal{E} \rightarrow \mathcal{M} \rightarrow 0$ of locally free sheaves, choose an isomorphism $\varphi$ between the generic fiber of $\mathcal{E}$ and $F^{2}$ so that the induced exact sequence of generic fibers $0 \rightarrow F \rightarrow F^{2} \rightarrow F \rightarrow 0$ is standard $\left(x \mapsto\binom{x}{0},\binom{x}{y} \mapsto y\right)$. The isomorphism $\varphi$ is defined uniquely up to left multiplication by an automorphism of $F^{2}$ of the form $\left(\begin{array}{ll}1 & t \\ 0 & 1\end{array}\right), t \in F$. The pair $(\mathcal{E}, \varphi)$ determines an element of $\operatorname{GL}(2, \mathbb{A}) / \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$, of the form $u=\left(\begin{array}{ll}1 & z \\ 0 & 1\end{array}\right)\left(\begin{array}{ll}a & 0 \\ 0 & b\end{array}\right)$, with $z$ in $\mathbb{A}$. Since $u$ is defined up to right multiplication by an element of $\mathrm{GL}(2, O), z$ is uniquely defined up to addition of an element of $\frac{a}{b} O_{\mathbb{A}}$. Replacing $\varphi$ by $\left(\begin{array}{cc}1 & t \\ 0 & 1\end{array}\right) \varphi$ with $t \in F$ replaces $z$ by $z+t$. Thus we get a bijection

$$
\operatorname{Ext}(\mathcal{M}, \mathcal{L}) \rightarrow \mathbb{A} /\left(F+\frac{a}{b} O_{\mathbb{A}}\right)
$$

This is an isomorphism of $\mathbb{F}_{q}$-vector spaces.
In summary, if the invertible sheaves $\mathcal{L}$ and $\mathcal{M}$ correspond to idèles $a$ and $b$, then $\operatorname{Ext}(\mathcal{M}, \mathcal{L}) \simeq \mathbb{A} /\left(F+\frac{a}{b} O_{\mathbb{A}}\right)$, and the $\operatorname{map} \operatorname{Ext}(\mathcal{M}, \mathcal{L}) \rightarrow \mathrm{Bun}_{2}$ which associates to the exact sequence $0 \rightarrow \mathcal{L} \rightarrow \mathcal{E} \rightarrow \mathcal{M} \rightarrow 0$ its middle term,
coincides with the map $\mathbb{A} /\left(F+\frac{a}{b} O_{\mathbb{A}}\right) \simeq H^{1}\left(X, \mathcal{M}^{-1} \mathcal{L}\right)$, see [S97], II. 5. The isomorphism $\mathbb{A} /\left(F+\frac{a}{b} O_{\mathbb{A}}\right) \xrightarrow{\sim} \operatorname{Ext}(\mathcal{M}, \mathcal{L})$ is $H^{1}\left(X, \mathcal{M}^{-1} \mathcal{L}\right) \xrightarrow{\sim} \operatorname{Ext}(\mathcal{M}, \mathcal{L})$.

### 2.3 The space $\operatorname{GL}(2, F) \backslash \operatorname{GL}(2, \mathbb{A})$

Proposition 2.4. Given $a \in \mathbb{A}^{\times}, \operatorname{deg} a \geq 2 g-1$, then $a O_{\mathbb{A}}+F=\mathbb{A}$.
Proof. If $\mathcal{L}$ is an invertible sheaf on $X$ associated with $a$, then $\mathbb{A} /\left(F+a O_{\mathbb{A}}\right)=$ $H^{1}(X, \mathcal{L})$. By Serre duality $H^{1}(X, \mathcal{L}) \simeq H^{0}\left(X, \mathcal{L}^{-1} \omega\right)$, where $\omega$ is the canonical bundle of degree $2 g-2$. Then $\operatorname{deg}\left(\mathcal{L}^{-1} \omega\right) \leq(2 g-2)-(2 g-1)=-1<0$, hence $H^{0}\left(X, \mathcal{L}^{-1} \omega\right)=\{0\}$.

Define a function

$$
\mathrm{ht}^{+}: \mathrm{GL}(2, \mathbb{A}) \rightarrow \mathbb{Z} \quad \text { by } \quad \mathrm{ht}^{+}\left(\left(\begin{array}{ll}
a & c \\
0 & b
\end{array}\right) k\right)=\operatorname{deg} a-\operatorname{deg} b
$$

for all $a, b \in \mathbb{A}^{\times}, c \in \mathbb{A}, k \in \operatorname{GL}\left(2, O_{\mathbb{A}}\right)$. It is clearly a well defined function on $B(F) \backslash \mathrm{GL}(2, \mathbb{A})$. For $x \in \mathrm{GL}(2, \mathbb{A})$, put

$$
\operatorname{ht}(x)=\max _{\gamma \in \operatorname{GL}(2, F)} \mathrm{ht}^{+}(\gamma x) .
$$

On $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ it is well defined.
Proposition 2.5. For any $x \in \mathrm{GL}(2, \mathbb{A})$ we have $-2 g \leq \mathrm{ht}(x)<\infty$.
Proof. This follows from Proposition 2.1 as if $\mathcal{E}$ is a rank two locally free sheaf of $\mathcal{O}_{X}$-modules associated to the image of $x$ in $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) / \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$, then $\operatorname{ht}(x)=\operatorname{ht}(\mathcal{E})$.

Put $H_{B}=\left\{x \in B(F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \mathrm{ht}^{+}(x)>0\right\}$ and

$$
H=\{x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \operatorname{ht}(x)>0\} .
$$

Proposition 2.6. (1) The restriction $p$ to $H_{B}$ of the natural projection $p^{\prime}$ : $B(F) \backslash \mathrm{GL}(2, \mathbb{A}) \rightarrow \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ is a homeomorphism $H_{B} \rightarrow H$.
(2) The set $\{x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \mathrm{ht}(x) \leq n\}$ is compact modulo the center $Z(\mathbb{A})$ of $\mathrm{GL}(2, \mathbb{A})$ for every integer $n$.

Proof. (1) The map $p$ is clearly onto. To show that $p$ is injective it suffices to show for any $x$ in $\operatorname{GL}(2, \mathbb{A}), \gamma \in \mathrm{GL}(2, F)$, that $\mathrm{ht}^{+}(x)>0$ and $\mathrm{ht}^{+}(\gamma x)>0$ implies $\gamma \in B(F)$. This is a typical application of the Harder-Narasimhan filtration. In simple, explicit terms, this follows from

Lemma 2.7. If $g \in \mathrm{GL}(2, F)-B(F)$ then $\mathrm{ht}^{+}(x)+\mathrm{ht}^{+}(g x) \leq 0$.
Proof. Write $g$ as $g_{1} w g_{2}$ with $g_{1}, g_{2}$ in $B(F), w=\left(\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right)$. Put $x^{\prime}=g_{2} x$. Then $\mathrm{ht}^{+}(x)=\mathrm{ht}^{+}\left(x^{\prime}\right), \mathrm{ht}^{+}(g x)=\mathrm{ht}^{+}\left(w x^{\prime}\right)$. Thus we need to show that $\mathrm{ht}^{+}\left(x^{\prime}\right)+\mathrm{ht}^{+}\left(w x^{\prime}\right) \leq 0$. Suppose $x^{\prime}=\left(\begin{array}{cc}a_{1} & c_{1} \\ 0 & b_{1}\end{array}\right) k_{1}, w x^{\prime}=\left(\begin{array}{cc}a_{2} & c_{2} \\ 0 & b_{2}\end{array}\right) k_{2}$ with
$k_{1}, k_{2} \in \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$. Put $k_{2} k_{1}^{-1}=\left(\begin{array}{cc}\alpha & \beta \\ \gamma & \delta\end{array}\right)$. Then $\left(\begin{array}{cc}a_{2} & c_{2} \\ 0 & b_{2}\end{array}\right)\left(\begin{array}{cc}\alpha & \beta \\ \gamma & \delta\end{array}\right)=w\left(\begin{array}{cc}a_{1} & c_{1} \\ 0 & b_{1}\end{array}\right)=$ $\left(\begin{array}{ccc}0 & b_{1} \\ a_{1} & c_{1}\end{array}\right)$. Hence $b_{2} \gamma=a_{1}$. Thus $\operatorname{deg} a_{1} \leq \operatorname{deg} b_{2}\left(\right.$ as $\operatorname{deg} \gamma \leq 0$, since $\left.\gamma \in O_{\mathbb{A}}\right)$. But $\operatorname{deg} a_{2} b_{2}=\operatorname{deg} a_{1} b_{1}$. Hence $\operatorname{deg} a_{2} \leq \operatorname{deg} b_{1}$. Then $\mathrm{ht}^{+}\left(x^{\prime}\right)+\mathrm{ht}^{+}\left(w x^{\prime}\right)=$ $\operatorname{deg} a_{1}-\operatorname{deg} b_{1}+\operatorname{deg} a_{2}-\operatorname{deg} b_{2} \leq 0$.

Now the natural projection $p^{\prime}: B(F) \backslash \mathrm{GL}(2, \mathbb{A}) \rightarrow \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ is open and $H_{B}$ is an open subset of $B(F) \backslash \mathrm{GL}(2, \mathbb{A})$. Hence the bijection $p=p^{\prime} \mid H_{B}$ : $H_{B} \rightarrow H$ is open. Since it is also continuous, $p$ is a homeomorphism.
(2) The image of the set $S=\left\{x \in B(F) \backslash \mathrm{GL}(2, \mathbb{A}) ;-2 g \leq \mathrm{ht}^{+}(x) \leq n\right\}$ in $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ under $p^{\prime}$ contains the set $\{x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \operatorname{ht}(x) \leq$ $n\}$. So it suffices to show that $S$ is compact $\bmod Z(\mathbb{A})$. Choose a compact $C$ in $\mathbb{A}^{\times}$with

$$
C F^{\times}=\left\{t \in \mathbb{A}^{\times} ;-2 g \leq \operatorname{deg} t \leq n\right\} .
$$

Choose an idèle $d$ with $\operatorname{deg} d \geq 2 g-1$. Put

$$
Y=\left\{\left(\begin{array}{ll}
1 & c \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
a & 0 \\
0 & b
\end{array}\right) k ; \quad k \in \mathrm{GL}\left(2, O_{\mathbb{A}}\right), \quad a, b \in \mathbb{A}^{\times}, \quad \frac{a}{b} \in C, c \in d O_{\mathbb{A}}\right\} .
$$

Lemma 2.8. The map $Y \rightarrow S$ is surjective.
Proof. Let $x \in \operatorname{GL}(2, \mathbb{A}),-2 g \leq \mathrm{ht}^{+}(x) \leq n$. We need to show that $x$ can be written as $h y$ with $y \in Y$ and $h \in B(F)$. Write $x$ as $\left(\begin{array}{cc}r & s \\ 0 & t\end{array}\right) K$ with $k \in$ $\mathrm{GL}\left(2, O_{\mathbb{A}}\right), r, t \in \mathbb{A}^{\times}, s \in \mathbb{A}$. It remains to show that $\left(\begin{array}{cc}r & s \\ 0 & t\end{array}\right)$ can be expressed as $\left(\begin{array}{ll}\alpha & \gamma \\ 0 & \beta\end{array}\right)\left(\begin{array}{ll}1 & c \\ 0 & 1\end{array}\right)\left(\begin{array}{ll}a & 0 \\ 0 & b\end{array}\right)$ with $a, b \in \mathbb{A}^{\times}, \frac{a}{b} \in C, c \in d O_{\mathbb{A}}, \alpha, \beta \in F^{\times}, \gamma \in F$. Thus we need to show the existence of $a, b, c, \alpha, \beta, \gamma$ such that

$$
\begin{array}{ll}
\left({ }^{*}\right) & a \alpha=r, \quad \beta b=t, \quad a, b \in \mathbb{A}^{\times}, \quad \alpha, \beta \in F^{\times}, \quad \frac{a}{b} \in C, \\
\left({ }^{* *}\right) & b(\alpha c+\gamma)=s, \quad c \in d O_{\mathbb{A}}, \quad \gamma \in F .
\end{array}
$$

By definition of $x, \operatorname{deg} r-\operatorname{deg} t$ lies in $[-2 g, n]$, so the existence of $a, b, \alpha, \beta$ satisfying $\left(^{*}\right)$ follows from the definition of $C$. The existence of $c \in d O_{\mathbb{A}}$ and $\gamma \in F$ satisfying $\alpha c+\gamma=s / b$ follows from: $c O_{\mathbb{A}}+F=\mathbb{A}$ if $\operatorname{deg} c \geq 2 g-1$.

Since $Y$ is compact $\bmod Z(\mathbb{A})$, so is $S$, and (2) follows.
In summary, the homogeneous space $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ is the union of the compact $\bmod Z(\mathbb{A})$ set $\{x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \operatorname{ht}(x) \leq 0\}$, and the set $H=\{x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \operatorname{ht}(x)>0\}$, whose structure is simpler. The set $H_{B}$, hence also the sets $H$ and $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$, are noncompact modulo $Z(\mathbb{A})$. Indeed the function $\mathrm{ht}^{+}$takes arbitrary large values.
The image of $H$ in $\operatorname{Bun}_{2}=\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) / \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$ is the set of nonsemistable locally free sheaves.
The set $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) / \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$ is analogous to the set

$$
\mathrm{SL}(2, \mathbb{Z}) \backslash \mathrm{SL}(2, \mathbb{R}) / \mathrm{SO}(2)=\mathrm{SL}(2, \mathbb{Z}) \backslash \mathfrak{h},
$$

where $\mathfrak{h}=\{z \in \mathbb{C} ; \operatorname{Im} z>0\}$, the upper half plane, is isomorphic to $\mathrm{SL}(2, \mathbb{R}) / \mathrm{SO}(2)$, by

$$
g \mapsto g(i)=(a i+b) /(c i+d)
$$

The set $B(F) \backslash \mathrm{GL}(2, \mathbb{A}) / \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$ is analogous to $N \backslash \mathfrak{h}$ where $N$ is the group of transformations $z \mapsto z+n(n \in \mathbb{Z})$ on $\mathfrak{h}$. The function $\mathrm{ht}^{+}$is analogous to the function $z \mapsto \ln \operatorname{Im} z$ on $N \backslash \mathfrak{h}$. The statement $-2 g \leq \operatorname{ht}(x)<\infty$ corresponds to the statement that the natural map from the half plane $\{z \in \mathbb{C} ; \operatorname{Im} z \geq \sqrt{3} / 2\}$ to $\mathrm{SL}(2, \mathbb{Z}) \backslash \mathfrak{h}$ is onto. The statement that $p: H_{B} \rightarrow H$ is homeomorphism corresponds to the statement that the map $N \backslash\{z \in \mathbb{C} ; \operatorname{Im} z>1\} \rightarrow \mathrm{SL}(2, \mathbb{Z}) \backslash \mathfrak{h}$ is injective, and the compactness of $\{x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \operatorname{ht}(x) \leq n\}$ corresponds to the statement that the complement in $\mathrm{SL}(2, \mathbb{Z}) \backslash \mathfrak{h}$ of the image of the half plane $\{z \in \mathbb{C} ; \operatorname{Im} z>h\}$ is compact.

## $2.4 \ell$-GROUPS

An $\ell$-space is a Hausdorff topological space such that each of its points has a fundamental system of open compact neighborhoods.
We shall consider on $\ell$-spaces only measures for which every open compact subset is measurable, and its volume is a rational number. If $d x$ is such a measure on an $\ell$-space $Y$, and $f$ is a locally constant compactly supported function on $Y$ with values in a field $E$ of characteristic zero, then $\int_{Y} f(x) d x$ reduces to a finite sum, and it is well defined.
On topological groups we consider only left- or right-invariant measures.
An $\ell$-group is a topological group with an $\ell$-space structure.
Proposition 2.9. Let $G$ be an $\ell$-group. Then (1) there exists a fundamental system of neighborhoods of the identity in $G$ consisting of open compact subgroups;
(2) there exists a left Haar measure on $G$ such that the volume of each open compact set is a rational number.

Proof. (1) Let $U$ be a neighborhood of the identity in $G$. We shall show that $U$ contains an open compact subgroup. Since $G$ is $\ell$-space, we may assume that $U$ is open and compact. Put $V=\{x \in G ; x U \subset U\}$. Then $V=\cap_{u \in U} U u^{-1}$, hence it is compact. Now for each $v$ in $V$ and $u$ in $U$, by continuity of multiplication $m$ there exists an open subset $W_{u}$ containing $v$, and $U_{u}$ in $U$ containing $u$, such that $m\left(W_{u}, U_{u}\right) \subset U$. As $U$ is compact and $U=\cup_{u \in U} U_{u}$, there are finitely many $u_{1}, \ldots, u_{n}$ in $U$ with $U=\cup_{1 \leq i \leq n} U_{u_{i}}$. Then $W=\cap_{1 \leq i \leq n} W_{u_{i}}$ is open in $V$ and it contains $v$. Thus $V$ is an open neighborhood of the identity, and $V \cdot V=V$. Then $V \cap V^{-1}$ is an open compact subgroup in $U$.
(2) Fix some left Haar measure on $G$. Denote the volume of an open compact subgroup $U$ by $|U|$. For two such groups, $U_{1}$ and $U_{2}$ we have

$$
\frac{\left|U_{1}\right|}{\left|U_{2}\right|}=\frac{\left|U_{1}\right|}{\left|U_{1} \cap U_{2}\right|} / \frac{\left|U_{2}\right|}{\left|U_{1} \cap U_{2}\right|}=\frac{\left[U_{1}: U_{1} \cap U_{2}\right]}{\left[U_{2}: U_{1} \cap U_{2}\right]} \in \mathbb{Q}
$$

Consequently the Haar measure on $G$ can be chosen to assign rational volume to every open compact subgroup of $G$. But then the volume of every open compact subset $K$ in $G$ is rational, since as in (1) for such $K$ there is a compact
open subgroup $U$ of $G$ with $K U \subset K$, and then $|K|=[K: U]|U|$ is rational, where $K$ is a disjoint union of $[K: U]$ translates of $U$.

Fix an $\ell$-group $G$ and a left Haar measure on $G$ such that the volume of any open compact set is a rational number. Fix a field $E$ of characteristic zero. The $E$ vector space $H_{G}$ of compactly supported locally constant functions $f: G \rightarrow E$ is an algebra under the convolultion $\left(f_{1} * f_{2}\right)(g)=\int_{G} f_{1}(h) f_{2}\left(h^{-1} g\right) d h$. For an open compact subgroup $U$ in $G$ the set of $U$-biinvariant functions in $H_{G}$ is a subalgebra $H_{G}^{U}$, called the Hecke algebra of $(G, U)$. Although $H_{G}$ has no unit (unless $G$ is discrete, when the $\delta$-function is in $H_{G}$ ), $H_{G}^{U}$ does: it is $\delta_{U}: G \rightarrow \mathbb{Q}$, the characteristic function of $U$ divided by $|U|$.
A representation $\pi$ of the group $G$ on a vector space $V$ is called smooth if the stabilizer of any vector of $V$ is open, and admissible if it is smooth and for any open subgroup $U$ of $G$ the space $V^{U}$ of $U$-fixed vectors in $V$ is finite dimensional.
If $\pi$ is a smooth representation of an $\ell$-group $G$ on a vector space $V$ over $E$, for each $f \in H_{G}$ define the operator $\pi(f): V \rightarrow V$ by $\pi(f) v=\int_{G} f(g) \pi(g) v d g$. This integral reduces to a finite sum since $\pi$ is smooth, and $\pi\left(f_{1} * f_{2}\right)=$ $\pi\left(f_{1}\right) \circ \pi\left(f_{2}\right)$. Then $V$ is naturally an $H_{G}$-module, and for any open compact subgroup $U$ of $G$, the space $V^{U}$ is a unital module over $H_{G}^{U}$.

Proposition 2.10. (1) A smooth $G$-module $V \neq\{0\}$ is irreducible iff for every open compact subgroup $U$ of $G$ either $V^{U}=0$ or $V^{U}$ is an irreducible $H_{G}^{U}$-module.
(2) Given an open compact subgroup $U$ of $G$ and an irreducible unital $H_{G}^{U}$ module $M$, there exists a smooth irreducible $G$-module $V$ such that $V^{U}$ is isomorphic to $M$ as an $H_{G}^{U}$-module, and $V$ is determined by this property up to isomorphism.

For a proof see [BZ76], 2.10. See [BZ76], 2.11 for
Schur's Lemma. Let $\pi$ be an irreducible admissible representation of $G$ in a vector space $V$ over an algebraically closed field $E$. Then any nonzero $G$-module morphism (intertwining operator) $V \rightarrow V$ is a scalar.

Proposition 2.11. Let $\pi$ be an irreducible admissible representation of $G$ in a vector space $V$ over an algebraically closed field $E$. For any field extension $E^{\prime}$ of $E$, the representation of $G$ in $V \otimes_{E} E^{\prime}$ is also irreducible.

Proof. By Proposition 2.10, the statement reduces to a similar statement for finite dimensional algebras, since $\pi$ is assumed to be admissible.

Let $E$ be a subfield of $\mathbb{C}$ invariant with respect to complex conjugation. A representation of $G$ on a vector space $V$ over $E$ is unitary if there is a $G$ invariant scalar product on $V$ (thus a bilinear function $(\cdot, \cdot): V \times V \rightarrow E$ with $\overline{(v, w)}=(w, v)$ and $(v, v)=0$ iff $v=0$, and $(g v, g w)=(v, w)$ for all $v, w$ in $V$ and $g$ in $G$ ).

Note that we do not require $V$ to be complete with respect to the scalar product, even in the case $E=\mathbb{C}$. If $E$ is algebraically closed and the representation of $G$ in $E$ is irreducible and admissible, then the $G$-invariant inner product on $V$ is unique up to a scalar multiple, if it exists.

Proposition 2.12. Let $\pi$ be an admissible unitary representation of $G$ in the $E$-space $V$. Fix a $G$-invariant scalar product on $V$. Let $L$ be an invariant subspace of $V$, and $L^{\perp}$ its orthogonal complement. Then $V=L \oplus L^{\perp}$.
Proof. Given $x \in V$, we need to express it as $x_{1}+x_{2}$ with $x_{1} \in L$ and $x_{2} \in$ $L^{\perp}$. Since $\pi$ is smooth there exists a compact open subgroup $U$ of $G$ with $x \in V^{U}$. Since $\pi$ is admissible, $\operatorname{dim}_{E} V^{U}$ is finite. Thus $x=x_{1}+x_{2}$ for some $x_{1} \in L^{U}, x_{2} \in V^{U}, x_{2}$ orthogonal to $L^{U}$. It remains to show that $x_{2}$ is orthogonal to the entire space $L$. Let $\delta_{U}$ be the unit in $H_{G}^{U}$. Then $\pi\left(\delta_{U}\right)$ is the orthogonal projector $V \mapsto V^{U}$. Hence for every $y$ in $L,\left(x_{2}, y\right)=\left(\pi\left(\delta_{U}\right) x_{2}, y\right)=$ $\left(x_{2}, \pi\left(\delta_{U}\right) y\right)=0$ since $\pi\left(\delta_{U}\right) y \in L^{U}$.

It follows that every admissible unitary representation of $G$ is a direct sum of irreducible representations. This sum is not necessarily finite. However, given an open compact subgroup $U$ of $G$, only finitely many summands contain nonzero $U$-invariant vectors.

### 2.5 Automorphic forms

Let $E$ be an algebraically closed field of characteristic zero. An automorphic form is a smooth function $\phi: \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) \rightarrow E$, where by smooth we mean that there is an open subgroup $\mathrm{U}_{\phi}$ of $\mathrm{GL}(2, \mathbb{A})$ such that $\phi(x u)=\phi(x)$ for all $u \in \mathrm{U}_{\phi}$ and $x \in \mathrm{GL}(2, \mathbb{A})$. A cusp form is an automorphic form $\phi$ with $\int_{\mathbb{A} / F} \phi\left(\left(\begin{array}{ll}1 & z \\ 0 & 1\end{array}\right) x\right) d z=0$ for all $x \in \mathrm{GL}(2, \mathbb{A})$.
Since $\phi$ is right locally constant ( $=$ smooth) and $\mathbb{A} / F$ is compact, the integral here is well defined and reduces to a finite sum.
Let $A_{0}^{E}$ be the space of cusp forms $\phi: \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) \rightarrow E$. The group $\mathrm{GL}(2, \mathbb{A})$ acts on $A_{0}^{E}$ by right translation: $(r(h) \phi)(g)=\phi(g h)$. By a character of an $\ell$-group $G$ with values in $E$ we mean a locally constant homomorphism $\chi: G \rightarrow E^{\times}$. If $E \subset \mathbb{C}$ such $\chi$ is called a unitary character if $|\chi(g)|=1$ for all $g$ in $G$.
Denote by $A_{0}^{E}(\chi)$ the space of $\phi \in A_{0}^{E}$ with $\phi(a x)=\chi(a) \phi(x), a \in \mathbb{A}^{\times}$(identified with the center of $\mathrm{GL}(2, \mathbb{A})), x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$. The space $A_{0}^{E}(\chi)$ is invariant under the GL $(2, \mathbb{A})$-action.
Let $\pi$ be an irreducible representation of $\mathrm{GL}(2, \mathbb{A})$ over $E$. By Schur's lemma, there is a character $\chi: \mathbb{A}^{\times} \rightarrow E^{\times}$such that for every $a$ in $\mathbb{A}^{\times}, \pi(a)$ is multiplication by $\chi(a)$. This $\chi$ is called the central character of $\pi$.
If $V \subset A_{0}^{E}$ is an irreducible admissible representation $\pi$ of $\mathrm{GL}(2, \mathbb{A})$ and $\chi$ is the central character of $V$, then $V \subset A_{0}^{E}(\chi)$. Since the center of GL $(2, F)$ acts trivially on $A_{0}^{E}, \chi$ is trivial on $F^{\times}$. Thus every irreducible admissible $\pi \subset A_{0}^{E}$ lies in $A_{0}^{E}(\chi)$, where $\chi$ is the central character of $\pi$, which is a character of $\mathbb{A}^{\times} / F^{\times}$. The following is known also e.g. for $\operatorname{GL}(n)$.

Proposition 2.13. Fix an open subgroup $U$ of $\mathrm{GL}(2, \mathbb{A})$. There exists a compact $\bmod Z(\mathbb{A})$ subset $K$ of $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ such that the support of any $U$-invariant cusp form is contained in $K$.

Proof. We first show that there is an integer $n$ such that given $z \in \mathbb{A}$ and $x \in \operatorname{GL}(2, \mathbb{A})$ with ht $^{+}(x) \geq n$, there exist $u \in U$ and $\beta \in F$ with $\left(\begin{array}{cc}1 & z \\ 0 & 1\end{array}\right) x=$ $\left(\begin{array}{cc}1 & \beta \\ 0 & 1\end{array}\right) x u$.
To see this, fix an effective divisor $-D=\sum_{v \in|X|} n_{v} v$ on $X$, put $d=\left(\boldsymbol{\pi}_{v}^{n_{v}}\right)$ and let $J_{D}=d O_{\mathbb{A}}$ be the corresponding ideal in $O_{\mathbb{A}}$. The groups $\Gamma(D)=\{\gamma \in$ $\left.\mathrm{GL}\left(2, O_{\mathbb{A}}\right) ; \gamma \equiv I \bmod J_{D}\right\}$ make a basis of neighborhoods of the identity in $\mathrm{GL}(2, \mathbb{A})$. Thus we may assume in this proof that $U=\Gamma(D)$. In this case we shall show that $n=2 g-1-\operatorname{deg}(d)$. Indeed, fix $z \in \mathbb{A}$ and $x=\left(\begin{array}{ll}a & c \\ 0 & b\end{array}\right) k$ with $k \in \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$ and $\mathrm{ht}^{+}(x)=\operatorname{deg} a-\operatorname{deg} b \geq 2 g-1-\operatorname{deg}(d)$ (note: $\left.\operatorname{deg}(d)=-\operatorname{deg} D=\sum_{v} n_{v} \operatorname{deg} v\right)$. Then $\frac{a d}{b} O_{\mathbb{A}}+F=\mathbb{A}$ and $z=\frac{a d}{b} t+\beta$ for some $\beta \in F$ and $t \in O_{\mathbb{A}}$. Put $u=k^{-1}\left(\begin{array}{cc}1 & t d \\ 0 & 1\end{array}\right) k$. Then $u \in \Gamma(D)$ and $\left(\begin{array}{ll}1 & z \\ 0 & 1\end{array}\right) x=\left(\begin{array}{ll}1 & \beta \\ 0 & 1\end{array}\right) x u$.
We claim the proposition holds with $K=\{x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \operatorname{ht}(x)<$ $n\}$. This $K$ is compact modulo $Z(\mathbb{A})$. Let $\phi$ be a $U$-invariant cusp form, $x \in \mathrm{GL}(2, \mathbb{A}), \operatorname{ht}(x) \geq n$. We shall show that $\phi(x)=0$. Replacing $x$ by $\gamma x$ for suitable $\gamma \in \mathrm{GL}(2, F)$, we assume that $\mathrm{ht}^{+}(x) \geq n$. By our choice of $n, \phi\left(\left(\begin{array}{ll}1 & z \\ 0 & 1\end{array}\right) x\right)=\phi(x)$ for all $z$ in $\mathbb{A}$. Since $\phi$ is a cusp form, $\phi(x)=0$.

Corollary 2.14. The representation of $\mathrm{GL}(2, \mathbb{A})$ in $A_{0}^{E}(\chi)$ is admissible.
Proposition 2.15. Let $E^{\prime}$ be an extension of $E$, and $\chi: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times} a$ character. Then $A_{0}^{E^{\prime}}(\chi)=A_{0}^{E}(\chi) \otimes_{E} E^{\prime}$.

Proof. The space $A_{0}^{E}(\chi) \otimes_{E} E^{\prime}$ consists of the functions $\phi$ in $A_{0}^{E^{\prime}}(\chi)$ whose values span a finite dimensional space over $E$, since $\phi \in A_{0}^{E}(\chi)$ takes finite number of values times the set $\Gamma$ of values of $\chi$. But every $\phi$ in $A_{0}^{E^{\prime}}(\chi)$ has this property, since the set of its values lies in finitely many cosets of $\Gamma$.

Given a representation $\pi$ of $\operatorname{GL}(2, \mathbb{A})$ over $E$ and a character $\omega: \mathbb{A}^{\times} \rightarrow E^{\times}$, write $\omega \pi$ or $\pi \omega$ or $\omega \otimes \pi$ or $\pi \otimes \omega$ for the representation $(\pi \omega)(x)=\omega(\operatorname{det} x) \pi(x)$ in the space of $\pi$.

Proposition 2.16. For any characters $\chi, \omega: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$, we have $A_{0}^{E}(\chi) \otimes \omega=A_{0}^{E}\left(\chi \omega^{2}\right)$.

Proof. We need to construct an invertible linear map $L: A_{0}^{E}(\chi) \rightarrow A_{0}^{E}\left(\chi \omega^{2}\right)$ such that for every $\phi \in A_{0}^{E}(\chi)$ and $h \in \operatorname{GL}(2, \mathbb{A})$ we have $r(h) L(\phi)=$ $\omega(\operatorname{det} h) L(r(h) \phi)$, where $(r(h) \phi)(x)=\phi(x h)$.
Such $L$ is $(L \phi)(x)=\phi(x) \omega(\operatorname{det} x)$.
Proposition 2.17. Given a character $\chi: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$there exists a character $\omega: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$such that $\chi(x) \omega(x)^{2}$ is a root of unity for every $x$ in $\mathbb{A}^{\times} / F^{\times}$.

Proof. Fix $\alpha \in \mathbb{A}^{\times} / F^{\times}$with $\operatorname{deg} \alpha=1$. Such $\alpha$ exists since in the finite field extension $F / \mathbb{F}_{q}(t)$, where $t \in F$ is transcendental over $\mathbb{F}_{q}$, there are always primes which split completely. Fix $c$ in the algebraically closed field $E$ with $c^{2}=\chi(\alpha)$. Define $\omega: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$by $\omega(x)=c^{-\operatorname{deg}(x)}$, put $\chi_{1}(x)=$ $\chi(x) \omega^{2}(x)$, put $\alpha^{\mathbb{Z}}=\left\{\alpha^{n} ; n \in \mathbb{Z}\right\}$. Then $\chi_{1}$ is a character of the profinite group $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$, hence the values of $\chi_{1}$ are roots of 1 .

Proposition 2.18. Let $E$ be a subfield of $\mathbb{C}$ invariant under complex conjugation, $\chi$ an $E^{\times}$-valued unitary character of $\mathbb{A}^{\times} / F^{\times}$. Then the representation of $\mathrm{GL}(2, \mathbb{A})$ in $A_{0}^{E}(\chi)$ is unitary.
Proof. The function $x \mapsto \phi_{1}(x) \bar{\phi}_{2}(x)$ on $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$, where $\phi_{1}$, $\phi_{2} \in A_{0}^{E}(\chi)$, is invariant under $Z(\mathbb{A})$ and is compactly supported as a function on $\operatorname{PGL}(2, F) \backslash \operatorname{PGL}(2, \mathbb{A})$. Let $d x$ be an invariant measure on $\mathrm{PGL}(2, F) \backslash \mathrm{PGL}(2, \mathbb{A})$. It exists since $\mathrm{PGL}(2, F)$ is a discrete subgroup of $\operatorname{PGL}(2, \mathbb{A})$, a group with a two-sided invariant measure. Then

$$
\left(\phi_{1}, \phi_{2}\right)=\int \phi_{1}(x) \bar{\phi}_{2}(x) d x \quad(x \in \operatorname{PGL}(2, F) \backslash \operatorname{PGL}(2, \mathbb{A}))
$$

is an invariant scalar product on $A_{0}^{E}(\chi)$.
Corollary 2.19. The representation of $\mathrm{GL}(2, \mathbb{A})$ in $A_{0}^{E}(\chi)$ is a direct sum of irreducible subrepresentations.

Note that we may assume that all values of $\chi$ are roots of unity, and that $E=\overline{\mathbb{Q}}$.
The multiplicity one theorem asserts that in $A_{0}^{E}(\chi)$ any irreducible representation of $\mathrm{GL}(2, \mathbb{A})$ occurs with multiplicity one.
An irreducible representation of $\operatorname{GL}(2, \mathbb{A})$ over an algebraically closed field $E$ is called cuspidal if it is isomorphic to a subrepresentation of $A_{0}^{E}$.

### 2.6 Factorizability

Irreducible admissible representations of $\mathrm{GL}(2, \mathbb{A})$ are factorizable, as we proceed to show. Let $E$ denote an algebraically closed subfield of $\mathbb{C}$. An irreducible representation of $\operatorname{GL}\left(2, F_{v}\right)$ in an $E$-space $V$ is unramified if $V$ contains a nonzero $\mathrm{GL}\left(2, O_{v}\right)$-invariant vector.

Proposition 2.20. The space of $\mathrm{GL}\left(2, O_{v}\right)$-invariant vectors $V^{\mathrm{GL}\left(2, O_{v}\right)}$ in an unramified representation $(\pi, V)$ of $\mathrm{GL}\left(2, F_{v}\right)$ is one dimensional.

Proof. Denote by $H_{v}=C_{c}\left(\mathrm{GL}\left(2, O_{v}\right) \backslash \mathrm{GL}\left(2, F_{v}\right) / \mathrm{GL}\left(2, O_{v}\right)\right)$ the Hecke convolution algebra of compactly supported $\mathrm{GL}\left(2, O_{v}\right)$-biinvariant $E$-valued functions on $\mathrm{GL}\left(2, F_{v}\right)$. We claim it is a commutative algebra. Indeed, for any $f \in H_{v}$, the function ${ }^{t} f(x)=f\left({ }^{t} x\right)$, where ${ }^{t} x$ is the transpose of $x$, is also in $H_{v}$. Since ${ }^{t}(x y)={ }^{t} y^{t} x$, we have ${ }^{t}\left(f_{1} * f_{2}\right)={ }^{t} f_{2} *^{t} f_{1}$ for all $f_{1}, f_{2} \in H_{v}$. By Cartan decomposition every $\mathrm{GL}\left(2, O_{v}\right)$-double coset in $\mathrm{GL}\left(2, F_{v}\right)$ contains a diagonal
matrix. Hence ${ }^{t} f=f$ for all $f \in H_{v}$, and $f_{1} * f_{2}={ }^{t}\left(f_{1} * f_{2}\right)={ }^{t} f_{2} *{ }^{t} f_{1}=f_{2} * f_{1}$ for all $f_{1}, f_{2} \in H_{v}$. If $V$ is unramified, $V^{\mathrm{GL}\left(2, O_{v}\right)}$ is a nonzero irreducible $H_{v^{-}}$ module. But $H_{v}$ is commutative, so $\operatorname{dim}_{E} V^{\mathrm{GL}\left(2, O_{v}\right)}$ is 1 .

Given an irreducible admissible representation $\pi_{v}$ of $\mathrm{GL}\left(2, F_{v}\right)$ in a space $V_{v}$ for every closed point $v \in|X|$ such that $\pi_{v}$ is unramified for all $v \in S, S \subset|X|$ finite, construct a representation $\pi=\otimes \pi_{v}$ of $\mathrm{GL}(2, \mathbb{A})$ as follows. For each $v \in|X|-S$ choose a nonzero vector $\xi_{v}^{0} \in V_{v}^{\mathrm{GL}\left(2, O_{v}\right)}$. For any finite set $S^{\prime} \supset S$ of closed points of $X$ put $V_{S^{\prime}}=\otimes_{v \in S^{\prime}} V_{v}$. If $S^{\prime \prime} \supset S^{\prime} \supset S$, define an inclusion $V_{S^{\prime}} \hookrightarrow V_{S^{\prime \prime}}$ by $x \mapsto\left(\otimes_{v \in S^{\prime \prime}-S^{\prime}} \xi_{v}^{0}\right) \otimes x$. Put $V=\lim _{\overrightarrow{S^{\prime} S} S} V_{S^{\prime}}$. It is the span of the vectors $\otimes_{v \in|X|} \xi_{v}, \xi_{v}=\xi_{v}^{0}$ for almost all $v$, and $\xi_{v} \in V_{v}$ for all $v \in|X|$. Then $V$ is a $\mathrm{GL}(2, \mathbb{A})$-module in a natural way; denote by $\pi$ the corresponding representation of $\operatorname{GL}(2, \mathbb{A})$. The vectors $\xi_{v}^{0}$ are determined uniquely up to a scalar multiple, hence $\pi$ is uniquely determined by the $\pi_{v}$ for all $v \in|X|$. Reducing to irreducible finite dimensional representations of tensor products of algebras, we have

Proposition 2.21. Given an irreducible admissible representation $\pi_{v}$ of $\mathrm{GL}\left(2, F_{v}\right)$ for every $v$ in $|X|$ which is unramified for almost all $v, \pi=\otimes_{v} \pi_{v}$ is an irreducible admissible representation of $\mathrm{GL}(2, \mathbb{A})$. Every irreducible admissible representation $\pi$ of $\mathrm{GL}(2, \mathbb{A})$ equals $\otimes_{v} \pi_{v}$ for some irreducible admissible representations $\pi_{v}$ of $\mathrm{GL}\left(2, F_{v}\right)$ which are almost all unramified. The representations $\pi_{v}$ are determined by $\pi$ uniquely up to isomorphism.

## 3 Looking For a trace formula

### 3.1 Trace formula in the compact case

Let $X$ be an $\ell$-space. Denote by $C^{\infty}(X)$ the space of $E$-valued locally constant (= smooth) functions on $X$. Here $E$ is a fixed algebraically closed subfield of $\mathbb{C}$. Let $C_{c}^{\infty}(X)$ be the space of smooth compactly supported $E$-valued functions on $X$. Let $r$ be an admissible representation of an $\ell$-group $G$ in an $E$-space $V$. Fix a Haar measure $d x$ on $G$. Given $f \in C_{c}^{\infty}(G)$, define $r(f)=\int_{G} f(x) r(x) d x$, an endomorphism of $V$. Since $f$ is $C^{\infty}$, that is smooth, it is right invariant under an open subgroup $U$ of $G$. Then $\operatorname{Im} r(f) \subset V^{U}$, so $\operatorname{Im} r(f)$ is finite dimensional, and the trace $\operatorname{tr} r(f)$ is well defined. Let $r$ be now the representation of $G$ on $C^{\infty}(\Gamma \backslash G)$ by right translation, where $\Gamma$ is a discrete cocompact subgroup of $G$. Since $r$ is admissible, $\operatorname{tr} r(f)$ is defined.

Proposition 3.1. Let $G$ be an $\ell$-group. Let $\Gamma$ be a discrete cocompact sugroup of $G$. Then $G$ has a two sided invariant measure and $\Gamma \backslash G$ has a $G$-invariant measure.

Proof. Since (see [BZ76]) $\Gamma \backslash G$ admits a measure which when translated by $x$ in $G$ is multiplied by $\Delta(x)$, where $\Delta$ is the modulus of $G$, we have $|\Gamma \backslash G|=$ $\Delta(x)|\Gamma \backslash G|$, thus $\Delta=1$.

Proposition 3.2. Let $X$ be an $\ell$-space, $d x$ a measure on $X, K \in C_{c}^{\infty}(X \times X)$. Define a linear endomorphism $A$ of $C^{\infty}(X)$ by $(A \phi)(y)=\int_{X} K(x, y) \phi(x) d x$. Then the image of $A$ is finite dimensional and $\operatorname{tr} A=\int_{X} K(x, x) d x$.
Proof. We may assume that $K(x, y)$ is of the form $\varphi(x) \psi(y)$, as such functions span $C_{c}^{\infty}(X \times X)$. In this case the claim is clear.

Proposition 3.3. Let $G$ be an $\ell$-group, $\Gamma$ a discrete cocompact subgroup, $r$ the representation of $G$ in $C^{\infty}(\Gamma \backslash G)$ by right translation, dx a Haar measure on $G, f \in C_{c}^{\infty}(G), S$ a set of representatives of the conjugacy classes in $\Gamma, Z_{\Gamma}(\gamma)$ the centralizer of $\gamma$ in $\Gamma$. Then $\operatorname{tr} r(f)=\sum_{\gamma \in S} \int_{G / Z_{\Gamma}(\gamma)} f\left(x \gamma x^{-1}\right) d x$.

Proof. We first show that for each $\gamma \in \Gamma$ the function $x \mapsto f\left(x \gamma x^{-1}\right)$ on $G / Z_{\Gamma}(\gamma)$ is compactly supported, and that there are at most finitely many $\gamma \in S$ for which $x \mapsto f\left(x \gamma x^{-1}\right)$ is not identically zero. For this, fix a compact subset $K$ in $G$ with $K \Gamma=G$. Given $x \in G$ there are $k \in K, \delta \in \Gamma$, with $x=k \delta$. Fix $\gamma \in \Gamma$. If $f\left(x \gamma x^{-1}\right) \neq 0$ then $k \delta \gamma \delta^{-1} k^{-1}$ lies in $\operatorname{supp} f$, thus $\delta \gamma \delta^{-1} \in K_{f}=K \cdot \operatorname{supp} f \cdot K$. Since $K_{f}$ is compact $K_{f} \cap \Gamma$ is finite, and there are only finite number of possibilities for $\delta \gamma \delta^{-1}$. Hence there are only a finite number of possibilities $\delta_{1}, \ldots, \delta_{n}$ for $\delta$ modulo $Z_{\Gamma}(\gamma)$. Then $f\left(x \gamma x^{-1}\right) \neq 0$ implies that $x \in K^{\prime} Z_{\Gamma}(\gamma)$, where $K^{\prime}=\cup_{1 \leq i \leq n} K \delta_{i}$ is compact. If $f\left(x \gamma x^{-1}\right) \neq$ 0 , the conjugacy class of $\gamma$ in $\Gamma$ intersects the finite set $K_{f} \cap \Gamma$. The number of such classes is finite. Thus the sum is finite and the integrals converge.
Now given $\phi$ in $C^{\infty}(\Gamma \backslash G)$, for any $y$ in $G$ we have

$$
(r(f) \phi)(y)=\int_{G} f(x) \phi(y x) d x=\int_{G} f\left(y^{-1} x\right) \phi(x) d x=\int_{\Gamma \backslash G} K_{f}(x, y) \phi(x) d x
$$

where $K_{f}(x, y)=\sum_{\gamma \in \Gamma} f\left(y^{-1} \gamma x\right)$. Then

$$
\begin{gathered}
\operatorname{tr} r(f)=\int_{\Gamma \backslash G} K_{f}(x, x) d x=\int_{\Gamma \backslash G} \sum_{\gamma \in \Gamma} f\left(x^{-1} \gamma x\right) d x \\
=\int_{\Gamma \backslash G} \sum_{\gamma \in S} \sum_{\delta \in Z_{\Gamma}(\gamma) \backslash \Gamma} f\left(x^{-1} \delta^{-1} \gamma \delta x\right) d x=\sum_{\gamma \in S} \int_{\Gamma \backslash G} \sum_{\delta \in Z_{\Gamma}(\gamma) \backslash \Gamma} f\left(x^{-1} \delta^{-1} \gamma \delta x\right) d x \\
=\sum_{\gamma \in S} \int_{Z_{\Gamma}(\gamma) \backslash G} f\left(x^{-1} \gamma x\right) d x
\end{gathered}
$$

### 3.2 Case of GL(2), oversimplified

Let now $A_{0}^{E}$ denote the space of $E$-valued cusp forms on $\mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$. The right-shifts representation of $\mathrm{GL}(2, \mathbb{A})$ on $A_{0}^{E}$ is not admissible since the center $Z(\mathbb{A})$ of $\mathrm{GL}(2, \mathbb{A})$ is not compact. Fix a degree-one idèle $\alpha$ and put $\alpha^{\mathbb{Z}}=\left\{\alpha^{n} ; n \in \mathbb{Z}\right\}$. It is a cyclic subgroup of $\mathbb{A}^{\times}$, and we view $\mathbb{A}^{\times}$as the
center of GL $(2, \mathbb{A})$. Denote by $A_{0, \alpha}^{E}$ the space of cusp forms in $A_{0}^{E}$ invariant under $\alpha$, and by $r_{0}$ the representation of $\mathrm{GL}(2, \mathbb{A})$ on $A_{0, \alpha}^{E}$ by right translation. Since $\mathbb{A}^{\times} / F^{\times} \alpha^{\mathbb{Z}}$ is compact and every $U$-invariant cusp form - where $U$ is an open subgroup of $G L(2, \mathbb{A})$ - is supported on some compact modulo $Z(\mathbb{A})$ set $K \subset \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$, the representation $r_{0}$ is admissible. Hence $\operatorname{tr} r_{0}(f)$ is defined for every $f \in C_{c}^{\infty}(\operatorname{GL}(2, \mathbb{A}))$.
Put $A_{c, \alpha}=C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right)$. Fix $f \in C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$. Let $r$ be the right representation of $\operatorname{GL}(2, \mathbb{A})$ on $A_{c, \alpha}$. We proceed to compute $\operatorname{tr} r(f)$ as if the space $\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ were compact, to see what needs to be corrected. This space is not compact and $r$ is not admissible, so that in fact $\operatorname{tr} r(f)$ makes no sense.
For any ring $R$ define $A(R)=\left\{\operatorname{diag}(a, b) ; a, b \in R^{\times}\right\}, A^{\prime}(R)=\{\operatorname{diag}(a, b) ; a, b \in$ $\left.R^{\times}, a \neq b\right\}, N(R)=\left\{\left(\begin{array}{cc}1 & a \\ 0 & 1\end{array}\right) ; a \in R\right\}$. Let $Q$ be the set of quadratic extensions of the field $F$. For each $L \in Q$ choose an embedding $L \hookrightarrow M(2, F)$; it exists and is unique up to an automorphism of $M(2, F)$; all automorphisms of $M(2, F)$ are inner. Given $\gamma \in \alpha^{\mathbb{Z}} \cdot \operatorname{GL}(2, F)$, denote by $Z(\gamma)$ the centralizer of $\gamma$ in $\alpha^{\mathbb{Z}} \mathrm{GL}(2, F)$.

Proposition 3.4. Every conjugacy class of $\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F)$ intersects precisely one of : $F^{\times} \cdot \alpha^{\mathbb{Z}} ; a\left(\begin{array}{ll}1 & 1 \\ 0 & 1\end{array}\right), a \in F^{\times} \cdot \alpha^{\mathbb{Z}} ; \alpha^{\mathbb{Z}} \cdot A^{\prime}(F) ; \alpha^{\mathbb{Z}} \cdot\left(L^{\times}-F^{\times}\right)$for some $L \in Q$. In the first two cases the number of intersection points is 1 , in the 3 rd case 2 , in the 4 th case: the number of automorphisms of $L$ over $F$. The centralizers $Z(\gamma)$ are $\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F), \alpha^{\mathbb{Z}} F^{\times} N(F), \alpha^{\mathbb{Z}} \cdot A(F), \alpha^{\mathbb{Z}} L^{\times}$, respectively.

Immitating the trace formula in the compact case, one may expect

$$
\operatorname{tr} r(f)=S_{1}(f)+\sum_{L \in Q} S_{2, L}(f)+S_{3}(f)+S_{4}(f)
$$

with

$$
\begin{gathered}
S_{1}(f)=\left|\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right|, \\
S_{2, L}(f)=\left|\operatorname{Aut}_{F}(L)\right|^{-1} \sum_{\gamma \in \alpha^{\mathbb{Z}} \cdot\left(L^{\times}-F^{\times}\right)} \int_{\alpha^{\mathbb{Z}} \cdot L^{\times} \backslash \mathrm{GL}(2, \mathbb{A})} f\left(x^{-1} \gamma x\right) d x, \\
S_{3}(f)=\frac{1}{2} \sum_{\gamma \in \alpha^{\mathbb{Z}} A^{\prime}(F)} \int_{\alpha^{\mathbb{Z}} A(F) \backslash \mathrm{GL}(2, \mathbb{A})} f\left(x^{-1} \gamma x\right) d x, \\
S_{4}(f)=\sum_{a \in \alpha^{\mathbb{Z}} \cdot F^{\times}} \int_{\alpha^{\mathbb{Z}} F^{\times} \times N(F) \backslash \operatorname{GL}(2, \mathbb{A})} f\left(x^{-1} a\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right) x\right) d x .
\end{gathered}
$$

The left side of this wrong trace formula is divergent. So is $S_{3}(f)$, since the homogeneous space $A(\mathbb{A}) / \alpha^{\mathbb{Z}} \cdot A(F)$ is not compact. We shall show that $S_{1}(f)$ and $\sum_{L \in Q} S_{2, L}(f)$ converge, and although $S_{4}(f)$ diverges, we shall show in which way it does.

Proposition 3.5. Given $f \in C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$, the number of conjugacy classes of $\gamma \in \alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F)$ with $x \in \mathrm{GL}(2, \mathbb{A})$ and $f\left(x \gamma x^{-1}\right) \neq 0$ is finite.

Proof. The sets $K_{1}=\{\operatorname{tr} h ; h \in \operatorname{supp} f\} \subset \mathbb{A}, K_{2}=\{\operatorname{det} h ; h \in \operatorname{supp} f\} \subset$ $\mathbb{A}^{\times}$are compact. It suffices to show that the set $\left\{\gamma \in \alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) ; \operatorname{tr} \gamma \in\right.$ $K_{1}$, $\left.\operatorname{det} \gamma \in K_{2}\right\}$ is a union of finitely many conjugacy classes. Put $\gamma=\alpha^{n} x$ for some $x \in \operatorname{GL}(2, F)$. Then $2 n=\operatorname{deg} \gamma$, so $n$ lies in a finite set. Fix $n$. Then $\operatorname{tr} x \in \alpha^{-n} K_{1}$, $\operatorname{det} x \in \alpha^{-2 n} K_{2}$. But the sets $F \cap \alpha^{-n} K_{1}$ and $F^{\times} \cap \alpha^{-2 n} K_{2}$ are finite. Hence the trace and determinant of $x$ can take only finitely many values. As the number of conjugacy classes of elements in $\mathrm{GL}(2, F)$ with fixed trace and determinant is at most two, we are done.

### 3.3 Central elements

Proposition 3.6. The volume $\left|\mathrm{GL}(2, F) \cdot \alpha^{\mathbb{Z}} \backslash \mathrm{GL}(2, \mathbb{A})\right|$ is finite.
Proof. This volume is equal to (below $x \in \alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) / \mathrm{GL}\left(2, O_{\mathbb{A}}\right)$ )

$$
\begin{aligned}
& \sum_{x}\left|\alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \cap x \mathrm{GL}\left(2, O_{\mathbb{A}}\right) x^{-1} \backslash x \mathrm{GL}\left(2, O_{\mathbb{A}}\right)\right| \\
= & \left|\mathrm{GL}\left(2, O_{\mathbb{A}}\right)\right| \sum_{x}\left|\alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \cap x \mathrm{GL}\left(2, O_{\mathbb{A}}\right) x^{-1}\right|^{-1} .
\end{aligned}
$$

For $x$ in $\operatorname{GL}(2, \mathbb{A}) / \operatorname{GL}\left(2, O_{\mathbb{A}}\right)$, let $\mathcal{E}=x O_{\mathbb{A}}^{2}$ be the associated rank 2 locally free sheaf on $X$. Then $\operatorname{Aut}(\mathcal{E})$ consists of the $g \in \mathrm{GL}(2, \mathbb{A})$ which map $(\mathcal{E}=) x O_{\mathbb{A}}^{2}$ to $x O_{\mathbb{A}}^{2}$ and the generic fiber $F^{2}$ to itself, thus Aut $\mathcal{E}$ is

$$
\mathrm{GL}(2, F) \cap x \mathrm{GL}\left(2, O_{\mathbb{A}}\right) x^{-1}=\alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \cap x \mathrm{GL}\left(2, O_{\mathbb{A}}\right) x^{-1}
$$

We then need to show the convergence of

$$
\sum_{\mathcal{E} \in \operatorname{Bun}_{2} / J}|\operatorname{Aut} \mathcal{E}|^{-1},
$$

$J$ being the image of $\alpha^{\mathbb{Z}}$ under the natural homomorphism $\mathbb{A}^{\times} \rightarrow \operatorname{Pic} X$. The number of $J$-orbits on the set of stable rank two locally free sheaves on $X$ is finite, so it remains to show that the sum of $\mid$ Aut $\left.\mathcal{E}\right|^{-1}$ over the set Bun $_{2}^{\text {un }}$ of $J$-orbits of unstable rank two locally free sheaves on $X$ is convergent.

Lemma 3.7. (1) A rank two locally free sheaf $\mathcal{E}$ on $X$ is very unstable $(\operatorname{ht}(\mathcal{E}) \geq$ $2 g-1)$ iff $\mathcal{E} \simeq \mathcal{L} \oplus \mathcal{M}$ where $\mathcal{L}, \mathcal{M}$ are invertible sheaves with $\operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{M} \geq$ $2 g-1$.
(2) If $\mathcal{L}, \mathcal{M} \in \operatorname{Pic} X$ and $\operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{M} \geq \max (2 g-1,1)$ then

$$
|\operatorname{Aut}(\mathcal{L} \oplus \mathcal{M})|=(q-1)^{2} q^{\operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{M}+1-g}
$$

(3) If $\mathcal{L} \oplus \mathcal{M} \simeq \mathcal{L}^{\prime} \oplus \mathcal{M}^{\prime}$ with $\operatorname{deg} \mathcal{L}>\operatorname{deg} \mathcal{M}, \operatorname{deg} \mathcal{L}^{\prime}>\operatorname{deg} \mathcal{M}^{\prime}$ then $\mathcal{L} \simeq$ $\mathcal{L}^{\prime}, \mathcal{M} \simeq \mathcal{M}^{\prime}$.

Proof. (1) If $\mathcal{L}$ is an invertible sheaf of $\mathcal{E}$ of maximal degree and $\mathcal{M}=\mathcal{E} / \mathcal{L}$, then $\mathcal{M}$ is invertible, and $\operatorname{Ext}(\mathcal{M}, \mathcal{L})=H^{1}\left(X, \mathcal{M}^{-1} \mathcal{L}\right)$ is 0 (by Serre duality) as

$$
\operatorname{deg} \mathcal{M}^{-1} \mathcal{L}=\operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{M}=2 \operatorname{deg} \mathcal{L}-\operatorname{deg} \mathcal{E}=\operatorname{ht}(\mathcal{E}) \geq 2 g-1
$$

The exact sequence

$$
0 \rightarrow \operatorname{Hom}(\mathcal{M}, \mathcal{L}) \rightarrow \operatorname{Aut}(\mathcal{L} \oplus \mathcal{M}) \rightarrow \operatorname{Aut} \mathcal{L} \times \operatorname{Aut} \mathcal{M} \rightarrow 0
$$

implies (2) since $\operatorname{Hom}(\mathcal{M}, \mathcal{L})=H^{0}\left(X, \mathcal{M}^{-1} \mathcal{L}\right)$ and $H^{1}\left(X, \mathcal{M}^{-1} \mathcal{L}\right)=\{0\}$, so the Riemann-Roch theorem implies that $\operatorname{dim} H^{0}\left(X, \mathcal{M}^{-1} \mathcal{L}\right)=\operatorname{deg}\left(\mathcal{M}^{-1} \mathcal{L}\right)+$ $1-g$. If the invertible sheaf $\mathcal{L}$ corrsponds to $a O_{\mathbb{A}}$, then Aut $\mathcal{L}$ consists of $g \in \mathbb{A}^{\times}$ which map the generic fiber $F$ onto itself (thus $g \in F^{\times}$) and map $a O_{\mathbb{A}}$ onto itself (thus $g \in O_{\mathbb{A}}^{\times}$). Then Aut $\mathcal{L}=F^{\times} \cap O_{\mathbb{A}}^{\times}=\mathbb{F}_{q}^{\times}$has cardinality $q-1$.
For (3), put $\mathcal{E}=\mathcal{L} \oplus \mathcal{M} \xrightarrow[\rightarrow]{\sim} \mathcal{L}^{\prime} \oplus \mathcal{M}^{\prime}$. Since $\operatorname{deg} \mathcal{L}>(\operatorname{deg} \mathcal{E}) / 2>\operatorname{deg} \mathcal{M}^{\prime}$, we have $\operatorname{Hom}\left(\mathcal{L}, \mathcal{M}^{\prime}\right)=\{0\}$. Hence the image of $\mathcal{L}$ under the isomorphism $\mathcal{L} \oplus \mathcal{M} \xrightarrow{\sim} \mathcal{L}^{\prime} \oplus \mathcal{M}^{\prime}$ lies in $\mathcal{L}^{\prime}$. Hence $\mathcal{L} \simeq \mathcal{L}^{\prime}$ and $\mathcal{M} \simeq \mathcal{E} / \mathcal{L} \simeq \mathcal{E} / \mathcal{L}^{\prime} \simeq \mathcal{M}^{\prime}$.

Assume $g \geq 1$, so that $2 g-1 \geq 1$ (the case $g=0$ is similar). The lemma implies

$$
\sum_{\mathcal{E} \in \operatorname{Bun}_{2}^{\mathrm{un}} / J}|\operatorname{Aut} \mathcal{E}|^{-1}=(q-1)^{-2}\left|\operatorname{Pic}^{0}(X)\right| \sum_{n \geq 2 g-1} q^{g-1-n}<\infty .
$$

Corollary 3.8. If the Haar measure on $\mathrm{GL}(2, \mathbb{A})$ is normalized so that $\left|\mathrm{GL}\left(2, O_{\mathbb{A}}\right)\right|$ is a rational number, then $\left|\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right| \in \mathbb{Q}$.

This follows from the proof of the last proposition.

### 3.4 Elliptic ElEments

Proposition 3.9. Let $L$ be a quadratic extension of $F, \gamma \in \alpha^{\mathbb{Z}} \cdot\left(L^{\times}-F^{\times}\right) \subset$ $\mathrm{GL}(2, \mathbb{A})$, and $f \in C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$. Then the function $x \mapsto f\left(x \gamma x^{-1}\right)$ on $\mathrm{GL}(2, \mathbb{A}) / \alpha^{\mathbb{Z}} \cdot L^{\times}$has compact support.

Proof. We need to show that the map $x \mapsto x \gamma x^{-1}$ on $\mathrm{GL}(2, \mathbb{A}) / \alpha^{\mathbb{Z}} \cdot L^{\times}$is proper (the preimage of a compact is compact). Since $\left(L \otimes_{F} \mathbb{A}\right)^{\times} / \alpha^{\mathbb{Z}} \cdot L^{\times}$is compact, it suffices to show that the map $\psi(x)=x \gamma x^{-1}, \psi: \operatorname{GL}(2, \mathbb{A}) / \mathbb{A}_{L}^{\times} \rightarrow \operatorname{GL}(2, \mathbb{A})$, is proper $\left(\mathbb{A}_{L}=L \otimes_{F} \mathbb{A}\right.$ is the ring of adèles of $L$ ).

Lemma 3.10. Let $F$ be a local field in this lemma. Suppose $\gamma \in M(2, F)$ is regular, i.e. the subalgebra $E=F[\gamma]$ generated by $\gamma$ is a field or is $F \times F$. Then the map $\psi: \mathrm{GL}(2, F) / E^{\times} \rightarrow \mathrm{GL}(2, F), x \mapsto x \gamma x^{-1}$, is proper. Moreover, if $\gamma \in \mathrm{GL}(2, O)$ and the ring $O[\gamma]$ is integrally closed, then $\psi^{-1}(\mathrm{GL}(2, O))=$ $\mathrm{GL}(2, O) / E^{\times} \cap \mathrm{GL}(2, O)$.

Proof. The conjugacy class $C$ of $\gamma$ is a closed subset of $\mathrm{GL}(2, F)$, since $\gamma$ is regular. So it suffices to show that $\psi$ maps $\mathrm{GL}(2, F) / E^{\times}$homeomorphically onto $C$. It is clear that $\psi$ is continuous, injective and $\operatorname{Im} \psi=C$. It remains to show that the map $\psi^{\prime}: \mathrm{GL}(2, F) \rightarrow C, x \mapsto x \gamma x^{-1}$, is open. For this, it suffices to show that $C$ is the set of $F$-points of a smooth variety $\mathbf{C}$ over $F$, and that $\psi^{\prime}$
is smooth, that is its differential is everywhere onto. Since $\mathbf{C}$ is a homogeneous space under a connected group $\mathbf{G}$ is suffices to show that the tangent map $d \psi^{\prime}$ of $\psi^{\prime}$ at the identity is onto. When verifying these properties of $\mathbf{C}$ and $\psi^{\prime}$, we may replace $F$ with an extension, thus we may assume that $\gamma$ is of the form $\operatorname{diag}(a, b)$ with $a \neq b$, or $\left(\begin{array}{cc}a & 1 \\ 0 & a\end{array}\right)$ (if $E$ is nonseparable over $F$ ). To compute the tangent map $d \psi^{\prime}$ : Lie $G \rightarrow T_{\gamma}(\mathbf{C})$ of $\psi^{\prime}(x)=x \gamma x^{-1}$ near the identity $x=1$, let $Y$ be in Lie $G$, and put $x=1+\epsilon Y$, where $\epsilon^{2}=0$. Then $x^{-1}=1-\epsilon Y$ and $\psi^{\prime}(x)=(1+\epsilon Y) \gamma(1-\epsilon Y)=1+\epsilon(Y \gamma-\gamma Y)$, so $d \psi^{\prime}(Y)=Y \gamma-\gamma Y$ is onto the tangent space $T_{\gamma}(\mathbf{C})$ of $\mathbf{C}$ at $\gamma$, and $\psi$ is proper.
If $x \in \mathrm{GL}(2, F)$ and $x \gamma x^{-1} \in \mathrm{GL}(2, O)$, put $M=x^{-1} O^{2}$. Then $\gamma M \subset M$. In addition, $\gamma \in \mathrm{GL}(2, O)$, so $\gamma O^{2} \subset O^{2}$. Thus $M$ and $O^{2}$ are $O[\gamma]$-submodules in $F^{2}$. Both modules are of finite type. As $F^{2}$ is a rank one free $E=F[\gamma]$-module, and we assume that $O[\gamma]$ is integrally closed, namely it is the ring of integers in $E=F[\gamma]$, both $M$ and $O^{2}$ are rank one torsion free over the discrete valuation ring $O[\gamma]$ (being rank two over $O$ ). Hence there exists $a \in E^{\times}$with $M=a O^{2}$. Thus $x a O^{2}=O^{2}$, that is $x a \in \mathrm{GL}(2, O)$.

Now for $\gamma$ as in the proposition, for almost all closed points in $X$ the component of $\alpha$ at $v$ is $1, \gamma \in \mathrm{GL}\left(2, O_{v}\right)$, and the ring $O_{v}[\gamma]$ is integrally closed. This and the lemma imply the proposition.

### 3.5 Regularization of the unipotent terms

To study the integral which occurs in $S_{4}(f)$, we regularize it as

$$
\theta_{a, f}(t)=\int_{\alpha^{\mathbb{Z}} \cdot F \times N(F) \backslash \mathrm{GL}(2, F)} f\left(a x^{-1}\left(\begin{array}{cc}
1 & 1 \\
0 & 1
\end{array}\right) x\right) t^{\mathrm{ht}^{+}(x)} d x
$$

Proposition 3.11. (1) For every $f \in C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$ and $a \in \mathbb{A}^{\times}$, the integral $\theta_{a, f}(t)$ converges as an element of $\mathbb{C}((t))$, and $\zeta_{F}\left(q^{-1} t\right)^{-1} \theta_{a, f}(t) \in \mathbb{C}\left[t, t^{-1}\right]$, where $\zeta_{F}(t)=\prod_{v \in|X|}\left(1-t_{v}\right)^{-1}, t_{v}=t^{\operatorname{deg} v}$.
(2) If $f$ is the characteristic function of $\mathrm{GL}\left(2, O_{\mathbb{A}}\right)$ in $\mathrm{GL}(2, \mathbb{A})$, then

$$
\theta_{1, f}(t)=\left|\mathrm{GL}\left(2, O_{\mathbb{A}}\right)\right| \cdot(q-1)^{-1} q^{g-1} \cdot\left|\operatorname{Pic}^{0}(X)\right| \zeta_{F}\left(q^{-1} t\right)
$$

Proof. (1) It suffices to consider $f(x)=\prod_{v} f_{v}\left(x_{v}\right), x=\left(x_{v}\right) \in \mathrm{GL}(2, \mathbb{A})$, where $f_{v} \in C_{c}^{\infty}\left(\mathrm{GL}\left(2, F_{v}\right)\right)$ for all $v \in|X|$ and $f_{v}$ is the characteristic function $f_{v}^{0}$ of $\mathrm{GL}\left(2, O_{v}\right)$ at almost all $v$, since such functions span $C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$. Normalize the measures on $F_{v}^{\times}$and $F_{v}$ so that $\left|O_{v}^{\times}\right|=1=\left|O_{v}\right|$. Denote by val $\left(x_{v}\right)$ the valuation of $x_{v} \in F_{v}^{\times}$, normalized by $\operatorname{val}_{v}\left(\boldsymbol{\pi}_{v}\right)=1$. Define a function

$$
h_{v}^{+}: \mathrm{GL}\left(2, F_{v}\right) \rightarrow \mathbb{Z} \text { by } h_{v}^{+}\left(\left(\begin{array}{cc}
a & b \\
0 & c
\end{array}\right) k\right)=\operatorname{val}_{v}(a)-\operatorname{val}_{v}(c), \quad k \in \mathrm{GL}\left(2, O_{v}\right) .
$$

Then $h_{v}^{+}$is well-defined and $\mathrm{ht}^{+}(x)=\sum_{v \in|X|} h_{v}^{+}\left(x_{v}\right) \operatorname{deg}(v)$. We have

$$
\theta_{a, f}(t)=\left|\mathbb{A}^{\times} / \alpha^{\mathbb{Z}} \cdot F^{\times}\right| \cdot|\mathbb{A} / F| \prod_{v} \theta_{a_{v}, f_{v}}\left(t_{v}\right)
$$

where

$$
\theta_{a_{v}, f_{v}}\left(t_{v}\right)=\int_{F_{v}^{\times} N\left(F_{v}\right) \backslash \operatorname{GL}\left(2, F_{v}\right)} f_{v}\left(a_{v} x^{-1}\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right) x\right) t^{h_{v}^{+}(x) \operatorname{deg} v} d x
$$

and $t_{v}=t^{\operatorname{deg}(v)}$. To compute it, note that $p_{n, v}=\operatorname{diag}\left(\boldsymbol{\pi}_{v}^{n}, 1\right)(n \in \mathbb{Z})$ make a set of representatives of the two sided coset space

$$
F_{v}^{\times} N\left(F_{v}\right) \backslash \mathrm{GL}\left(2, F_{v}\right) / \mathrm{GL}\left(2, O_{v}\right) .
$$

Then

$$
\begin{aligned}
& \theta_{a_{v}, f_{v}}\left(t_{v}\right)=\sum_{n \in \mathbb{Z}} t_{v}^{n} \int_{F_{v}^{\times} N\left(F_{v}\right) \cap p_{n, v}^{-1} \mathrm{GL}\left(2, O_{v}\right) p_{n, v} \backslash p_{n, v}^{-1} \mathrm{GL}\left(2, O_{v}\right)} f_{v}\left(a_{v} x^{-1}\left(\begin{array}{cc}
1 & 1 \\
0 & 1
\end{array}\right) x\right) d x \\
& =\sum_{n \in \mathbb{Z}} t_{v}^{n}\left|F_{v}^{\times} N\left(F_{v}\right) \cap p_{n, v}^{-1} \mathrm{GL}\left(2, O_{v}\right) p_{n, v}\right|^{-1} \int_{p_{n, v}^{-1} \mathrm{GL}\left(2, O_{v}\right)} f_{v}\left(a_{v} x^{-1}\left(\begin{array}{cc}
1 & 1 \\
0 & 1
\end{array}\right) x\right) d x \\
& =\sum_{n \in \mathbb{Z}} q_{v}^{-n} t_{v}^{n} \int_{\mathrm{GL}\left(2, O_{v}\right)} f_{v}\left(a_{v} y p_{n, v}\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right) p_{n, v}^{-1} y^{-1}\right) d y=\sum_{n \in \mathbb{Z}} \tau_{n}\left(f_{v}\right) q_{v}^{-n} t_{v}^{n},
\end{aligned}
$$

where $\tau_{n}\left(f_{v}\right)=\int_{\mathrm{GL}\left(2, O_{v}\right)} f_{v}\left(a_{v} y\left(\begin{array}{cc}1 & \pi_{v}^{n} \\ 0 & 1\end{array}\right) y^{-1}\right) d y$ is 0 if $n \ll 0$ and $\tau_{n}\left(f_{v}\right)=$ $f_{v}\left(a_{v}\right)$ for $n \gg 0$.
If $a_{v} \in O_{v}^{\times}$and $f_{v}$ is the characteristic function of GL $\left(2, O_{v}\right)$, then $\tau_{n}\left(f_{v}\right)=$ $\left|\mathrm{GL}\left(2, O_{v}\right)\right|$ for $n \geq 0$ and $u_{n, v}=0$ for $n<0$, so

$$
\theta_{a_{v}, f_{v}}\left(t_{v}\right)=\left|\mathrm{GL}\left(2, O_{v}\right)\right|\left(1-t_{v} / q_{v}\right)^{-1}
$$

(2) It remains to compute (note that $\left|O_{\mathbb{A}}^{\times}\right|=1$ and $\left|O_{\mathbb{A}}\right|=1$ ):

$$
\left|\mathbb{A}^{\times} N(\mathbb{A}) / \alpha^{\mathbb{Z}} F^{\times} N(F)\right|=\left(\left|\mathbb{A}^{\times} / \alpha^{\mathbb{Z}} F^{\times}\right| /\left|O_{\mathbb{A}}^{\times}\right|\right)\left(|\mathbb{A} / F| /\left|O_{\mathbb{A}}\right|\right) .
$$

The exact sequence $1 \rightarrow \mathbb{F}_{q}^{\times} \rightarrow O_{\mathbb{A}}^{\times} \rightarrow \mathbb{A}^{\times} / \alpha^{\mathbb{Z}} F^{\times} \rightarrow \operatorname{Pic} X / \alpha^{\mathbb{Z}}\left(=\operatorname{Pic}^{0}(X)\right) \rightarrow 1$ implies that the first factor on the right is $\left|\operatorname{Pic}^{0}(X)\right| /(q-1)$. The exact sequence

$$
0 \rightarrow \mathbb{F}_{q} \rightarrow O_{\mathbb{A}} \rightarrow \mathbb{A} / F \rightarrow H^{1}\left(X, O_{X}\right) \rightarrow 0
$$

implies that the second factor on the right is $q^{g-1}$.

## 4 Intertwining operators and Eisenstein series

### 4.1 Intertwining operators

Let $E$ be an algebraically closed field of characteristic zero, and $v \in|X|$ a closed point of $X$. Denote by $|a|_{v}$ the absolute value of $a \in F_{v}^{\times}$normalized by $\left|\boldsymbol{\pi}_{v}\right|=q_{v}^{-1}$. It is an $E^{\times}$-valued character of $F_{v}^{\times}$. Fix a square root $\sqrt{q}=q^{1 / 2}$ of $q$ in $E$. If $E \subset \mathbb{C}$ we choose $q^{1 / 2}>0$. For $E$-valued characters
$\mu_{1}, \mu_{2}$ of $F_{v}^{\times}$denote by $I\left(\mu_{1}, \mu_{2}\right)$ both the space of right locally constant functions $\phi: \operatorname{GL}\left(2, F_{v}\right) \rightarrow E$ with $\phi\left(\left(\begin{array}{cc}a_{1} & b \\ 0 & a_{2}\end{array}\right) x\right)=\left|a_{1} / a_{2}\right|_{v}^{1 / 2} \mu_{1}\left(a_{1}\right) \mu_{2}\left(a_{2}\right) \phi(x)(x \in$ $\left.\mathrm{GL}\left(2, F_{v}\right) ; a_{1}, a_{2} \in F_{v}^{\times} ; b \in F_{v}\right)$, and the action of the group GL $\left(2, F_{v}\right)$ by right translation on $I\left(\mu_{1}, \mu_{2}\right)$. The induced representation $I\left(\mu_{1}, \mu_{2}\right)$ is admissible by the Iwasawa decomposition $G=B K$. It is unitarizable when $\mu_{1}, \mu_{2}$ are unitary. It is possible to work with $I\left(|\cdot|{ }_{v}^{1 / 2} \mu_{1},|\cdot|{ }_{v}^{1 / 2} \mu_{2}\right)$, in whose definition the factor $\left|a_{1} / a_{2}\right|_{v}^{1 / 2} \mu_{1}\left(a_{1}\right) \mu_{2}\left(a_{2}\right)$ becomes $\left|a_{1}\right|_{v} \mu_{1}\left(a_{1}\right) \mu_{2}\left(a_{2}\right)$, but later we shall need to multiply back by $|\cdot|_{v}^{-1 / 2}$. The following is a standard basic result.
Proposition 4.1. If $\mu_{1} / \mu_{2} \neq|\cdot|_{v},|\cdot|_{v}^{-1}$, then the representations of $\operatorname{GL}\left(2, F_{v}\right)$ in $I\left(\mu_{1}, \mu_{2}\right)$ and $I\left(\mu_{2}, \mu_{1}\right)$ are irreducible and isomorphic. If $\mu_{1} / \mu_{2}=|\cdot|_{v}$ or $|\cdot|_{v}^{-1}$ then $I\left(\mu_{1}, \mu_{2}\right)$ contains a unique proper invariant subspace $I^{\prime}\left(\mu_{1}, \mu_{2}\right)$ and there is a $\mathrm{GL}\left(2, F_{v}\right)$-isomorphism $I^{\prime}\left(\mu_{1}, \mu_{2}\right) \simeq I\left(\mu_{2}, \mu_{1}\right) / I^{\prime}\left(\mu_{2}, \mu_{1}\right)$. If $\mu_{2} / \mu_{1}=$ $|\cdot|_{v}$, the subspace $I^{\prime}\left(\mu_{1}|\cdot|_{v}^{-1 / 2}, \mu_{1}|\cdot|{ }_{v}^{1 / 2}\right)$ is one dimensional; $x \in \operatorname{GL}\left(2, F_{v}\right)$ acts on $I^{\prime}\left(\mu_{1}|\cdot|{ }_{v}^{-1 / 2}, \mu_{1}|\cdot|_{v}^{1 / 2}\right)$ via multiplication by $\mu_{1}(x)$. The subspace

$$
I^{\prime}\left(\mu_{2}|\cdot|{ }_{v}^{1 / 2},\left.\mu_{2}|\cdot|\right|_{v} ^{-1 / 2}\right) \quad \text { is denoted by } \quad \operatorname{St}\left(\mu_{2}\right)=\operatorname{St}\left(\mu_{2}|\cdot|{ }_{v}^{1 / 2},\left.\mu_{2}|\cdot|\right|_{v} ^{-1 / 2}\right)
$$

It is isomorphic to $I\left(\mu_{2}|\cdot|{ }_{v}^{-1 / 2}, \mu_{2}|\cdot|{ }_{v}^{1 / 2}\right) / I^{\prime}\left(\mu_{2}|\cdot|{ }_{v}^{-1 / 2}, \mu_{2}|\cdot|{ }_{v}^{1 / 2}\right)$. It consists of

$$
\phi \in I\left(\left.\mu_{2}|\cdot|\right|_{v} ^{1 / 2}, \mu_{2}|\cdot|{ }_{v}^{-1 / 2}\right) \quad \text { with } \quad \int_{\mathrm{GL}\left(2, O_{v}\right)} \mu_{2}(\operatorname{det} x)^{-1} \phi(x) d x=0
$$

If $I\left(\mu_{1}, \mu_{2}\right) \simeq I\left(\mu_{1}^{\prime}, \mu_{2}^{\prime}\right)$ then $\left\{\mu_{1}, \mu_{2}\right\}=\left\{\mu_{1}^{\prime}, \mu_{2}^{\prime}\right\}$, the representations $I\left(\mu_{1}, \mu_{2}\right)$ $\left(\mu_{1} / \mu_{2} \neq|\cdot|_{v}\right.$ or $\left.|\cdot|_{v}^{-1}\right)$ and $\operatorname{St}\left(\mu_{2}^{\prime}\right)$ are infinite dimensional and inequivalent, and $\operatorname{St}\left(\mu_{1}\right) \simeq \operatorname{St}\left(\mu_{2}\right)$ implies $\mu_{1}=\mu_{2}$.

We proceed to describe the operator intertwining $I\left(\mu_{1}, \mu_{2}\right)$ and $I\left(\mu_{2}, \mu_{1}\right)$.
Proposition 4.2. If $\left|\mu_{1}\left(\boldsymbol{\pi}_{v}\right) / \mu_{2}\left(\boldsymbol{\pi}_{v}\right)\right|<1$ the integral

$$
(M \phi)(x)=\int_{F_{v}} \phi\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)\left(\begin{array}{ll}
1 & y \\
0 & 1
\end{array}\right) x\right) d y
$$

converges for each $\phi \in I\left(\mu_{1}, \mu_{2}\right)$ and $x \in \mathrm{GL}\left(2, F_{v}\right)$, and $M \phi \in I\left(\mu_{2}, \mu_{1}\right)$.
Proof. As $\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right)\left(\begin{array}{ll}1 & y \\ 0 & 1\end{array}\right)=\left(\begin{array}{cc}y^{-1} & -1 \\ 0 & y\end{array}\right)\left(\begin{array}{cc}1 & 0 \\ y^{-1} & 1\end{array}\right)$, the integrand is

$$
\mu_{2}(y) \mu_{1}(y)^{-1}|y|_{v}^{-1} \phi\left(\left(\begin{array}{cc}
\frac{1}{1} & 0 \\
y^{-1} & 1
\end{array}\right) x\right),
$$

which is 0 if $|y|_{v}$ is small, and $\mu_{2}(y) \mu_{1}(y)^{-1}|y|_{v}^{-1} \phi(x)$ if $|y|_{v}$ is big enough. For sufficiently large $n$ then the part of the integral over $|y|_{v} \geq q_{v}^{n}$ is bounded by $\phi(x)$ times

$$
\int_{|y|_{v} \geq q_{v}^{n}}\left|\mu_{2}(y) / \mu_{1}(y)\right| \cdot|y|_{v}^{-1} d y=\left|O_{v}^{\times}\right| \sum_{k \geq n}\left|\mu_{1}\left(\boldsymbol{\pi}_{v}\right) / \mu_{2}\left(\boldsymbol{\pi}_{v}\right)\right|^{k}<\infty .
$$

It is clear that $(M \phi)\left(\left(\begin{array}{ll}1 & c \\ 0 & 1\end{array}\right) x\right)=(M \phi)(x)\left(c \in F_{v}\right)$ and $(M \phi)\left(\left(\begin{array}{ll}a & 0 \\ 0 & b\end{array}\right) x\right)$ equals

$$
\int_{F_{v}} \phi\left(\left(\begin{array}{cc}
b & 0 \\
0 & a
\end{array}\right)\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)\left(\begin{array}{cc}
1 & y b / a \\
0 & 1
\end{array}\right) x\right) d y=\mu_{1}(b) \mu_{2}(a)\left|\frac{b}{a}\right|_{v}^{1 / 2}\left|\frac{a}{b}\right|_{v}(M \phi)(x) .
$$

We obtained, if $\left|\mu_{1}\left(\boldsymbol{\pi}_{v}\right) / \mu_{2}\left(\boldsymbol{\pi}_{v}\right)\right|<1$, a GL $\left(2, F_{v}\right)$-equivariant map

$$
M=M\left(\mu_{1}, \mu_{2}\right): I\left(\mu_{1}, \mu_{2}\right) \rightarrow I\left(\mu_{2}, \mu_{1}\right) .
$$

Let $\nu_{t}$ be the unramified character of $F_{v}^{\times}$with $\nu_{t}\left(\boldsymbol{\pi}_{v}\right)=t$. Put $M\left(\mu_{1}, \mu_{2}, t\right)=$ $M\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$. It converges for any $\mu_{1}, \mu_{2}$, provided $t \in \mathbb{C}$ is small enough in absolute value. To define $M\left(\mu_{1}, \mu_{2}\right)$ as the value at $t=1$ of the analytic continuation of $M\left(\mu_{1}, \mu_{2}, t\right)$, we need these operators to be defined on the same space, which we will take to be

$$
\begin{gathered}
I_{0}\left(\mu_{1}, \mu_{2}\right)=\left\{\phi \in C^{\infty}\left(\mathrm{GL}\left(2, O_{v}\right)\right) ; \phi\left(\left(\begin{array}{cc}
a_{1} & b \\
0 & a_{2}
\end{array}\right) x\right)=\mu_{1}\left(a_{1}\right) \mu_{2}\left(a_{2}\right) \phi(x),\right. \\
\left.a_{1}, a_{2} \in O_{v}^{\times}, b \in O_{v}, x \in \operatorname{GL}\left(2, O_{v}\right)\right\} .
\end{gathered}
$$

By the Iwasawa decomposition $G=B K$, the restriction map $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t-1}\right)$ $\rightarrow I_{0}\left(\mu_{1}, \mu_{2}\right)$ is bijective for any $t$. Identifying these spaces, the operator $M\left(\mu_{1}, \mu_{2}, t\right)$ becomes a map $I_{0}\left(\mu_{1}, \mu_{2}\right) \rightarrow I_{0}\left(\mu_{2}, \mu_{1}\right)$.
Write $L(\mu, t)$ for $\left(1-\mu\left(\boldsymbol{\pi}_{v}\right) t\right)^{-1}$ if $\mu$ is unramified, and $L(\mu, t)=1$ if $\mu$ is a ramified character of $F_{v}{ }^{\times}$.

Proposition 4.3. The operator valued function $M\left(\mu_{1}, \mu_{2}, t\right)$ is rational in $t \in$ $\mathbb{C}^{\times}$. In fact the function $t \mapsto L\left(\mu_{1} / \mu_{2}, t^{2}\right)^{-1}\left(M\left(\mu_{1}, \mu_{2}, t\right) \phi\right)(x)$ is a polynomial in $t$ for all $\phi \in I_{0}\left(\mu_{1}, \mu_{2}\right), x \in \mathrm{GL}\left(2, O_{v}\right)$. If $\mu_{1}, \mu_{2}$ are unramified and the restrictions of $\phi \in I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ and $\psi \in I\left(\mu_{2} \nu_{t^{-1}}, \mu_{1} \nu_{t}\right)$ to $\mathrm{GL}\left(2, O_{v}\right)$ are 1, then $M\left(\mu_{1}, \mu_{2}, t\right) \phi=\frac{L\left(\mu_{1} / \mu_{2}, t^{2}\right)}{L\left(\mu_{1} / \mu_{2}, q_{v}^{-1} t^{2}\right)} \psi$.

Proof. Put $\phi_{t}=M\left(\mu_{1}, \mu_{2}, t\right) \phi$ and $a_{1}=\int_{|y|_{v} \leq 1} \phi\left(\left(\begin{array}{cc}0 & -1 \\ 1 & y\end{array}\right) x\right) d y$ where $x \in$ $\mathrm{GL}\left(2, O_{v}\right)$. Then

$$
\phi_{t}(x)=a_{1}+\int_{|y|_{v}>1} \mu_{2}(y) \mu_{1}(y)^{-1}|y|_{v}^{-1} \nu_{t}(y)^{-2} \phi\left(\left(\begin{array}{cc}
1 & 0 \\
y^{-1} & 1
\end{array}\right) x\right) d y .
$$

We shall show that this is the Taylor series of a rational function.
If $n$ is large enough, $\phi\left(\left(\begin{array}{rr}1 & 0 \\ y^{-1} & 1\end{array}\right) x\right)=\phi(x)$ for $|y|_{v} \geq q_{v}^{n}$. Then $\phi_{t}(x)=$ $a_{1}+a_{2}(t)+a_{3}(t)$ with

$$
\begin{aligned}
& a_{2}(t)=\int_{1<|y|_{v}<q_{v}^{n}} \mu_{2}(y) \mu_{1}(y)^{-1}|y|_{v}^{-1} \nu_{t}(y)^{-2} \phi\left(\left(\begin{array}{cc}
1 & 0 \\
y^{-1} & 1
\end{array}\right) x\right) d y \\
& a_{3}(t)=\phi(x) \int_{|y|_{v} \geq q_{v}^{n}} \mu_{2}(y) \mu_{1}(y)^{-1}|y|_{v} \nu_{t}(y)^{-2} d y
\end{aligned}
$$

Clearly $a_{2}(t)$ is a polynomial in $t$ (since $\nu_{t}\left(\boldsymbol{\pi}_{v}^{-1}\right)^{-1}=t$ ) and $a_{3}(t)=$ $c t^{2 n} L\left(\mu_{1} / \mu_{2}, t^{2}\right)$.
If $\mu_{1}, \mu_{2}$ are unramified and $x \in \mathrm{GL}\left(2, O_{v}\right), a_{1}=1$ and the expression for $\phi_{t}(x)$ is

$$
\begin{aligned}
\phi_{t}(x)=1 & +\int_{|y|_{v}>1} \mu_{2}(y) \mu_{1}(y)^{-1}|y|_{v}^{-1} \nu_{t}(y)^{-2} d y \\
& =1-\left(1-q_{v}^{-1}\right) \sum_{k \geq 1}\left(\mu_{1}\left(\boldsymbol{\pi}_{v}\right) / \mu_{2}\left(\boldsymbol{\pi}_{v}\right)\right)^{k} t^{2 k} \\
& =1+\frac{\left(1-q_{v}^{-1}\right)\left(\mu_{1}\left(\boldsymbol{\pi}_{v}\right) / \mu_{2}\left(\boldsymbol{\pi}_{v}\right)\right) t^{2}}{1-\left(\mu_{1}\left(\boldsymbol{\pi}_{v}\right) / \mu_{2}\left(\boldsymbol{\pi}_{v}\right)\right) t^{2}}=\frac{L\left(\mu_{1} / \mu_{2}, t^{2}\right)}{L\left(\mu_{1} / \mu_{2}, q_{v}^{-1} t^{2}\right)}
\end{aligned}
$$

The operator $M\left(\mu_{1}, \mu_{2}, t\right): I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right) \rightarrow I\left(\mu_{2} \nu_{t^{-1}}, \mu_{1} \nu_{t}\right)$ intertwines the GL $\left(2, F_{v}\right)$-modules for every $t$ where it is defined. It can be regarded as a rational function of $t$ (in fact, of $t^{2}$ ) with values in the set of operators $I_{0}\left(\mu_{1}, \mu_{2}\right) \rightarrow I_{0}\left(\mu_{2}, \mu_{1}\right)$. Indeed,

$$
M\left(\mu_{1}, \mu_{2}, t\right)=M\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)=M\left(\mu_{1} \nu_{t^{2}}, \mu_{2}\right)
$$

Define

$$
R\left(\mu_{1}, \mu_{2}, t\right)=\frac{L\left(\mu_{1} / \mu_{2}, q_{v}^{-1} t^{2}\right)}{L\left(\mu_{1} / \mu_{2}, t^{2}\right)} M\left(\mu_{1}, \mu_{2}, t\right)
$$

Corollary 4.4. Suppose $\mu_{1}$ and $\mu_{2}$ are unramified and $\varphi \in I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$, $\psi \in I\left(\mu_{2} \nu_{t^{-1}}, \mu_{1} \nu_{t}\right)$ are the functions whose restrictions to $\mathrm{GL}\left(2, O_{v}\right)$ are one, then $R\left(\mu_{1}, \mu_{2}, t\right) \varphi=\psi$.

Given characters $\mu_{1}, \mu_{2}$ of $\mathbb{A}^{\times}$, write $I\left(\mu_{1}, \mu_{2}\right)$ for the space of right locally constant functions $\phi$ on $\operatorname{GL}(2, \mathbb{A})$ which satisfy

$$
\phi\left(\left(\begin{array}{cc}
a_{1} & b \\
0 & a_{2}
\end{array}\right) x\right)=\mu_{1}\left(a_{1}\right) \mu_{2}\left(a_{2}\right)\left|a_{1} / a_{2}\right|^{1 / 2} \phi(x) . \quad \text { Put } \quad \nu(a)=q^{\operatorname{deg}(a)}
$$

Then $I\left(\mu_{1}, \mu_{2}\right)$ is the restricted tensor product of the spaces $I\left(\mu_{1 v}, \mu_{2 v}\right)$ where $\mu_{i v}$ is the component of $\mu_{i}$ at $v$ (the restriction of $\mu_{i}$ to $F_{v}^{\times} \hookrightarrow \mathbb{A}^{\times}$); it is spanned by $\otimes_{v} \phi_{v}$ with $\phi_{v} \in I\left(\mu_{1 v}, \mu_{2 v}\right)$ for all $v$ and $\phi_{v} \mid \mathrm{GL}\left(2, O_{v}\right)=1$ for almost all $v$, where $\mu_{i v} \mid O_{v}^{\times}=1$, i.e. $\mu_{i v}$ are unramified. Define the character $\nu_{t}$ of $\mathbb{A}^{\times}$by $\nu_{t}(a)=t^{\operatorname{deg}(a)}$. Then the restriction of $\nu_{t}$ to $F_{v}^{\times}$is $\nu_{t_{v}}$, the unramified character of $F_{v}^{\times}$with $\nu_{t_{v}}\left(\boldsymbol{\pi}_{v}\right)=t_{v}\left(=t^{\operatorname{deg}(v)}\right)$. As in the local case, we identify the spaces $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ with $I_{0}\left(\mu_{1}, \mu_{2}\right)$ for all $t$. The operator $R\left(\mu_{1}, \mu_{2}, t\right)$ from $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ to $I\left(\mu_{2} \nu_{t^{-1}}, \mu_{1} \nu_{t}\right)$ defined by $R\left(\mu_{1}, \mu_{2}, t\right)=$ $\otimes_{v} R\left(\mu_{1 v}, \mu_{2 v}, t_{v}\right)$ is rational in $t$. On any element in $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ at most finitely many components $R\left(\mu_{1 v}, \mu_{2 v}, t_{v}\right)$ do not act as the identity. Also write $m(\mu, t)$ for $L(\mu, t) / L(\mu, t / q)$.

### 4.2 Eisenstein series

Write $A_{\alpha}=C^{\infty}\left(\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right)$,

$$
A_{c, \alpha}=C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right), \quad Y=A(F) N(\mathbb{A}) \backslash \mathrm{GL}(2, \mathbb{A})
$$

and $Y_{\alpha}=Y / \alpha^{\mathbb{Z}}$. Normalize the Haar measure on $N(\mathbb{A}) \simeq \mathbb{A}$ by $|N(\mathbb{A}) / N(F)|=$ $|\mathbb{A} / F|=1$. The Haar measure on $N(\mathbb{A})$ is invariant with respect to conjugation by the elements of $A(F)$ by the product formula. So it extends to a two-sided invariant measure on the space $\alpha^{\mathbb{Z}} \cdot A(F) N(\mathbb{A})$. This, and the two-sided Haar measure on $\operatorname{GL}(2, \mathbb{A})$ induce an invariant measure on $Y_{\alpha}$.
Let $\varphi$ and $\psi$ be locally constant functions on $Y_{\alpha}$, at least one of which is compactly supported. Put $(\varphi, \psi)=\int_{Y_{\alpha}} \varphi(x) \bar{\psi}(x) d x$. On $\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ a scalar product is similarly defined. Define the map $E^{*}: A_{\alpha} \rightarrow C^{\infty}\left(Y_{\alpha}\right)$ by

$$
\phi \mapsto \phi_{N}, \quad \phi_{N}(x)=\int_{N(F) \backslash N(\mathbb{A})} \phi(n x) d n, \quad x \in \mathrm{GL}(2, \mathbb{A}) .
$$

Note that $N(F) \backslash N(\mathbb{A})$ is compact, so the integral converges. Note that ker $E^{*}$ is the space $A_{0, \alpha}$ of cusp forms invariant under $\alpha$. For any $f \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ define a function $E f$ on $\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ by

$$
(E f)(x)=\sum_{\gamma \in A(F) N(F) \backslash \mathrm{GL}(2, F)} f(\gamma x), \quad x \in \mathrm{GL}(2, \mathbb{A}) .
$$

Proposition 4.5. The sum defining $(E f)(x)$ converges. For $f \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ and $\phi \in A_{\alpha}$ we have $(E f, \phi)=\left(f, E^{*} \phi\right)$.
Proof. Consider the diagram

$$
Y_{\alpha} \stackrel{r}{\leftarrow} \alpha^{\mathbb{Z}} \cdot A(F) N(F) \backslash \mathrm{GL}(2, \mathbb{A}) \xrightarrow{s} \alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) .
$$

Since $N(F) \backslash N(\mathbb{A})$ is compact, the map $r$ is proper. Hence the natural embedding $r^{*}$ maps $C_{c}^{\infty}\left(Y_{\alpha}\right)$ to $C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} \cdot A(F) N(F) \backslash \mathrm{GL}(2, \mathbb{A})\right)$. Given

$$
\psi \in C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} A(F) N(F) \backslash \mathrm{GL}(2, \mathbb{A})\right),
$$

define a function $s_{*} \psi$ on $\alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$ by

$$
\left(s_{*} \psi\right)(x)=\sum_{\gamma \in A(F) N(F) \backslash \operatorname{GL}(2, F)} \psi(\gamma x), \quad x \in \mathrm{GL}(2, \mathbb{A}) .
$$

The sum is finite since $\psi$ is compactly supported, and

$$
s_{*} \psi \in C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right) .
$$

The sum which defines $(E f)(x)$ converges since $E=s_{*} r^{*}$.
Now define $E^{*}=r_{*} s^{*}$, where $s^{*}$ is the natural embedding, and

$$
r_{*}: C^{\infty}\left(\alpha^{\mathbb{Z}} A(F) N(F) \backslash \mathrm{GL}(2, \mathbb{A})\right) \rightarrow C^{\infty}\left(Y_{\alpha}\right)
$$

is defined by $\left(r_{*} h\right)(x)=\int_{N(F) \backslash N(\mathbb{A})} h(n x) d n, x \in \mathrm{GL}(2, \mathbb{A})$. Since $\left(r^{*}, r_{*}\right)$ and $\left(s_{*}, s^{*}\right)$ are adjoint pairs, so is $\left(E=s_{*} r^{*}, E^{*}=r_{*} s^{*}\right)$.

The image $A_{E, \alpha}$ of the Eisenstein map $E=s_{*} r^{*}: C_{c}^{\infty}\left(Y_{\alpha}\right) \rightarrow A_{c, \alpha}$ is called the Eisenstein part of $A_{c, \alpha}$. The maps $E$ and $E^{*}$ intertwine the GL $(2, \mathbb{A})$-action; $A_{E, \alpha}$ is an invariant subspace of $A_{c, \alpha}$.

Proposition 4.6. The space $A_{c, \alpha}$ is an orthogonal direct sum of the space $A_{0, \alpha}$ of cusp forms and of $A_{E, \alpha}$.

Proof. Cusp forms are compactly supported. Since $A_{0, \alpha}=\operatorname{ker} E^{*}$ and $A_{E, \alpha}=$ $\operatorname{im} E$, we have $A_{0, \alpha} \perp A_{E, \alpha}$. Given a compact open subgroup $U$ in $\operatorname{GL}(2, \mathbb{A})$, put $A_{\alpha}^{U}$ for the space of $U$-invariant functions in $A_{\alpha}$, and

$$
A_{c, \alpha}^{U}=A_{c, \alpha} \cap A_{\alpha}^{U}, \quad A_{0, \alpha}^{U}=A_{0, \alpha} \cap A_{\alpha}^{U}, \quad A_{E, \alpha}^{U}=A_{E, \alpha} \cap A_{\alpha}^{U}
$$

It remains to show that $A_{0, \alpha}^{U}+A_{E, \alpha}^{U}=A_{c, \alpha}^{U}$. If not there exists a nonzero linear form $\ell: A_{c, \alpha}^{U} \rightarrow \mathbb{C}$ which is zero on $A_{0, \alpha}^{U}+A_{E, \alpha}^{U}$. There exists $f \in A_{\alpha}^{U}$ such that $\ell(\phi)=(\phi, f)$ for every $\phi \in A_{c, \alpha}^{U}$. For any $U$-invariant function $\psi \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ we have $\left(\psi, E^{*} f\right)=(E \psi, f)=\ell(E \psi)=0$. Hence $E^{*} f=0$, thus $f \in A_{0, \alpha}^{U}$. This however is impossible since $f$ is orthogonal to the space $A_{0, \alpha}^{U}$ of $U$-invariant cusp forms.

Given $\phi \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ and $x \in \operatorname{GL}(2, \mathbb{A})$, put $(M \phi)(x)=\int_{N(\mathbb{A})} \phi\left(\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x\right) d n$. The integral converges, by

Proposition 4.7. The map $N(\mathbb{A}) \rightarrow Y_{\alpha}, n \mapsto \alpha^{\mathbb{Z}} A(F) N(\mathbb{A})\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x$, is proper.

Proof. It suffices to consider the case of $x=1$. The function

$$
\mathrm{ht}^{+}: Y_{\alpha} \rightarrow \mathbb{Z}, \quad\left(\begin{array}{cc}
a & c \\
0 & b
\end{array}\right) k \mapsto \operatorname{deg} a-\operatorname{deg} b,
$$

is continuous. Thus it suffices to show that the map $\varphi(a)=\mathrm{ht}^{+}\left(\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right)\left(\begin{array}{ll}1 & a \\ 0 & 1\end{array}\right)\right)$, $\varphi: \mathbb{A} \rightarrow \mathbb{Z}$, is proper. But $\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right)\left(\begin{array}{cc}1 & a_{v} \\ 0 & 1\end{array}\right)$ is in $\operatorname{GL}\left(2, O_{v}\right)$ if $\left|a_{v}\right|_{v} \leq 1$; otherwise it is $=\left(\begin{array}{cc}a_{v}^{-1} & -1 \\ 0 & a_{v}\end{array}\right)\left(\begin{array}{cc}1 & 0 \\ a_{v}^{-1} & 1\end{array}\right)$. If $a=\left(a_{v}\right)$, then $\varphi(a)=-2 \sum_{v} \max \left(0, \log _{q}\left|a_{v}\right|_{v}\right)$, as $\log _{q}\left|a_{v}\right|_{v}=-\operatorname{val}_{v}\left(a_{v}\right) \operatorname{deg}(v)$. Hence $\varphi$ is proper.

By definition, $x \mapsto(M \phi)(x)$ is invariant under left translation by $N(\mathbb{A})$, and also by $\alpha^{\mathbb{Z}} \cdot A(F)$. Indeed,

$$
(M \phi)\left(\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right) x\right)=\int_{\mathbb{A}} \phi\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right) x\right) d y=\left|\frac{a}{b}\right| \int_{N(\mathbb{Z})} \phi\left(\left(\begin{array}{cc}
b & 0 \\
0 & a
\end{array}\right)\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) d n
$$

and $|a / b|=q^{\operatorname{deg}(a / b)}$. Thus $M$ maps $C_{c}^{\infty}\left(Y_{\alpha}\right)$ to $C^{\infty}\left(Y_{\alpha}\right)$.
Proposition 4.8. Denote by $I$ the natural embedding of $C_{c}^{\infty}\left(Y_{\alpha}\right)$ in $C^{\infty}\left(Y_{\alpha}\right)$. Then

$$
E^{*} E=I+M
$$

Proof. By the Bruhat decomposition, an element of GL $(2, F)$ which is not in $A(F) N(F)$ has a unique decomposition $n_{1} a\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n_{2}$ with $n_{i} \in N(F), a \in$ $A(F)$. Thus, for any $\phi \in C_{c}^{\infty}\left(Y_{\alpha}\right), x \in \mathrm{GL}(2, \mathbb{A})$, we have

$$
(E \phi)(x)=\sum_{\gamma \in A(F) N(F) \backslash G L(2, F)} \phi(\gamma x)=\phi(x)+\sum_{\nu \in N(F)} \phi\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) \nu x\right)
$$

Hence

$$
\begin{aligned}
\left(E^{*} E \phi\right)(x) & =|N(\mathbb{A}) / N(F)| \phi(x)+\int_{N(F) \backslash N(\mathbb{A})} \sum_{\nu \in N(F)} \phi\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) \nu n x\right) d n \\
& =\phi(x)+\int_{N(\mathbb{A})} \phi\left(\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) d n=\phi(x)+(M \phi)(x)\right.
\end{aligned}
$$

Proposition 4.9. Let $\mu_{1}, \mu_{2}$ be characters of $\mathbb{A}^{\times} / F^{\times}$. Ift is sufficiently small, for all $\phi \in I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ and $x \in \mathrm{GL}(2, \mathbb{A})$, the integral $\left(M\left(\mu_{1}, \mu_{2}, t\right) \phi\right)(x)=$ $\int_{N(\mathbb{A})} \phi\left(\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x\right) d n$ converges and defines a function in $I\left(\mu_{2} \nu_{t^{-1}}, \mu_{1} \nu_{t}\right)$. Moreover, $M\left(\mu_{1}, \mu_{2}, t\right)=q^{1-g} m\left(\mu_{1} / \mu_{2}, t^{2}\right) R\left(\mu_{1}, \mu_{2}, t\right)$.
Proof. Recall that $|a|=q^{\operatorname{deg}(a)}$ and that $I\left(\mu_{1}, \mu_{2}\right)$ consists of the $\phi$ in $C^{\infty}(\mathrm{GL}(2, \mathbb{A}))$ with

$$
\phi\left(\left(\begin{array}{cc}
a_{1} & 0 \\
0 & a_{2}
\end{array}\right) x\right)=\left|a_{1} / a_{2}\right|^{1 / 2} \mu_{1}\left(a_{1}\right) \mu_{2}\left(a_{2}\right) \phi(x)
$$

while $\nu_{t}(a)=t^{\operatorname{deg} a}$. We put $t_{v}=t^{\operatorname{deg}(v)}$. We may assume that $\phi(x)=$ $\prod_{v} \phi_{v}\left(x_{v}\right)$ with $\phi_{v} \in I\left(\mu_{1 v} \nu_{t_{v}}, \mu_{2 v} \nu_{t_{v}^{-1}}\right)$. For almost all $v$, the restriction of $\phi_{v}$ to $\operatorname{GL}\left(2, O_{v}\right)$ is 1 . We may replace $\phi_{v}, \mu_{i}, t$ by their complex absolute values to assume $t>0$ and $\phi_{v}, \mu_{i}$ take real nonnegative values. Then $\left(M\left(\mu_{1}, \mu_{2}, t\right) \phi\right)(x)=c \prod_{v} \tau_{v}$, with $\tau_{v}=\int_{N\left(F_{v}\right)} \phi_{v}\left(\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x_{v}\right) d n=$ $\int_{F_{v}} \phi_{v}\left(\left(\begin{array}{cc}0 & -1 \\ 1 & z\end{array}\right) x_{v}\right) d z$. The measure $d n_{v}$ on $N\left(F_{v}\right)$ is normalized by $\left|N\left(O_{v}\right)\right|=1$, and $c=|N(\mathbb{A}) / N(F)|$ in the measure $\otimes_{v} d n_{v}$ on $N(\mathbb{A})$.
We saw that for small enough $t$ the integral which defines $\tau_{v}$ converges for all $v$. For almost all $v$ we have $\tau_{v}=L\left(\mu_{1 v} / \mu_{2 v}, t_{v}^{2}\right) / L\left(\mu_{1 v} / \mu_{2 v}, q_{v}^{-1} t_{v}^{2}\right)$, so the product $\prod_{v} \tau_{v}$ converges for small $t$. Now $M\left(\mu_{1}, \mu_{2}, t\right)=c \prod_{v} M\left(\mu_{1 v}, \mu_{2 v}, t_{v}\right)$. Each factor here is $\frac{L\left(\mu_{1 v} / \mu_{2 v}, t_{v}^{2}\right)}{L\left(\mu_{1 v} / \mu_{2 v}, q_{v}^{-1} t_{v}^{2}\right)} R\left(\mu_{1 v}, \mu_{2 v}, t_{v}\right)$. Put $R\left(\mu_{1}, \mu_{2}, t\right)=\otimes_{v} R\left(\mu_{1 v}, \mu_{2 v}, t_{v}\right)$, and $m(\mu, t)=\frac{L(\mu, t)}{L\left(q^{-1} t, \mu\right)}$, where $L(\mu, t)=\prod_{v} L\left(\mu_{v}, t_{v}\right)$. Note that $c$ is $|O|=q^{1-g}$, using $0 \rightarrow \mathbb{F}_{q} \rightarrow O \rightarrow \mathbb{A} / F \rightarrow H^{1}\left(X, O_{X}\right) \rightarrow 0$.
It follows (since $L(\mu, t)$ is a rational function of $t$ ) that after identifying the spaces $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ for all $t$, the operator

$$
M\left(\mu_{1}, \mu_{2}, t\right): I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right) \rightarrow I\left(\mu_{2} \nu_{t^{-1}}, \mu_{1} \nu_{t}\right)
$$

(defined for small $t$ ) depends on $t$ rationally. Hence $M\left(\mu_{1}, \mu_{2}, t\right)$ is defined for almost all $t$, and it commutes with the action of $\operatorname{GL}(2, \mathbb{A})$.

## 4.3 $L$-FUNCTIONS

Let us review the theory of $L$-functions for GL(2). Let $E$ be an algebraically closed field of characteristic zero. The valuation $\operatorname{val}_{v}(a)$ of $a \in F_{v}^{\times}$is the largest integer $n$ with $a \in \pi_{v}^{n} O_{v}$. For any character $\psi: F_{v} \rightarrow E^{\times}, \psi \neq 1$, let $r(\psi)$ be the largest $n$ such that $\psi\left(\boldsymbol{\pi}_{v}^{-n} O_{v}\right)=1$. Normalize the Haar measure on $F_{v}$ by $\left|O_{v}\right|=1$. The conductor of a character $\chi: F_{v}^{\times} \rightarrow E^{\times}$is $n=0$ if $\chi\left(O_{v}^{\times}\right)=1$, i.e., $\chi$ is unramified; otherwise it is the smallest $n \geq 1$ such that $\chi\left(1+\pi_{v}^{n} O_{v}\right)=1$. Given $\chi$, put $L(t, \chi)=\left(1-\chi\left(\boldsymbol{\pi}_{v}\right) t\right)^{-1}$ if $\chi$ is unramified, $L(t, \chi)=1$ is $\chi$ is ramified. Given $\psi \neq 1$, put

$$
\Gamma(\chi, \psi, t)=\int_{F_{v}^{\times}} \chi(x)^{-1} \psi(x) t^{-\operatorname{val}_{v}(x)} d x, \quad \psi: F_{v} \rightarrow E^{\times}
$$

This $\Gamma(\chi, \psi, t)$ is a formal power series in $t$ which contains positive and negative powers of $t$. Tate's thesis (see [Lg94], VII, section 3-4) establishes

Proposition 4.10. The formal series $\Gamma(\chi, \psi, t)$ has finitely many positive powers of $t$. It is a rational function of $t$, namely a Laurent series of a rational function of $t$ at $t=\infty$. Put $\varepsilon(\chi, \psi, t)=\frac{L(\chi, t) \Gamma(\chi, \psi, t)}{L\left(\chi^{-1}, q_{v}^{-1} t^{-1}\right)}$. It has the form $c(\chi, \psi) t^{n(\chi, \psi)}$. If $r(\psi)=0$ then $n(\chi, \psi)$ is the conductor of $\chi$. If in addition $\chi$ is unramified then $\varepsilon(\chi, \psi, t)$ is 1 . If $a \in F_{v}^{\times}, \psi_{a}(x)=\psi(a x)$, then $\varepsilon\left(\chi, \psi_{a}, t\right)=\chi(a)\left(q_{v} t\right)^{\operatorname{val}_{v}(a)} \varepsilon(\chi, \psi, t)$.

Note that $L$ and $\varepsilon$ are usually considered, in the case where $E=\mathbb{C}$, as functions of $s$, where $t=q_{v}^{-s}$, rather than of $t$. The Haar measure on $F_{v}$ is usually normalized by $\left|O_{v}\right|=q_{v}^{-r(\psi) / 2}$, as this measure is self-dual with respect to the pairing $F_{v} \times F_{v} \rightarrow E^{\times},(x, y) \mapsto \psi(x y)$. This choice of measure is not convenient if $E \neq \mathbb{C}$ since $E$ has no distinguished square root of $q$.
Given a character $\chi$ of $\mathbb{A}^{\times}$, denote its restriction to $F_{v}^{\times}$by $\chi_{v}$. The restriction to $F_{v}$ of a character $\psi$ of $\mathbb{A}$ is denoted $\psi_{v}$. For a closed point $v$ of $X$, we write $\operatorname{deg}(v)$ for the dimension of the residue field at $v$ over $\mathbb{F}_{q}$, and $q_{v}=q^{\operatorname{deg}(v)}$. Given a character $\chi: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$, put $L(\chi, t)=\prod_{v} L\left(\chi_{v}, t_{v}\right)$, where $t_{v}=$ $t^{\operatorname{deg}(v)}$; the product converges in $E[[t]]$. Let $\psi: \mathbb{A} / F \rightarrow E^{\times}$be a character $\neq 1$. Then $\varepsilon(\chi, t)=q^{1-g} \prod_{v} \varepsilon\left(\chi_{v}, \psi_{v}, t_{v}\right)$ converges as almost all factors are 1 , and $\varepsilon(\chi, t)$ is independent of $\psi$ by Proposition 4.10.

Proposition 4.11. For any character $\chi: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$the formal series $L(\chi, t)$ is rational in $t$, and $L(\chi, t)=\varepsilon(\chi, t) L\left(\chi^{-1}, q^{-1} t^{-1}\right)$. If the restriction of $\chi$ to the group of $x \in \mathbb{A}^{\times} / F^{\times}$with $\operatorname{deg}(x)=0$ is nontrivial, then $L(\chi, t)$ is a polynomial. If the restriction is trivial, $\chi$ is given by $\chi(x)=u^{\operatorname{deg}(x)}$, and then $L(\chi, t)$ has precisely two poles: $t=u^{-1}$ and $t=q^{-1} u^{-1}$, both poles are simple. If $\chi: \mathbb{A}^{\times} / F^{\times} \rightarrow \mathbb{C}^{\times}$is a unitary character $(|\chi(x)|=1$ for all $x)$ then the zeroes of $L(\chi, t)$ lie in the doughnut $\left\{t \in \mathbb{C} ; q^{-1}<|t|<1\right\}$.

The proof of this is also in [Lg94], Chapter VII, sections 7-8. The following is due to [W45].

Theorem 4.12. (A. Weil). For any unitary character $\chi: \mathbb{A}^{\times} / F^{\times} \rightarrow \mathbb{C}^{\times}$, all zeroes of $L(\chi, t)$ lie on the circle $|t|=q^{-1 / 2}$.

Given a character $\psi: \mathbb{A} / F \rightarrow E^{\times}, \psi \neq 1$, let $W(\psi)$ be the space of locally constant functions $\phi: \mathrm{GL}\left(2, F_{v}\right) \rightarrow E$ with $\phi\left(\left(\begin{array}{cc}1 & z \\ 0 & 1\end{array}\right) x\right)=\psi(z) \phi(x)$ for all $z \in F_{v}, x \in \mathrm{GL}\left(2, F_{v}\right)$. The group $\mathrm{GL}\left(2, F_{v}\right)$ acts on $W(\psi)$ by right translation. Fix a Haar measure $d^{\times} x$ on $F_{v}^{\times}$. For any $\phi \in W(\psi)$ put

$$
\Lambda_{\phi}(t)=\int \phi\left(\left(\begin{array}{ll}
a & 0 \\
0 & 1
\end{array}\right)\right)\left(q_{v} t\right)^{\operatorname{val}_{v}(a)} d^{\times} a, \quad \tilde{\Lambda}_{\phi}(t)=\int \phi\left(\left(\begin{array}{ll}
0 & 1 \\
a & 0
\end{array}\right)\right)\left(q_{v} t\right)^{\operatorname{val}_{v}(a)} d^{\times} a .
$$

Both $\Lambda_{\phi}(t)$ and $\tilde{\Lambda}_{\pi}(t)$ are formal power series in $t$, containing positive and negative powers of $t$.
Let $\pi$ be an irreducible admissible representation of $\mathrm{GL}\left(2, F_{v}\right)$ over $E$. Then $\pi\left(\left(\begin{array}{ll}a & 0 \\ 0 & a\end{array}\right)\right)$ is the operator of multiplication by a scalar $\eta(a) \in E^{\times}$. The character $\eta: F_{v}^{\times} \rightarrow E^{\times}$is called the central character of $\pi$.

Proposition 4.13. Let $\pi$ be an irreducible admissible infinite dimensional representation over $E$ of $\mathrm{GL}\left(2, F_{v}\right)$. Let $\eta$ be the central character of $\pi$. (1) There exists a unique $\mathrm{GL}\left(2, F_{v}\right)$-invariant subspace $W(\pi, \psi)$ of $W(\psi)$ equivalent to $\pi$. (2) If $\phi \in W(\pi, \psi)$ then $\Lambda_{\phi}(t)$ is the Laurent series at $t=0$ of a rational function, and $\tilde{\Lambda}_{\phi}(t)$ is the Laurent series at $t=\infty$ of a rational function. (3) There exists a nonzero polynomial $P \in E[t]$ such that for any $\phi \in W(\pi, \psi)$ we have $P(t) \Lambda_{\phi}(t) \in E\left[t, t^{-1}\right]$. There exists $\phi \in W(\pi, \psi)$ with $\Lambda_{\phi}(t) \neq 0$. (4) The quotient $\tilde{\Lambda}_{\phi}(t) / \Lambda_{\phi}(t)$ of rational functions in $t$ does not depend on the choice of $\phi$ in $W(\pi, \psi)$ with $\Lambda_{\phi}(t) \neq 0$. (5) The lowest degree polynomial $P \in E[t]$ which satisfies (3) and $P(0)=1$ is independent of $\psi$. (6) Put $\Gamma(\pi, \psi, t)=\tilde{\Lambda}_{\phi}(t) / \Lambda_{\phi}(t)$ and $\varepsilon(\pi, \psi, t)=\frac{\Gamma(\pi, \psi, t) L(\pi, t)}{L\left(\pi \otimes \eta^{-1}, q_{v}^{-2} t^{-1}\right)}$ where $L(\pi, t)=P(t)^{-1}$ with $P$ of (5). Then $\varepsilon(\pi, \psi, t)$ has the form $c(\pi, \psi) t^{n(\pi, \psi)}, c(\pi, \psi)$ in $E^{\times}$and $n(\pi, \psi)$ in $\mathbb{Z}$. (7) If $\psi_{a}(x)$ is $\psi(a x)$ for $a \in F_{v}^{\times}$, then $\varepsilon\left(\pi, \psi_{a}, t\right)=\eta(a)\left(q_{v} t\right)^{2 \operatorname{val}_{v}(a)} \varepsilon(\pi, \psi, t)$.

This is [JL70], Theorem 2.18. Our $L$ and $\varepsilon$ relate to those $L_{J L}, \varepsilon_{J L}$ of JacquetLanglands by $L_{J L}(\pi, s)=L\left(\pi, t_{v}\right), t_{v}=q_{v}^{-s}, \varepsilon_{J L}(\pi, \psi, s)=\varepsilon\left(\pi, \psi, t_{v}\right)$. Note that the proof of [JL70], which claims that $\Lambda_{\phi}(t)$ is a Laurent series of a meromorphic function in $\mathbb{C}-\{0\}$, shows that $\Lambda_{\phi}(t)$ is rational. In general, the meromorpic functions of $s$ over $p$-adic and global function fields are rational functions of $q^{s}$. Every smooth finite dimensional irreducible representation of $\mathrm{GL}\left(2, F_{v}\right)$ is one dimensional, of the form $x \mapsto \chi(\operatorname{det} x)$, where $\chi: F_{v}^{\times} \rightarrow E^{\times}$ is a character ([JL70], Proposition 2.7).

Proposition 4.14. Let $\pi, \pi^{\prime}$ be irreducible admissible infinite dimensional representations of $\mathrm{GL}\left(2, F_{v}\right)$ with equal central characters. If there is a character $\psi: F_{v} \rightarrow E^{\times}$such that for every character $\omega: F_{v}^{\times} \rightarrow E^{\times}$we have $\Gamma(\pi \omega, \psi, t)=\Gamma\left(\pi^{\prime} \omega, \psi, t\right)$, then $\pi \simeq \pi^{\prime}$.

For a proof see [JL70], Corollary 2.19.

The conductor of an irreducible admissible infinite dimensional representation $\pi$ of $\mathrm{GL}\left(2, F_{v}\right)$ is the integer $n(\pi, \psi)$, with $\psi$ normalized by $r(\psi)=0$. It is well defined, as from (7) above, the integer $n(\pi, \psi)$ of (6) is not changed if $\psi$ is replaced by $\psi_{a}: x \mapsto \psi(a x)$.

Proposition 4.15. The conductor of $\pi$ is the least integer $n$ such that the representation space of $\pi$ contains a nonzero vector invariant under the group $H_{n}=\left\{\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in \mathrm{GL}\left(2, O_{v}\right) ; c \in \boldsymbol{\pi}_{v}^{n} O_{v}, d \in 1+\boldsymbol{\pi}_{v}^{n} O_{v}\right\}$. For this $n, \operatorname{dim}_{E} \boldsymbol{\pi}^{H_{n}}=$ 1.

For a proof see Casselman, Math. Ann. 201 (1973), 301-314.
Proposition 4.16. Let $\pi$ be an irreducible admissible infinite dimensional representation, with central character $\eta$, of $\mathrm{GL}\left(2, F_{v}\right)$. Let $\psi: F_{v} \rightarrow E^{\times}$be a nontrivial character. Then there exists an integer $m_{\pi}$ such that if $\chi: F_{v}^{\times} \rightarrow E^{\times}$is any character with conductor $>m_{\pi}$, then $L(\pi \chi, t)=1$ and

$$
\varepsilon(\pi \chi, \psi, t)=\varepsilon(\chi, \psi, t) \varepsilon\left(\chi \eta, \psi, q_{v} t\right) q_{v}^{-r(\psi)}
$$

For a proof see [JL70], Proposition 3.8. See [JL70], Proposition 3.5, 3.6, for a proof of:

Proposition 4.17. Let $\mu_{1}, \mu_{2}$ be characters of $F_{v}^{\times}$, and $\psi \neq 1$ a character of $F_{v}$. If $\mu_{1} / \mu_{2} \neq|\cdot|_{v}^{ \pm 1}$ then $L\left(I\left(\mu_{1}, \mu_{2}\right), t\right)=L\left(\mu_{1}, t\right) L\left(\mu_{2}, t\right)$ and

$$
\varepsilon\left(I\left(\mu_{1}, \mu_{2}\right), \psi, t\right)=\varepsilon\left(\mu_{1}, \psi, t\right) \varepsilon\left(\mu_{2}, \psi, t\right) q_{v}^{-r(\psi)}
$$

If $\mu_{2} / \mu_{1}=|\cdot|_{v}$, then

$$
\begin{gathered}
L\left(\operatorname{St}\left(\mu_{1}|\cdot|{ }_{v}^{-1 / 2}, \mu_{1}|\cdot|_{v}^{1 / 2}\right), t\right)=L\left(\mu_{1}|\cdot|{ }_{v}^{1 / 2}, t\right) \\
\varepsilon\left(\operatorname{St}\left(\mu_{1}|\cdot|_{v}^{-1 / 2}, \mu_{1}|\cdot|{ }_{v}^{1 / 2}\right), \psi, t\right)=\frac{L\left(\mu_{1}^{-1}, t^{-1}\right)}{L\left(\mu_{1}, t\right)} \varepsilon\left(\mu_{1}, \psi, t\right) \varepsilon\left(\mu_{1}|\cdot|_{v}, \psi, t\right) q_{v}^{-r(\psi)}
\end{gathered}
$$

If $\pi$ is a cuspidal representation of $\mathrm{GL}\left(2, F_{v}\right)$ then $L(\pi, t)$ is 1.
Recall that an irreducible admissible infinite dimensional representation $\pi$ of $\mathrm{GL}\left(2, F_{v}\right)$ on a vector space $V$ is called unramified if its space $V^{K}$ of $K=$ $\mathrm{GL}\left(2, O_{v}\right)$-fixed vectors is nonzero. In this case $V^{K}$ is one dimensional, and $\pi=I\left(\mu_{1}, \mu_{2}\right)$ with unramified $\mu_{1}, \mu_{2}$ and $\mu_{1} / \mu_{2} \neq|\cdot|^{ \pm 1}$.

Corollary 4.18. Let $\pi$ be an unramified irreducible admissible infinite dimensional representation of $\operatorname{GL}\left(2, F_{v}\right)$ and $\psi \neq 1$ with $r(\psi)=0$. Then $\varepsilon(\pi, \psi, t)=1$.

Proof. Here $\pi=I\left(\mu_{1}, \mu_{2}\right)$ with unramified $\mu_{1}, \mu_{2}$, so the claim follows from the last proposition and Tate's Thesis.

Let $\pi$ be an admissible irreducible representation of $\mathrm{GL}(2, \mathbb{A})$ whose local components are all infinite dimensional. Put $L(\pi, t)=\prod_{v} L\left(\pi_{v}, t_{v}\right), t_{v}=t^{\operatorname{deg}(v)}$; the infinite product converges in $E[[t]]$. For any character $\psi: \mathbb{A} / F \rightarrow E^{\times}, \psi \neq$ 1 , put $\varepsilon(\pi, \psi, t)=\prod_{v} \varepsilon\left(\pi_{v}, \psi_{v}, t_{v}\right)$; almost all factors here are 1. From (7) it follows that if the central character of $\pi$ is trivial on $F^{\times}$, then $\varepsilon(\pi, \psi, t)$ is independent of the choice of $\psi: \mathbb{A} / F \rightarrow E^{\times}$. We denote it in this case by $\varepsilon(\pi, t)$.
Theorems 11.1, 11.3 of [JL70] assert:
Theorem 4.19. Let $\pi$ be an irreducible admissible representation of GL(2, $\mathbb{A}$ ) over $E$. Denote by $\eta: \mathbb{A}^{\times} \rightarrow E^{\times}$its central character. Then $\pi$ is cuspidal iff (1) $\eta$ is trivial on $F^{\times}$; (2) all local components of $\pi$ are infinite dimensional; (3) for any character $\omega: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$, the formal series $L(\pi \omega, t)$ is a polynomial in $t$, and (4) $L(\pi \omega, t)=\varepsilon(\pi \omega, t) L\left(\pi \eta^{-1} \omega^{-1}, q^{-2} t^{-1}\right)$.
Note that (4) makes sense due to (3). In [JL70], (3) is formulated as stating that the product $\prod_{v} L\left(\pi_{v} \omega_{v}, t_{v}\right)$ converges absolutely for sufficiently small $t$, and its value has an analytic continuation to a holomorphic function in $\mathbb{C}-\{0\}$. But the argument of [JL70] can be modified to lead to (3) in our case of $E$ which is not $\mathbb{C}$, over a function field $F$. Note that (4) is not $\prod_{v} \Gamma\left(\pi_{v} \omega_{v}, \psi_{v}, t_{v}\right)=1$; indeed the product here does not converge.
Proposition 4.20. If $\pi, \pi^{\prime}$ are cuspidal representations of $\mathrm{GL}(2, \mathbb{A})$ and $\pi_{v} \simeq$ $\pi_{v}^{\prime}$ for almost all $v$, then $\pi \simeq \pi^{\prime}$.

Proof. Let $S$ be a finite set of closed points of $X$ with $\pi_{v} \simeq \pi_{v}^{\prime}$ at $v \notin S$. Let $\eta, \eta^{\prime}$ be the central characters of $\pi, \pi^{\prime}$, and $\eta_{v}, \eta_{v}^{\prime}$ their components at $v$ (restrictions to $F_{v}^{\times}$). By our assumption, $\eta_{v}^{\prime}=\eta_{v}$ for all $v \notin S$. But the groups $F_{v}^{\times}, v \notin S$, generate a dense subgroup of $\mathbb{A}^{\times} / F^{\times}$. Hence $\eta^{\prime}=\eta$. By the Theorem 4.19, of [JL70], above, fixing a character $\psi: \mathbb{A} / F \rightarrow E^{\times}, \psi \neq 1$, for any character $\omega: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$one has

$$
\begin{aligned}
\prod_{v} L\left(\pi_{v} \omega_{v}, t_{v}\right) & =\prod_{v} \varepsilon\left(\pi_{v} \omega_{v}, \psi_{v}, t_{v}\right) L\left(\pi_{v} \eta_{v}^{-1} \omega_{v}^{-1}, q_{v}^{-2} t_{v}^{-1}\right), \\
\prod_{v} L\left(\pi_{v}^{\prime} \omega_{v}, t_{v}\right) & =\prod_{v} \varepsilon\left(\pi_{v}^{\prime} \omega_{v}, \psi_{v}, t_{v}\right) L\left(\pi_{v}^{\prime} \eta_{v}^{\prime-1} \omega_{v}^{-1}, q_{v}^{-2} t_{v}^{-1}\right) .
\end{aligned}
$$

Since $\pi_{v} \simeq \pi_{v}^{\prime}$ at all $v \notin S$, we conclude

$$
\begin{aligned}
& \prod_{v \in S} \Gamma\left(\pi_{v} \omega_{v}, \psi_{v}, t_{v}\right)=\prod_{v \in S} \frac{\varepsilon\left(\pi_{v} \omega_{v}, \psi_{v}, t_{v}\right) L\left(\pi_{v} \eta_{v}^{-1} \omega_{v}^{-1}, q_{v}^{-2} t_{v}^{-1}\right)}{L\left(\pi_{v} \omega_{v}, t_{v}\right)} \\
= & \prod_{v \in S} \frac{\varepsilon\left(\pi_{v}^{\prime} \omega_{v}, \psi_{v}, t_{v}\right) L\left(\pi_{v}^{\prime} \eta_{v}^{\prime-1} \omega_{v}^{-1}, q_{v}^{-2} t_{v}^{-1}\right)}{L\left(\pi_{v}^{\prime} \omega_{v}, t_{v}\right)}=\prod_{v \in S} \Gamma\left(\pi_{v}^{\prime} \omega_{v}, \psi_{v}, t_{v}\right) .
\end{aligned}
$$

Since $\eta=\eta^{\prime}$, it follows from Proposition 4.16 that for each $v \in S$ there exists $m_{v}>0$ such that if $\chi: F_{v}^{\times} \rightarrow E^{\times}$is any character whose conductor is $\geq m_{v}$,
then $\Gamma\left(\pi_{v} \chi, \psi_{v}, t\right)=\Gamma\left(\pi_{v}^{\prime} \chi, \psi_{v}, t\right)$. Fix $v \in S$ and a character $\chi$ of $F_{v}^{\times}$. By Proposition 4.14, it suffices to show $\Gamma\left(\pi_{v} \chi, \psi_{v}, t\right)=\Gamma\left(\pi_{v}^{\prime} \chi, \psi_{v}, t\right)$. For this, it suffices to choose a character $\omega: \mathbb{A}^{\times} / F^{\times} \rightarrow E^{\times}$in the last displayed equation with $\omega_{v}=\chi$ and such that for each $u \in S-\{v\}$, the conductor of $\omega_{u}$ is bigger than $m_{u}$. But the group $H=F_{v}^{\times} \prod_{u \in S-\{v\}} O_{u}^{\times}$maps isomorphically and homeomorphically onto its image in $\mathbb{A}^{\times} / F^{\times}$. Hence any character of $H$ extends to a character of $\mathbb{A}^{\times} / F^{\times}$.

Proposition 4.21. Let $\eta$ be a character of $\mathbb{A}^{\times} / F^{\times}$, $S$ a finite set of closed points of $X, \psi \neq 1$ a character of $\mathbb{A} / F$ with $r\left(\psi_{u}\right)=0$ for all $u$ in $S$. Suppose that for any closed point $v \in|X|-S, \pi_{v}$ is an irreducible admissible infinite dimensional representation of $\mathrm{GL}\left(2, F_{v}\right)$ with central character $\eta_{v}$ such that almost all $\pi_{v}$ are unramified, there is no pair $\mu_{1}, \mu_{2}$ of characters of $\mathbb{A}^{\times} / F^{\times}$with $\pi_{v}=\pi\left(\mu_{1 v}, \mu_{2 v}\right)$ for almost all $v \in|X|-S$, and for any character $\omega$ of $\mathbb{A}^{\times} / F^{\times}$ which is unramified at all points of $S$, the formal series $\prod_{v \notin S} L\left(\pi_{v} \omega_{v}, t_{v}\right)$ and $\prod_{v \notin S} L\left(\pi_{v} \eta_{v}^{-1} \omega_{v}^{-1}, t_{v}\right)$ are polynomials, and there exists a number $c \in E^{\times}$and integers $n_{u}>0(u \in S)$ such that

$$
\prod_{v \notin S} L\left(\pi_{v} \omega_{v}, t_{v}\right)=c \prod_{u \in S}\left(\omega\left(\boldsymbol{\pi}_{u}\right) t_{u}\right)^{n_{u}} \prod_{v \notin S} \varepsilon\left(\pi_{v} \omega_{v}, \psi_{v}, t_{v}\right) L\left(\pi_{v} \eta_{v}^{-1} \omega_{v}^{-1}, q_{v}^{-2} t_{v}^{-1}\right) .
$$

Then there exists a cuspidal representation $\pi$ of $\mathrm{GL}(2, \mathbb{A})$ with central character $\eta$ such that for every $v \in|X|-S$ the local component of $\pi$ at $v$ is $\pi_{v}$.

A proof is in [JL70], Theorem 11, Corollary 11.6, proof of Theorem 12.2. The representation $\pi$ is unique by Proposition 4.20.

### 4.4 Intertwining Again

We can now return to the study of the intertwining operators.
Proposition 4.22. Let $\mu_{1}, \mu_{2}$ be characters of $F_{v}^{\times}$. Let $\psi \neq 1$ be a character of $F_{v}$. Then

$$
R\left(\mu_{1}, \mu_{2}, t\right) R\left(\mu_{2}, \mu_{1}, t^{-1}\right)=\varepsilon\left(\frac{\mu_{1}}{\mu_{2}}, \psi, q_{v}^{-1} t^{2}\right) \varepsilon\left(\frac{\mu_{2}}{\mu_{1}}, \psi, q_{v}^{-1} t^{-2}\right)
$$

Proof. By the transformation formula for the $\varepsilon$-factors, the right hand side does not depend on $\psi$. We then choose $\psi$ with $\operatorname{ker} \psi \supset O_{v}$ and $\operatorname{ker} \psi \not \supset \pi_{v}^{-1} O_{v}$. We can rewrite the asserted equality as

$$
M\left(\mu_{1}, \mu_{2}, t\right) M\left(\mu_{2}, \mu_{1}, t^{-1}\right)=\Gamma\left(\frac{\mu_{2}}{\mu_{1}}, \psi, q_{v}^{-1} t^{2}\right) \Gamma\left(\frac{\mu_{2}}{\mu_{1}}, \psi, q_{v}^{-1} t^{-2}\right)
$$

The restriction map $I\left(\mu_{1}, \mu_{2}\right) \rightarrow I\left(\mu_{1} / \mu_{2}\right)$, where

$$
I(\mu)=\left\{f \in C^{\infty}\left(\mathrm{SL}\left(2, F_{v}\right)\right) ; f\left(\left(\begin{array}{cc}
a & b \\
0 & 1 / a
\end{array}\right) x\right)=\mu(a)|a|_{v} f(x)\right\},
$$

is an isomorphism ( $\mu: F_{v}^{\times} \rightarrow E^{\times}$is a character). The group $\operatorname{SL}\left(2, F_{v}\right)$ acts transitively on $F_{v}^{2}-\{(0,0)\}$ on the right. The stabilizer of the vector $(0,1)$ is $N\left(F_{v}\right)$. Then $N\left(F_{v}\right) \backslash \mathrm{SL}\left(2, F_{v}\right)$ can be identified with $F_{v}^{2}-\{(0,0)\}$ by $\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \mapsto$ $(c, d) \in F_{v}^{2}-\{(0,0)\}$. Using this we identify $I(\mu)$ with

$$
\begin{gathered}
V(\mu)=\left\{f \in C^{\infty}\left(F_{v}^{2}-\{(0,0)\}\right)\right. \\
\left.f(a x)=\mu(a)^{-1}|a|_{v}^{-1} f(x), a \in F_{v}^{\times}, x \in F_{v}^{2}-\{(0,0)\}\right\}
\end{gathered}
$$

so $I\left(\mu_{1}, \underline{\mu_{2}}\right)$ with $V\left(\mu_{1} / \mu_{2}\right)$. The operator $M\left(\mu_{1}, \mu_{2}, t\right)$ corresponds to the operator $\bar{M}\left(\mu_{1} / \mu_{2}, t^{2}\right)$ where

$$
\bar{M}(\mu, s): V\left(\mu \nu_{s}\right) \rightarrow V\left(\mu^{-1} \nu_{s^{-1}}\right), \quad(\bar{M}(\mu, s) f)(x)=\int_{\{y ; x \wedge y=1\}} f(y) d y
$$

Here $\wedge$ denotes the symplectic form $(a, b) \wedge(c, d)=a d-b c$ on $F_{v}^{2}$. The measure on the line $\ell_{x}=\left\{y \in F_{v}^{2} ; x \wedge y=1\right\}$ is transferred from the Haar measure on $F_{v}$ via the map $F_{v} \rightarrow \ell_{x}$ given by $a \mapsto y_{0}+a x$ where $y_{0}$ is a fixed point on $\ell_{x}$. So we need to show:

$$
\bar{M}(\mu, s) \bar{M}\left(\mu^{-1}, s^{-1}\right)=\Gamma\left(\mu, \psi, q_{v}^{-1} s\right) \Gamma\left(\mu^{-1}, \psi, q_{v}^{-1} s^{-1}\right)
$$

For sufficiently small $s \in \mathbb{C}^{\times}$define operators $A_{s}: C_{c}^{\infty}\left(F_{v}^{2}\right) \rightarrow V\left(\mu \nu_{s}\right)$ and $B_{s}: C_{c}^{\infty}\left(F_{v}^{2}\right) \rightarrow V\left(\mu^{-1} \nu_{s}\right)$ by

$$
\left(A_{s} f\right)(x)=\int_{F_{v}} f(a x) \mu(a) \nu_{s}(a) d a, \quad\left(B_{s} f\right)(x)=\int_{F_{v}} f(a x) \mu(a)^{-1} \nu_{s}(a) d a .
$$

Restriction defines an isomorphism $V\left(\mu \nu_{s}\right) \rightarrow V_{0}(\mu)$, where

$$
\begin{gathered}
V_{0}(\mu)=\left\{f \in C^{\infty}\left(O_{v}^{2}-\{(0,0)\}\right)\right. \\
\left.f(a x)=\mu(a)^{-1} f(x), x \in O_{v}^{2}-\{(0,0)\}, a \in O_{v}^{\times}\right\},
\end{gathered}
$$

so we can identify the spaces $V\left(\mu \nu_{s}\right)$ as $s$ varies.
The operators $A_{s}$ and $B_{s}$, defined above for small $s$, depend rationally on $s$. Hence they can be extended to all $s$.
Consider the Fourier transform

$$
F: C_{c}^{\infty}\left(F_{v}^{2}\right) \rightarrow C_{c}^{\infty}\left(F_{v}^{2}\right), \quad(F f)(y)=\int_{F_{v}^{2}} f(x) \psi(x \wedge y) d x
$$

Lemma 4.23. We have $\bar{M}(\mu, s) A_{s}=\Gamma\left(\mu^{-1}, \psi, q_{v}^{-1} s^{-1}\right) B_{s^{-1}} F$,

$$
\bar{M}\left(\mu^{-1}, s^{-1}\right) B_{s^{-1}}=\Gamma\left(\mu, \psi, q_{v}^{-1} s\right) A_{s} F .
$$

Proof. Given $f \in C_{c}^{\infty}\left(F_{v}^{2}\right), x \in F_{v}^{2}-\{(0,0)\}$, we first show

$$
\Gamma\left(\mu^{-1}, \psi, q_{v}^{-1} s^{-1}\right)\left(B_{s^{-1}} F f\right)(x)=\left(\bar{M}(\mu, s) A_{s} f\right)(x) .
$$

The operators $F, A_{s}, B_{s}$ commute with the action of $\mathrm{SL}\left(2, F_{v}\right)$. This action is transitive on $F_{v}^{2}-\{(0,0)\}$, so we may assume $x=(0,1)$. We compute

$$
\begin{aligned}
\left(B_{s^{-1}} F f\right)((0,1)) & =\int_{F_{v}}(F f)((0, a)) \mu(a)^{-1} \nu_{s^{-1}}(a) d a \\
(F f)((0, a)) & =\int_{F_{v}^{2}} f(y, z) \psi(y a) d y d z=\hat{\varphi}(-a), \\
\hat{\varphi}(a) & =\int \varphi(y) \psi(-y a) d y, \quad \varphi(y)=\int f(y, z) d z
\end{aligned}
$$

Tate's functional equation (see [L], VII, section 3-4) is

$$
\Gamma\left(\mu^{-1}, \psi, q_{v}^{-1} s^{-1}\right) \int \hat{\varphi}(a) \mu^{-1}(a) \nu_{s^{-1}}(a) d a=\int \varphi(y) \mu(y) \nu_{s}(y) \frac{d y}{|y|}
$$

(Formally this can be deduced from the definition of the $\Gamma$-function and the inversion formula $\varphi(y)=\int \hat{\varphi}(a) \psi(a y) d a$. However the left side converges for large $|s|$, while the right for small $|s|$, so one has to show both sides are rational in $s$ ).
We conclude that the left side of the equation to be shown is

$$
\int \varphi(y) \mu(-y) \nu_{s}(y)|y|^{-1} d y=\iint f(y, z) \mu(-y) \nu_{s}(y)|y|^{-1} d y d z
$$

while the right side is (recall: $x=(0,1)$, so $(0,1) \wedge(y, z)=-y)$

$$
\int\left(A_{s} f\right)(-1, z) d z=\iint f(-y, y z) \mu(y) \nu_{s}(y) d y d z
$$

The proof of the second identity of the lemma is similar.
The inverse Fourier transform coincides with $F$ since the form $(x, y) \mapsto x \wedge y$ in the definition of $F$ is skew-symmetric. Hence $F^{2}=1$, and it follows from the Lemma that

$$
\bar{M}(\mu, s) \bar{M}\left(\mu^{-1}, s^{-1}\right) B_{s^{-1}}=\Gamma\left(\mu, \psi, q_{v}^{-1} s\right) \Gamma\left(\mu^{-1}, \psi, q_{v}^{-1} s^{-1}\right) B_{s^{-1}}
$$

However, the operator $B_{s^{-1}}$ is onto for those $s$ where it is defined (even its restriction to $C_{c}^{\infty}\left(F_{v}^{2}-\{(0,0)\}\right)$ is onto), as $V\left(\mu \nu_{s}\right)$ is irreducible, so the proposition follows.

Proposition 4.24. For any characters $\mu_{1}, \mu_{2}$ of $\mathbb{A}^{\times} / F^{\times}$we have

$$
M\left(\mu_{1}, \mu_{2}, t\right) M\left(\mu_{2}, \mu_{1}, t^{-1}\right)=1
$$

Proof. From Proposition 4.21, $M\left(\mu_{1}, \mu_{2}, t\right) M\left(\mu_{2}, \mu_{2}, t^{-1}\right)$ is equal to

$$
q^{2-2 g} m\left(\mu_{1} / \mu_{2}, t^{2}\right) m\left(\mu_{2} / \mu_{1}, t^{-2}\right) R\left(\mu_{1}, \mu_{2}, t\right) R\left(\mu_{2}, \mu_{1}, t^{-1}\right)
$$

while Proposition 4.22 implies, for any character $\psi \neq 1$ of $\mathbb{A} / F$, that

$$
R\left(\mu_{1}, \mu_{2}, t\right) R\left(\mu_{2}, \mu_{1}, t^{-1}\right)
$$

is

$$
\begin{aligned}
\prod_{v} & {\left[\varepsilon\left(\mu_{1 v} / \mu_{2 v}, \psi_{v}, q_{v}^{-1} t_{v}^{2}\right) \varepsilon\left(\mu_{2 v} / \mu_{1 v}, \psi_{v}, q_{v}^{-1} t_{v}^{-2}\right)\right] } \\
& =q^{2 g-2} \varepsilon\left(\mu_{1} / \mu_{2}, q^{-1} t^{2}\right) \varepsilon\left(\mu_{2} / \mu_{1}, q^{-1} t^{-2}\right)
\end{aligned}
$$

As $\varepsilon(\chi, t)=q^{1-g} \prod_{v} \varepsilon\left(\chi_{v}, \psi_{v}, t_{v}\right)$ satisfies the functional equation $L(\chi, t)=$ $\varepsilon(\chi, t) L\left(\chi^{-1}, q^{-1} t^{-1}\right)$, we have that

$$
\varepsilon\left(\mu_{1} / \mu_{2}, q^{-1} t^{2}\right) \varepsilon\left(\mu_{2} / \mu_{1}, q^{-1} t^{-2}\right) m\left(\mu_{1} / \mu_{2}, t^{2}\right) m\left(\mu_{2} / \mu_{1}, t^{-2}\right)
$$

which is equal to

$$
\frac{\varepsilon\left(\mu_{1} / \mu_{2}, q^{-1} t^{2}\right) L\left(\mu_{2} / \mu_{1}, t^{2}\right)}{L\left(\mu_{1} / \mu_{2}, q^{-1} t^{2}\right)} \cdot \frac{\varepsilon\left(\mu_{1} / \mu_{2}, q^{-1} t^{-2}\right) L\left(\mu_{1} / \mu_{2}, t^{2}\right)}{L\left(\mu_{2} / \mu_{1}, q^{-1} t^{-2}\right)}
$$

is equal to 1 .

## $4.5 M^{2}=1$ VIA Mellin transform

We shall next study the relationship between $M: C_{c}^{\infty}\left(Y_{\alpha}\right) \rightarrow C^{\infty}\left(Y_{\alpha}\right)$ and $M\left(\mu_{1}, \mu_{2}, t\right): I\left(\mu_{1} \nu^{t}, \mu_{2} \nu^{-t}\right) \rightarrow I\left(\mu_{2} \nu^{-t}, \mu_{1} \nu^{t}\right)$, and conclude that $M^{2}=1$. Both are defined by the same integral formula. Here $\mu_{1}, \mu_{2}$ are characters of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$. Put

$$
\eta\left(\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right)\right)=\mu_{1}(a) \mu_{2}(b)|a / b|^{1 / 2} \nu_{t}(a / b), \quad \eta: A(\mathbb{A}) / A(F) \cdot \alpha^{\mathbb{Z}} \rightarrow E^{\times} .
$$

It is a character. Recall that $Y_{\alpha}=\alpha^{\mathbb{Z}} N(\mathbb{A}) A(F) \backslash \mathrm{GL}(2, \mathbb{A})$ and $(M f)(x)=$ $\int_{N(\mathbb{A})} f\left(\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x\right) d n$. Suppose that $f \in C_{c}^{\infty}\left(Y_{\alpha}\right)$, and $t \in E^{\times}$. Define a function $T\left(f, \mu_{1}, \mu_{2}, t\right): \operatorname{GL}(2, \mathbb{A}) \rightarrow \mathbb{C}$ by

$$
\left(T\left(f, \mu_{1}, \mu_{2}, t\right)\right)(x)=\int_{\alpha^{\mathbb{Z}} A(F) \backslash A(\mathbb{A})} f\left(a^{-1} x\right) \eta(a) d^{\times} a
$$

Then $T\left(f, \mu_{1}, \mu_{2}, t\right) \in I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{-t}\right)$ is called the Mellin transform of $f$. The notation $T$ can be used also when $f \in C^{\infty}\left(Y_{\alpha}\right)$ is not compactly supported, whenever the integral converges.

Proposition 4.25. For $\varphi \in C_{c}^{\infty}\left(Y_{\alpha}\right)$, characters $\mu_{1}, \mu_{2}: \mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}} \rightarrow$ $E^{\times}$and large enough $t \in \mathbb{C}^{\times}$, the integral defining $T$ converges, and $T\left(M \varphi, \mu_{1}, \mu_{2}, t\right)=M\left(\mu_{2}, \mu_{1}, t^{-1}\right) T\left(\varphi, \mu_{2}, \mu_{1}, t^{-1}\right)$.
Proof. By definition,

$$
\left(T\left(f, \mu_{1}, \mu_{2}, t\right)\right)(x)=\iint f\left(\left(\begin{array}{cc}
a & 0 \\
0 & b
\end{array}\right)^{-1} x\right) \mu_{1}(a) \mu_{2}(b)|a / b|^{1 / 2} \nu_{t}(a / b) d^{\times} a d^{\times} b .
$$

Put $f=M \varphi$, so $f\left(\left(\begin{array}{ll}a & 0 \\ 0 & b\end{array}\right)^{-1} x\right)=|b / a| \int_{N(\mathbb{A})} \varphi\left(\left(\begin{array}{cc}b & 0 \\ 0 & a\end{array}\right)^{-1}\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x\right) d n$. Hence $\left(T\left(f, \mu_{1}, \mu_{2}, t\right)\right)(x)$ equals

$$
\begin{aligned}
& \iiint \varphi\left(\left(\begin{array}{cc}
b & 0 \\
0 & a
\end{array}\right)^{-1}\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) \mu_{1}(a) \mu_{2}(b)|b / a|^{1 / 2} \nu_{t}(a / b) d^{\times} a d^{\times} b d n \\
= & \int_{N(\mathbb{A})}\left(T\left(\varphi, \mu_{2}, \mu_{1}, t^{-1}\right)\right)\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) d n \\
= & \left(M\left(\mu_{2}, \mu_{1}, t^{-1}\right) T\left(\varphi, \mu_{2}, \mu_{1}, t^{-1}\right)\right)(x) .
\end{aligned}
$$

If $t$ is large enough, the integral which defines $M\left(\mu_{2}, \mu_{1}, t^{-1}\right)$ converges, and so is the integral which defines $T\left(f, \mu_{1}, \mu_{2}, t\right)$, which justifies the computation.

Proposition 4.26. If $\varphi \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ then $M \varphi \in C^{\infty}\left(Y_{\alpha}\right)$. If $M \varphi \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ then $M^{2} \varphi=\varphi$.

Proof. Put $f=M \varphi$ and $h=M f=M^{2} \varphi$ ( $h$ is defined if $f \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ ). By Proposition 4.25,

$$
\begin{aligned}
& T\left(h, \mu_{1}, \mu_{2}, t\right)=M\left(\mu_{2}, \mu_{1}, t^{-1}\right) T\left(f, \mu_{2}, \mu_{1}, t^{-1}\right) \\
& T\left(f, \mu_{2}, \mu_{1}, t^{-1}\right)=M\left(\mu_{1}, \mu_{2}, t\right) T\left(\varphi, \mu_{1}, \mu_{2}, t\right)
\end{aligned}
$$

The first equation holds only for large enough $t$, and the second only for small enough $t$. However, both sides of the second equality depend rationally on $t$ (for the left side, this is true since $f=M \varphi$ is compactly supported), hence it holds for all $t$ in $\mathbb{C}^{\times}$. Hence for large enough $t$, by Proposition $4.24 T\left(h, \mu_{1}, \mu_{2}, t\right)=$ $T\left(\varphi, \mu_{1}, \mu_{2}, t\right)$ for all $\mu_{1}, \mu_{2}$. This implies $h=\varphi$.

### 4.6 Poles, zeroes and values of $R$ and $M$

Recall that $\nu_{t}(x)=t^{\operatorname{deg}(x)}$ is a character of $\mathbb{A}^{\times} / F^{\times}$with $\nu_{t}\left(\boldsymbol{\pi}_{v}\right)=t_{v}\left(=t^{\operatorname{deg}(v)}\right)$, and locally we write $\nu_{t}$ for the unramified character of $F_{v}^{\times}$with $\nu_{t}\left(\boldsymbol{\pi}_{v}\right)=t$. Let $\mu_{1}, \mu_{2}$ be characters of $F_{v}^{\times}$. Recall:

$$
R\left(\mu_{1}, \mu_{2}, t\right)=\frac{L\left(\mu_{1} / \mu_{2}, q_{v}^{-1} t^{2}\right)}{L\left(\mu_{1} / \mu_{2}, t^{2}\right)} M\left(\mu_{1}, \mu_{2}, t\right)
$$

Proposition 4.27. (1) The function $R\left(\mu_{1}, \mu_{2}, t\right)$ is regular at $t=0$.
It has a pole at $\tau \in \mathbb{C}^{\times}$iff $\mu_{2} \nu_{\tau^{-1}} / \mu_{1} \nu_{\tau}=\nu\left(\right.$ with $\left.\nu\left(\boldsymbol{\pi}_{v}\right)=q_{v}^{-1}\right)$. This pole has order 1.
The function $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ has a pole at $\tau \in \mathbb{C}^{\times}$iff $\mu_{1} \nu_{\tau} / \mu_{2} \nu_{\tau^{-1}}=\nu$. This pole has order 1 .
(2) Suppose $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ has a pole at $\tau \in \mathbb{C}^{\times}$. Then the function $R\left(\mu_{1}, \mu_{2}, t\right)$ is regular at $t=\tau$. Put $L=\lim _{t \rightarrow \tau}(t-\tau) R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ and $Q=R\left(\mu_{1}, \mu_{2}, \tau\right)$. The operators $Q: I\left(\mu_{1} \nu_{\tau}, \mu_{2} \nu_{\tau^{-1}}\right) \rightarrow I\left(\mu_{2} \nu_{\tau^{-1}}, \mu_{1} \nu_{\tau}\right)$ and $L: I\left(\mu_{2} \nu_{\tau^{-1}}, \mu_{1} \nu_{\tau}\right) \rightarrow I\left(\mu_{1} \nu_{\tau}, \mu_{2} \nu_{\tau^{-1}}\right)$ intertwine the GL $\left(2, F_{v}\right)$-action. The representations of $\mathrm{GL}\left(2, F_{v}\right)$ in the spaces $\operatorname{ker} Q, \operatorname{coker} Q, \operatorname{im} L$ are isomorphic
to the square integrable $\operatorname{St}\left(\mu_{1} \nu_{\tau}, \mu_{2} \nu_{\tau^{-1}}\right)$. The representations of $\mathrm{GL}\left(2, F_{v}\right)$ in the spaces $\operatorname{ker} L$, coker $L, \operatorname{im} Q$ are isomorphic to the one dimensional $x \mapsto \mu_{2}(x)\left(\nu \nu_{\tau^{-1}}\right)(x)=\mu_{1}(x) \nu_{\tau}(x)$.
(3) The statement (2) remains true with $R\left(\mu_{1}, \mu_{2}, t\right)$ replaced by $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$.

Proof. From the first part of the proof of Proposition 4.3 it follows that

$$
M\left(\mu_{1}, \mu_{2}, t\right) / L\left(\mu_{1} / \mu_{2}, t^{2}\right)=R\left(\mu_{1}, \mu_{2}, t\right) / L\left(\mu_{1} / \mu_{2}, q_{v}^{-1} t^{2}\right)
$$

is regular. So $R\left(\mu_{1}, \mu_{2}, t\right)$ could have a pole at $t \in \mathbb{C}^{\times}$only if $L\left(\mu_{1} / \mu_{2}, q_{v}^{-1} t^{2}\right)$ is $\infty$, that is $\mu_{2} \nu_{\tau^{-1}} / \mu_{1} \nu_{\tau}=\nu$ (recall: $\nu(x)=|x|$ ), and the order of the pole is at most 1 .
A similar statement holds for $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}=c\left(\mu_{1}, \mu_{2}\right) t^{n\left(\mu_{1}, \mu_{2}\right)} R\left(\mu_{2}, \mu_{1}, t^{-1}\right)$. (The last equality follows from Proposition 4.22. In fact $n\left(\mu_{1}, \mu_{2}\right)=0$, but we do not need this.) Namely $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ has a pole at $\tau \in \mathbb{C}^{\times}$iff $\mu_{1} \nu_{\tau} / \mu_{2} \nu_{\tau^{-1}}=\nu$. This pole has order 1 .
Suppose $\mu_{1} \nu_{\tau} / \mu_{2} \nu_{\tau^{-1}}=\nu$. Then $\mu_{2} \nu_{\tau^{-1}} / \mu_{1} \nu_{\tau} \neq \nu$ so that $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ is regular at $t=\tau$. With $L, Q$ defined as in the proposition, it is clear they commute with the GL $\left(2, F_{v}\right)$-action. If $L=0$ then $Q=R\left(\mu_{1}, \mu_{2}, \tau\right)$ has no pole, in fact it is an isomorphism. If $Q=0$ then $L$ would be an isomorphism, as the operator $\lim _{t \rightarrow \tau} R\left(\mu_{1}, \mu_{2}, t\right) /(t-\tau)$ would be the inverse of $L$. However, the representations of $\operatorname{GL}\left(2, F_{v}\right)$ in $I\left(\mu_{1} \nu_{\tau}, \mu_{2} \nu_{\tau^{-1}}\right)$ and $I\left(\mu_{2} \nu_{\tau^{-1}}, \mu_{1} \nu_{\tau}\right)$ are not equivalent, hence $L \neq 0, Q \neq 0$. As $L \neq 0$, the function $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ does have a pole at $t=\tau$. From the description of the invariant subspaces of $I\left(\mu_{1} \nu_{\tau}, \mu_{2} \nu_{\tau^{-1}}\right)$ and $I\left(\mu_{2} \nu_{\tau^{-1}}, \mu_{1} \nu_{\tau}\right)$ the claims in the proposition on the description of the action of $\operatorname{GL}\left(2, F_{v}\right)$ follow. The regularity of $R\left(\mu_{1}, \mu_{2}, t\right)$ at $t=0$ follows from that of $L\left(\mu_{1} / \mu_{2}, q_{v}^{-1} t^{2}\right)^{-1} R\left(\mu_{1}, \mu_{2}, t\right)$.
In conclusion, the representation of $\mathrm{GL}\left(2, F_{v}\right)$ in $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ is reducible iff $R\left(\mu_{1}, \mu_{2}, t\right)$ or $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ has a pole at $t=\tau$. These last operators are regular at $t \in \mathbb{C}^{\times}$if $\mu_{1} / \mu_{2}$ is ramified. If $\mu_{1} / \mu_{2}$ is unramified and $\left(\mu_{1} / \mu_{2}\right)\left(\boldsymbol{\pi}_{v}\right)=$ $a$, then the poles of $R\left(\mu_{1}, \mu_{2}, t\right)$ are at $\pm \sqrt{q_{v} / a}$, and those of $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ are at $\pm \sqrt{a / q_{v}}$.

Corollary 4.28. Let $\mu_{1}$, $\mu_{2}$ be characters of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$. If $R\left(\mu_{1}, \mu_{2}, t\right)$ has a pole at $t=\tau \in \mathbb{C}^{\times}$, then $|\tau|=\sqrt{q}$. If $R\left(\mu_{1}, \mu_{2}, t\right)^{-1}$ has a pole at $t=\tau \in \mathbb{C}^{\times}$ then $|\tau|=q^{-1 / 2}$.

Indeed, a character of $\mathbb{A}^{\times} / F^{\times}$which takes the value 1 at $\alpha$ is unitary, thus $|a|=1$.

Proposition 4.29. Let $\mu_{1}, \mu_{2}$ be characters of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$ and $\tau \in \mathbb{C}^{\times}$, $|\tau| \leq 1$. If $M\left(\mu_{1}, \mu_{2}, t\right)$ has a pole at $t=\tau$ then $\mu_{1}=\mu_{2}$ and $\tau= \pm q^{-1 / 2}$. If $\mu_{1}=\mu_{2}$ is denoted $\mu$ and $\tau= \pm q^{-1 / 2}$ then $M(\mu, \mu, t)$ has an order 1 pole at $\tau$. The image of the operator $C=\lim _{t \rightarrow \tau}(t-\tau) M(\mu, \mu, t)$ in this case is one dimensional and is spanned by the function $f(x)=\mu(\operatorname{det} x) \nu_{\tau}(\operatorname{det} x)$ in $I\left(\mu \nu_{\tau^{-1}}, \mu \nu_{\tau}\right)$. Further, $M\left(\mu_{1}, \mu_{2}, t\right)$ is regular at $t=0$.

Proof. Recall that $M\left(\mu_{1}, \mu_{2}, t\right)=q^{1-g} m\left(\mu_{1} / \mu_{2}, t^{2}\right) R\left(\mu_{1}, \mu_{2}, t\right)$ where $m(\mu, t)=$ $L(\mu, t) / L(\mu, t / q)$. Let $\tau \in \mathbb{C}^{\times},|\tau| \leq 1$. By Corollary 4.28, the function $R\left(\mu_{1}, \mu_{2}, t\right)$ is regular at $\tau$. By Proposition 4.11, the function $m\left(\mu_{1} / \mu_{2}, t^{2}\right)$ is not regular at $\tau$ only if $\mu_{1}=\mu_{2}$ and $\tau= \pm q^{-1 / 2}$. In these cases it has a simple pole. Hence $M\left(\mu_{1}, \mu_{2}, t\right)$ is regular at $t=\tau(0<|\tau| \leq 1)$ unless $\mu_{1}=\mu_{2}$ and $\tau= \pm q^{-1 / 2}$ where the order of the pole is at most 1 . When $\mu_{1}=\mu_{2}=\mu$ and $\tau= \pm q^{-1 / 2}$, the operator $C=\lim _{t \rightarrow \tau}(t-\tau) M(\mu, \mu, t)$ is a scalar multiple of $R(\mu, \mu, t)=\otimes_{v} R\left(\mu_{v}, \mu_{v}, \tau_{v}\right), \tau_{v}=\tau^{\operatorname{deg}(v)}$.
From (1) in Proposition 4.27, the function $R\left(\mu_{v}, \mu_{v}, \tau_{v}\right)^{-1}$ has a pole at $t=\tau$ ( $t_{v}=\tau_{v}$ ). Its statement (2) implies that the image of $R\left(\mu_{v}, \mu_{v}, \tau_{v}\right)$ is one dimensional and GL $\left(2, F_{v}\right)$ acts on it via the character $x \mapsto \mu_{v}(\operatorname{det} x) \nu_{\tau}(\operatorname{det} x)^{\operatorname{deg} v}$. This implies the proposition, except the final claim, which follows from the regularily of $R\left(\mu_{1}, \mu_{2}, t\right)$ at $t=0$, and that of $m\left(\mu_{1} / \mu_{2}, t^{2}\right)$ at $t=0$.

Let $\mu_{1}, \mu_{2}$ be characters of $\mathbb{A}^{\times} / F^{\times}$. The operator $M\left(\mu_{1}, \mu_{2}, t\right)$ maps $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ into the space $I\left(\mu_{2} \nu_{t-1}, \mu_{1} \nu_{t}\right)$, which in general is different from $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$. However, when $\mu_{1}=\mu_{2}=\mu$ and $t= \pm 1$, then $M\left(\mu_{1}, \mu_{2}, t\right)$ maps $I\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}\right)$ to itself; $M(\mu, \mu, t)$ is regular at $t= \pm 1$. The representation of $\mathrm{GL}(2, \mathbb{A})$ in $I\left(\mu \nu_{\tau}, \mu \nu_{\tau^{-1}}\right), \tau= \pm 1$, is irreducible, and hence $M(\mu, \mu, \tau)$ is a scalar operator. Moreover, from Proposition 4.26, $M(\mu, \mu, \tau)^{2}=1$ at $\tau= \pm 1$.
Proposition 4.30. If $\mu$ is a character of $\mathbb{A}^{\times} / F^{\times}$and $\tau= \pm 1$, then $M(\mu, \mu, \tau)=-1$.
Proof. In view of the relation between $M$ and $R$, it suffices to verify that

$$
\lim _{t \rightarrow 1} \frac{L(1, t)}{L(1, t / q)}=-q^{g-1} \quad \text { and } \quad R(\mu, \mu, \tau)=1
$$

In fact, for any character $\omega$ of $F_{v}^{\times}, R(\omega, \omega, \tau)$ is 1 at $\tau= \pm 1$. Indeed, suppose first $\omega$ is unramified. Then there exists a function $f$ in $I\left(\omega \nu_{\tau}, \omega \nu_{\tau}\right)$ whose restriction to $\mathrm{GL}\left(2, O_{v}\right)$ is 1. By the normalization of the intertwining operator (Proposition $4.3(2)), R(\omega, \omega, \tau) f=f$. However, the representation of $\mathrm{GL}\left(2, F_{v}\right)$ on $I\left(\omega \nu_{\tau}, \omega \nu_{\tau}\right)$ is irreducible, so $R(\omega, \omega, \tau)=1$ if $\omega$ is unramified. The general case reduces to the case where $\omega$ is unramified, or even $\omega=1$, by the commutativity of the diagram

$$
\begin{gathered}
I\left(\omega \nu_{\tau}, \omega \nu_{\tau}\right) \\
\uparrow \\
\stackrel{R(\omega, \omega, \tau)}{\longrightarrow} I\left(\omega \nu_{\tau}, \omega \nu_{\tau}\right) \\
\uparrow \\
I\left(\nu_{\tau}, \nu_{\tau}\right) \otimes \omega \\
\xrightarrow{R(1,1, \tau)} I\left(\nu_{\tau}, \nu_{\tau}\right) \otimes \omega
\end{gathered}
$$

To compute the limit of the ratio of $L$-functions, we use the functional equation $L(1, t / q)=\varepsilon(1, t / q) L\left(1, t^{-1}\right)$. Then

$$
\lim _{t \rightarrow 1} L(1, t) / L(1, t / q)=\varepsilon(1,1 / q)^{-1} \lim _{t \rightarrow 1} L(1, t) / L\left(1, t^{-1}\right) .
$$

By the definition of the global $\varepsilon$-function and its properties (Proposition 6.1, 6.3), $\varepsilon(1,1 / q)=q^{1-g}$. Since $L(1, t)$ has a pole of order one at $t=1$, by L'Hôpital rule $\lim _{t \rightarrow 1} L(1, t) / L\left(1, t^{-1}\right)$ is -1 .

### 4.7 Global Eisenstein approach

These proofs of $M^{2}=1$ and rationality of $M\left(\mu_{1}, \mu_{2}, t\right)$ are based on local computations (normalization of the intertwining operators by $L$-functions and $\varepsilon$-factors), and the functional equation of the $L$-function. The following alternative proof of these results is based on properties of the Eisenstein map.
The alternative approach of this subsection, the following subsection 4.8, and the computation of traces in subsection 5.2 are motivated by Tate [T68]. They are the newest part of this paper, which - as noted in the introduction - cries out for generalization from our context of $\operatorname{GL}(2)$, and for further study.
We shall use the maps $\mathrm{ht}^{+}: Y_{\alpha} \rightarrow \mathbb{Z}$ and ht $: \alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) \rightarrow \mathbb{Z}$. Both maps are proper. However, $\mathrm{ht}^{+}$is onto while the image of ht contains the positive integers but only finitely many negatives. So in some sense $Y_{\alpha}$ is less compact than $\alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})$, so the map $E: C_{c}^{\infty}\left(Y_{\alpha}\right) \rightarrow$ $C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right)$ should have a big kernel. For $\varphi$ in $\operatorname{ker} E$ we have $(1+M) \varphi=E^{*} E \varphi=0$. Hence $M^{2} \varphi=\varphi$. Unlike $M$, the operator $M^{2}$ commutes with the action of $A(\mathbb{A})$ on $C_{c}^{\infty}\left(Y_{\alpha}\right)$ by left translation. Hence $M^{2} \varphi=\varphi$ not only for $\varphi \in \operatorname{ker} E$ but also for $\varphi$ in the span of $A(\mathbb{A})$-translates of $\varphi$ in $\operatorname{ker} E$. The number of such linear combinations is already sufficiently large to imply $M^{2}=1$. We now turn to rigorous proofs.

Proposition 4.31. Let $M: \mathbb{C}\left[z, z^{-1}\right]^{n} \rightarrow \mathbb{C}((z))^{n}$ be a $\mathbb{C}$-linear map with $M(z u)=z^{-1} M(u)$ for all $u \in \mathbb{C}\left[z, z^{-1}\right]^{n}$. Let $I$ denote the natural embedding $\mathbb{C}\left[z, z^{-1}\right]^{n} \hookrightarrow \mathbb{C}((z))^{n}$. Put $B=I+M$. Suppose there is some $k \in \mathbb{Z}$ for which the vector space $(\operatorname{Im} B) / B\left(z^{k} \mathbb{C}\left[z^{-1}\right]^{n}\right)$ is finite dimensional. Then there is some

$$
P(z) \in \mathrm{GL}(n, \mathbb{C}(z)) \subset \operatorname{GL}(n, \mathbb{C}((z)))
$$

with $P\left(z^{-1}\right)=P(z)^{-1}$ and $(M u)(z)=P(z) u\left(z^{-1}\right)$ for all $u(z) \in \mathbb{C}\left[z, z^{-1}\right]^{n}$.
Proof. Denote by $e_{i}$ the column in $\mathbb{C}^{n}$ with nonzero entry only at the $i$ th row, where it is 1 . From $M\left(\sum_{i}\left(\sum_{j} c_{i j} z^{j}\right) e_{i}\right)=\sum_{i}\left(\sum_{j} c_{i j} z^{-j}\right) M e_{i}$, we see that $(M u)(z)=P(z) u\left(z^{-1}\right)$ where $P(z)$ is the $n \times n$ matrix with columns $M e_{1}, \ldots, M e_{n}$ whose entries are in $\mathbb{C}((z))$. If $u$ is in the kernel of $B=I+M$, then $P(z) u\left(z^{-1}\right)=-u(z)$. Since $\operatorname{Im} B=\cup_{m \geq 1} B\left(z^{m} \mathbb{C}\left[z^{-1}\right]^{n}\right)$ and there is some $k \geq 0$ such that $B\left(z^{k} \mathbb{C}\left[z^{-1}\right]^{n}\right)$ has finite codimension in $\operatorname{Im} B$, there is some $\ell$ with $B\left(z^{\ell} \mathbb{C}\left[z^{-1}\right]^{n}\right)=\operatorname{Im} B$. Then ker $B+z^{\ell} \mathbb{C}\left[z^{-1}\right]^{n}=\mathbb{C}\left[z, z^{-1}\right]^{n}$. For each $i(1 \leq i \leq n), z^{\ell+1} e_{i} \in \operatorname{ker} B+z^{\ell} \mathbb{C}\left[z^{-1}\right]^{n}$. Hence there is a matrix $W \in$ $M\left(n, \mathbb{C}\left[z, z^{-1}\right]\right)$ whose columnes are in ker $B$ and $W-z^{\ell+1} \operatorname{Id} \in z^{\ell} M\left(n, \mathbb{C}\left[z^{-1}\right]\right)$, where Id is the identity matrix. But then $W \in \operatorname{GL}(n, \mathbb{C}(z))$, and since the columns of $W$ are in ker $B$, we have $P(z) W\left(z^{-1}\right)=-W(z)$. Then $P(z)=$ $-W(z) W\left(z^{-1}\right)^{-1}$, and $P\left(z^{-1}\right)=-W\left(z^{-1}\right) W(z)^{-1}=P(z)^{-1}$.

Corollary 4.32. A $\mathbb{C}$-linear map $M: \mathbb{C}\left[z, z^{-1}\right] \rightarrow \mathbb{C}\left[z, z^{-1}\right]$ which satisfies the conditions of Proposition 4.31 has $M^{2}=\mathrm{Id}$.

Recall that $Y_{\alpha}=\alpha^{\mathbb{Z}} A(F) N(\mathbb{A}) \backslash \mathrm{GL}(2, \mathbb{A})$. Write $C_{+}^{\infty}\left(Y_{\alpha}\right)$ for the space of the $E$-valued functions $f$ on $Y_{\alpha}$ with (1) $f(x)=0$ if ht $^{+}(x)$ is large enough, and (2) $f$ is invariant under right translation by some open subgroup $U$ of GL(2, $\mathbb{A})$.
Note that $C_{c}^{\infty}\left(Y_{\alpha}\right) \subset C_{+}^{\infty}\left(Y_{\alpha}\right) \subset C^{\infty}\left(Y_{\alpha}\right)$.
Proposition 4.33. The image of $C_{c}^{\infty}\left(Y_{\alpha}\right)$ under $M$ lies in $C_{+}^{\infty}\left(Y_{\alpha}\right)$.
Proof. For $f \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ there exists an integer $m$ such that $f(x)=0$ if ht ${ }^{+}(x)<$ $-m$. We shall show that for such $f,(M f)(x)=\int_{N(\mathbb{A})} f\left(\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x\right) d x$ is zero if $\mathrm{ht}^{+}(x)>m$. It suffices to show then that for $x \in \mathrm{GL}(2, \mathbb{A})$ with $\mathrm{ht}^{+}(x)>m$, and any $n \in N(\mathbb{A})$, we have $h^{+}\left(\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x\right)<-m$. But by Lemma 2.7 we have

$$
\mathrm{ht}^{+}(x)+\mathrm{ht}^{+}\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right)=\mathrm{ht}^{+}(n x)+\mathrm{ht}^{+}\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) \leq 0 .
$$

Proposition 4.34. Let $U$ be an open subgroup of $\mathrm{GL}(2, O)$. For every integer $m \geq 1$ define

$$
\begin{gathered}
W_{m}^{U}=\left\{\varphi \in C_{c}^{\infty}\left(Y_{\alpha}\right)^{U} ; \varphi(x)=0 \text { if } \mathrm{ht}^{+}(x)<m\right\} \\
Y_{m}^{U}=\left\{\varphi \in C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right)^{U} ; \varphi(x)=0 \text { if } \mathrm{ht}^{+}(x)<m\right\} .
\end{gathered}
$$

Then $E\left(W_{m}^{U}\right)=Y_{m}^{U}$ for large enough $m$.
Proof. Put

$$
Z_{m}^{U}=\left\{\varphi \in C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} \cdot A(F) N(F) \backslash \mathrm{GL}(2, \mathbb{A})\right)^{U} ; \varphi(x)=0 \text { if } \mathrm{ht}^{+}(x)<m\right\}
$$

Recall that

$$
E=s_{*} r^{*}, \quad s_{*}(x)=\sum_{\gamma} \psi(\gamma x), \quad \gamma \in A(F) N(F) \backslash \mathrm{GL}(2, F)
$$

It is clear that $s_{*}\left(Z_{m}^{U}\right)=Y_{m}^{U}$. It suffices to show that $r^{*}\left(W_{m}^{U}\right)=Z_{m}^{U}$ for sufficiently large $m$. In fact, we showed, as the first claim in the proof of Proposition 2.13, that for an open subgroup $U$ of $\mathrm{GL}(2, \mathbb{A})$, that there is an integer $m$ with the property that if $z \in \mathbb{A}, x \in \mathrm{GL}(2, \mathbb{A})$, $\mathrm{ht}^{+}(x) \geq m$, then there is $u \in U, \beta \in F$, with $\left(\begin{array}{cc}1 & z \\ 0 & 1\end{array}\right) x=\left(\begin{array}{ll}1 & \beta \\ 0 & 1\end{array}\right) x u$. In other words, if $x \in \operatorname{GL}(2, \mathbb{A})$ and $\mathrm{ht}^{+}(x)$ is large enough, then $N(\mathbb{A}) x \subset N(F) x U$.

We shall now give a different proof of Proposition 4.26.
Proposition 4.35. If $\varphi \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ and $M \varphi \in C_{c}^{\infty}\left(Y_{\alpha}\right)$ then $M^{2} \varphi=\varphi$.

Proof. Let us introduce a structure of $\mathbb{C}\left[z, z^{-1}\right]$-module on $C^{\infty}\left(Y_{\alpha}\right)$ by

$$
(z f)(x)=\frac{1}{\sqrt{q}} f\left(\left(\begin{array}{cc}
\alpha & 0 \\
0 & 1
\end{array}\right) x\right), \quad f \in C^{\infty}\left(Y_{\alpha}\right), \quad x \in \mathrm{GL}(2, \mathbb{A}) .
$$

From

$$
(M \phi)\left(\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right) x\right)=\left|\frac{a}{b}\right| \int_{N(\mathbb{A})} \phi\left(\left(\begin{array}{cc}
b & 0 \\
0 & a
\end{array}\right)\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) d n
$$

it follows that $M(z f)=z^{-1} M(f)$; recall that $|\alpha|=q$, and $f$ is invariant under $\alpha$. This is the reason for introducing the factor $\sqrt{q}$. Let $U$ be an open subgroup of GL $(2, O)$. Put

$$
W_{c}^{U}=C_{c}^{\infty}\left(Y_{\alpha}\right)^{U}, \quad W_{+}^{U}=C_{+}^{\infty}\left(Y_{\alpha}\right)^{U} .
$$

Both are $\mathbb{C}\left[z, z^{-1}\right]$-submodules in $C^{\infty}\left(Y_{\alpha}\right)$. Denote by $W_{0}^{U}$ the set of functions $f \in C^{\infty}\left(Y_{\alpha}\right)^{U}$ such that $f(x)=0$ if $\mathrm{ht}^{+}(x) \neq 0$. Then the natural map $W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right] \rightarrow W_{c}^{U}$ is an isomorphism. In the same way we have a canonical isomorphism $W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}((z)) \rightarrow W_{+}^{U}$. The operator

$$
M: W_{c}=C_{c}^{\infty}\left(Y_{\alpha}\right) \rightarrow W_{+}=C_{+}^{\infty}\left(Y_{\alpha}\right)
$$

maps $W_{c}^{U}$ into $W_{+}^{U}$. Hence it defines a map

$$
M: W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right] \rightarrow W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}((z))
$$

satisfying the first condition of Proposition 4.31.
It remains to check the second condition of that Proposition. The space $W_{m}^{U}$ can be identified with $W_{0}^{U} \otimes_{\mathbb{C}} z^{-m} \mathbb{C}\left[z^{-1}\right]$, and then the operator $B=I+$ $M$ is just $E^{*} E$. Thus it suffices to show that for some $m \in \mathbb{Z}$, the space $E^{*} E\left(W_{c}^{U}\right) / E^{*} E\left(W_{m}^{U}\right)$ is finite dimensional. Since $E\left(W_{m}^{U}\right)=Y_{m}^{U}$ for large $m$, and $\{x \in \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \operatorname{ht}(x) \leq m\}$ is compact $\bmod Z(\mathbb{A})$, it follows that the subspace

$$
E\left(W_{m}^{U}\right) \subset C_{c}^{\infty}\left(\alpha^{\mathbb{Z}} \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A})\right)^{U}
$$

has finite codimension. Thus $M$ satisfies both conditions of Proposition 4.31, and our claim follows from Corollary 4.32.

To use Proposition 4.31 to give another proof of the rationality of $M\left(\mu_{1}, \mu_{2}, t\right)$, we take a different view of the Mellin transform and the relationship between the operators $M$ and $M\left(\mu_{1}, \mu_{2}, t\right)$. Let $I_{c}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right)$ be the space of locally constant functions $f: \operatorname{GL}(2, \mathbb{A}) \rightarrow \mathbb{C}\left[z, z^{-1}\right]$ with

$$
f\left(\left(\begin{array}{ll}
a & c \\
0 & b
\end{array}\right) x\right)=\mu_{1}(a) \mu_{2}(b) \nu_{z}(b / a)|a / b|^{1 / 2} f(x)
$$

Let $I_{+}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right)$ be

$$
I_{c}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right) \otimes_{\mathbb{C}\left[z, z^{-1}\right]} \mathbb{C}((z))
$$

The group $\alpha^{\mathbb{Z}} \subset \mathrm{GL}(2, \mathbb{A})$ acts trivially on these $I_{c}$ and $I_{+}$. We put

$$
I_{c}=\oplus I_{c}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right), \quad I_{+}=\oplus I_{+}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right)
$$

where the sums range over all characters $\mu_{1}, \mu_{2}$ of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$.
Proposition 4.36. There exists an isomorphism of $\mathbb{C}((z))$-modules $I_{+} \xrightarrow{\sim}$ $C_{+}^{\infty}\left(Y_{\alpha}\right)$ which is $\mathrm{GL}(2, \mathbb{A})$-equivariant and maps $I_{c}$ to $C_{c}^{\infty}\left(Y_{\alpha}\right)$.

Proof. Define a map $F: I_{+} \rightarrow C_{+}^{\infty}\left(Y_{\alpha}\right)$ by mapping

$$
\varphi=\left\{\varphi_{\mu_{1}, \mu_{2}}\right\} \in I_{+}, \quad \varphi_{\mu_{1}, \mu_{2}} \in I_{c}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right)
$$

to

$$
(F \varphi)(x)=\text { constant term of the formal series } \sum_{\mu_{1}, \mu_{2}} \varphi_{\mu_{1}, \mu_{2}}(x) \in \mathbb{C}((z)),
$$

for any $x \in \mathrm{GL}(2, \mathbb{A})$. The map $F$ is well defined, commutes with the actions of $\mathbb{C}((z))$ and $\mathrm{GL}(2, \mathbb{A})$. The inverse of $F$ exists, as follows. If $\psi \in C_{+}^{\infty}\left(Y_{\alpha}\right)$ then $F^{-1}(\psi)=\left\{\varphi_{\mu_{1}, \mu_{2}}\right\}$ with $\varphi_{\mu_{1}, \mu_{2}} \in I_{+}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right)$ given by

$$
\varphi_{\mu_{1}, \mu_{2}}(x)=\int_{A(\mathbb{A}) / \alpha^{\mathbb{Z}} \cdot A(F)} \psi\left(h^{-1} x\right) \eta(h) d h
$$

where

$$
\eta: A(\mathbb{A}) \rightarrow \mathbb{C}((z))^{\times}, \quad \eta(\operatorname{diag}(a, b))=\mu_{1}(a) \mu_{2}(b) \nu_{z}(a / b) .
$$

The last integral converges in the field $\mathbb{C}((z))$. A base of the topology is given by $z^{n} \mathbb{C}[[z]], n>0$. The map $F$ maps $I_{c}$ to $C_{c}^{\infty}\left(Y_{\alpha}\right)$.

Put $I_{0}=\oplus_{\mu_{1}, \mu_{2}} I_{0}\left(\mu_{1}, \mu_{2}\right)$, with

$$
I_{0}\left(\mu_{1}, \mu_{2}\right)=\left\{f \in C^{\infty}(\operatorname{GL}(2, O)) ; f\left(\left(\begin{array}{cc}
a & c \\
0 & b
\end{array}\right) x\right)=\mu_{1}(a) \mu_{2}(b) f(x)\right\} .
$$

Denote by $M(z)$ the map $I_{0} \rightarrow I_{0}$ which takes $I_{0}\left(\mu_{1}, \mu_{2}\right)$ to $I_{0}\left(\mu_{2}, \mu_{1}\right)$ via $M\left(\mu_{1}, \mu_{2}, z\right)$. We use the isomorphism $F$ to identify the spaces $I_{+}$and $C_{+}^{\infty}\left(Y_{\alpha}\right)$, as well as $I_{c}$ and $C_{c}^{\infty}\left(Y_{\alpha}\right)$. The natural isomorphism

$$
I_{c}\left(\mu_{1} \nu_{z-1}, \mu_{2} \nu_{z}\right) \xrightarrow{\sim} I_{0}\left(\mu_{1}, \mu_{2}\right) \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]
$$

and

$$
I_{+}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right) \xrightarrow{\sim} I_{0}\left(\mu_{1}, \mu_{2}\right) \otimes_{\mathbb{C}} \mathbb{C}((z))
$$

permit us to identify $I_{c}$ and $I_{0} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]$ as well as $I_{+}$and $I_{0} \otimes_{\mathbb{C}} \mathbb{C}((z))$. Thus the map $M: C_{c}^{\infty}\left(Y_{\alpha}\right) \rightarrow C_{+}^{\infty}\left(Y_{\alpha}\right)$ induces an operator

$$
M_{0}: I_{0} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right] \rightarrow I_{0} \otimes_{\mathbb{C}} \mathbb{C}((z))
$$

Proposition 4.37. Regard the elements of $I_{0} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]$ as functions of $z$ with values in $I_{0}$ and the elements of $I_{0} \otimes_{\mathbb{C}} \mathbb{C}((z))$ as formal series in $z$ with coefficients in $I_{0}$. Then for any $u \in I_{0} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]$ one has $\left(M_{0} u\right)(z)=$ $M(z) u\left(z^{-1}\right), M(z)$ is viewed as a formal series in $z$.

Proof. Write $\iota$ for the automorphism of $\mathbb{C}\left[z, z^{-1}\right]$ which maps $z$ to $z^{-1}$. Given a function $f: \operatorname{GL}(2, \mathbb{A}) \rightarrow \mathbb{C}((z))$, denote by $f_{0}$ the function $\mathrm{GL}(2, \mathbb{A}) \rightarrow \mathbb{C}$ such that $f_{0}(x)$ is the constant term of $f(x)$.
Define an operator

$$
M^{\prime \prime}: I_{0} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right] \rightarrow I_{0} \otimes_{\mathbb{C}} \mathbb{C}((z)) \quad \text { by } \quad\left(M^{\prime \prime} u\right)(z)=M(z) u\left(z^{-1}\right)
$$

We claim that $M_{0}=M^{\prime \prime}$. Consider $M^{\prime \prime}$ as a map $I_{c} \rightarrow I_{+}$. We have to show that for every $f \in I_{c}$, we have $F M^{\prime \prime} f=M F f$, for the isomorphism $F: I_{+} \xrightarrow{\sim} C_{+}^{\infty}\left(Y_{\alpha}\right)$. As $I_{c}$ is the sum over $\mu_{1}, \mu_{2}$ of $I_{c}\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right)$, it suffices to consider $f$ in one of these summands.
For $x \in \operatorname{GL}(2, \mathbb{A})$, we have $\left(M^{\prime \prime} f\right)(x)=\int_{N(\mathbb{A})} \iota f\left(\left(\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right) n x\right) d n$. Then

$$
\begin{gathered}
\left(F M^{\prime \prime} f\right)(x)=\left(M^{\prime \prime} f\right)_{0}(x)=\int_{N(\mathbb{A})} f_{0}\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) d n \\
(M F f)(x)=\int_{N(\mathbb{A})} F f\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) d n=\int_{N(\mathbb{A})} f_{0}\left(\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right) n x\right) d n
\end{gathered}
$$

are equal, as required.
4.8 A second proof of the rationality of $M\left(\mu_{1}, \mu_{2}, t\right)$ and of the FUNCTIONAL EQUATION $M\left(\mu_{1}, \mu_{2}, t\right) M\left(\mu_{2}, \mu_{1}, t^{-1}\right)=1$
Let $U, W^{U}, A$ be as in the proof of Proposition 4.35. Then $W^{U}=$ $\oplus_{\mu_{1}, \mu_{2}} W_{\mu_{1}, \mu_{2}}^{U}$, where $W_{\mu_{1}, \mu_{2}}^{U}$ is the space of functions $f \in W^{U}$ with

$$
f\left(\left(\begin{array}{ll}
a & 0 \\
0 & b
\end{array}\right) x\right)=\mu_{1}(a)^{-1} \mu_{2}(b)^{-1} f(x)
$$

whenever $\operatorname{deg}(a)=\operatorname{deg}(b)=0$. The natural maps $I_{0}\left(\mu_{2}, \mu_{1}\right)^{U} \xrightarrow{\sim} W_{\mu_{1}, \mu_{2}}^{U}$ permit one to identify $W^{U}$ and the space $I_{0}^{U}$. The map

$$
M: W^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right] \rightarrow W^{U} \otimes_{\mathbb{C}} \mathbb{C}((z))
$$

is induced by the operator

$$
M_{0}: I_{0} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right] \rightarrow I_{0} \otimes_{\mathbb{C}} \mathbb{C}((z))
$$

The proof of Proposition 4.35 implies that the operator $M$ satisfies the conditions of Proposition 4.31. Then $M$ is given by a formula of the form $(M u)(z)=P(z) u\left(z^{-1}\right)$, where $P(z)$ is an automorphism of $V$ which depends on $z$ rationally, and $P\left(z^{-1}\right)=P(z)^{-1}$. From Proposition 4.37 it follows that
$P(z)$ is just the restriction of $M(z)$ to $I_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]$. The group $U$ may be arbitrarily small. Hence $M(z)$ is a rational function of $z$, and $M(z) M\left(z^{-1}\right)=1$. Hence for any characters $\mu_{1}, \mu_{2}$, of $\alpha^{\mathbb{Z}} \cdot F^{\times} \backslash \mathbb{A}^{\times}$, the operator $M\left(\mu_{1}, \mu_{2}, z\right)$ depends rationally on $z$, and

$$
M\left(\mu_{1}, \mu_{2}, z\right) M\left(\mu_{1}, \mu_{2}, z^{-1}\right)=1
$$

The same is true for any characters $\mu_{1}, \mu_{2}$ of $\mathbb{A}^{\times} / F^{\times}$, which are not necessarily trivial at $\alpha$. To see this, it suffices to use the identities $M\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t}, z\right)=$ $M\left(\mu_{1}, \mu_{2}, z\right)$ and $M\left(\mu_{1} \nu_{t}, \mu_{2} \nu_{t^{-1}}, z\right)=M\left(\mu_{1}, \mu_{2}, t z\right)$.

## 5 Proof of the trace formula

### 5.1 The geometric part

Our aim is to compute the trace $\operatorname{tr} r_{0}(f)$, where $f \in C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$ and $r_{0}$ is the representation of $\mathrm{GL}(2, \mathbb{A})$ by right translation on the space $A_{0, \alpha}$ of cusp forms invariant under $\alpha$. Recall that the space $A_{c, \alpha}$ of $\alpha$-invariant automorphic forms is equal to the direct sum of $A_{0, \alpha}$ and $A_{E, \alpha}=\operatorname{Im}\left(E: C_{c}^{\infty}\left(Y_{\alpha}\right) \rightarrow A_{c, \alpha}\right)$. The corresponding representations of $\mathrm{GL}(2, \mathbb{A})$ are denoted by $r$ and $r_{E}$. Had $r$ been admissible, we would have had $\operatorname{tr} r_{0}(f)=\operatorname{tr} r(f)-\operatorname{tr} r_{E}(f)$, and the computation of $\operatorname{tr} r_{0}(f)$ would have reduced to that of $\operatorname{tr} r(f)$ and $\operatorname{tr} r_{E}(f)$. But $r$ and $r_{E}$ are not admissible, so $\operatorname{tr} r(f)$ and $\operatorname{tr} r_{E}(f)$ make no sense.
Suppose $f$ is right invariant under the open subgroup $U$ of $\operatorname{GL}(2, O)$. Denote by $A_{0}^{U}, A_{c}^{U}, A_{E}^{U}$ the spaces of $U$-invariant vectors in $A_{0, \alpha}, A_{c, \alpha}, A_{E, \alpha}$. Since $\operatorname{Im} r_{0}(f) \subset A_{0}^{U}$, we have $\operatorname{tr} r_{0}(f)=\operatorname{tr} r_{0}^{U}(f)$, where $r_{0}^{U}(f)$ is the restriction of $r_{0}(f)$ to $A_{0}^{U}$.
Denote by $\chi_{m}$ the characteristic function of the set

$$
\left\{x \in \alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathbb{A}) ; \operatorname{ht}(x)<m\right\}, \quad m>0
$$

Denote by $\theta_{m}$ the operator of multiplication by $\chi_{m}$ on $A_{c, \alpha}$.
Proposition 5.1. (1) For any $m>0, \operatorname{dim} \theta_{m}\left(A_{c}^{U}\right)<\infty$.
(2) If $m \gg 1$ then (a) $\theta_{m}$ acts as the identity on $A_{0}^{U}$, and (b) $\theta_{m}\left(A_{E}^{U}\right) \subset A_{E}^{U}$.

Proof. (1) The support of $\chi_{m}$ is compact $\bmod Z(\mathbb{A})$, the quotient by the open $U$ is then finite. (2a) $A_{0}^{U}$ is finite dimensional, consisting of compactly supported forms. (2b) By (2a), $\left(1-\theta_{m}\right) A_{E}^{U}=\left(1-\theta_{m}\right) A_{c}^{U}$. This lies in $A_{E}^{U}$ as $U$-invariant cusp forms are uniformly compactly supported. Hence $\theta_{m}\left(A_{E}^{U}\right) \subset A_{E}^{U}$.

Denote by $r^{U}(f)$ and $r_{E}^{U}(f)$ the restrictions of $r(f)$ to $A_{c}^{U}$ and $A_{E}^{U}$. For $m$ such that $\theta_{m}\left(A_{E}^{U}\right) \subset A_{E}^{U}$, denote the restriction of $\theta_{m}$ to $A_{E}^{U}$ again by $\theta_{m}$. Then for $m \gg 1$,
$\operatorname{tr} r_{0}(f)=\operatorname{tr} r_{0}^{U}(f)=\operatorname{tr}\left(\theta_{m} r^{U}(f)\right)-\operatorname{tr}\left(\theta_{m} r_{E}^{U}(f)\right)=\operatorname{tr}\left(\theta_{m} r(f)\right)-\operatorname{tr}\left(\theta_{m} r_{E}^{U}(f)\right)$.
We then proceed to compute $\operatorname{tr}\left(\theta_{m} r(f)\right)$ and $\operatorname{tr}\left(\theta_{m} r_{E}^{U}(f)\right)$.

Proposition 5.2. There exist $c_{f} \in E$ and $\alpha_{m} \in E$ with $\lim _{n \rightarrow \infty} \alpha_{m}=0$, and

$$
\operatorname{tr}\left(\theta_{m} r(f)\right)=\sum_{1 \leq i \leq 4} S_{i}(f)+c_{f}\left(m-\frac{1}{2}\right)+\alpha_{m}
$$

Proof. The map $\theta_{m} r(f): A_{c, \alpha} \rightarrow A_{c, \alpha}$ is an integral operator with kernel $\chi_{m}(y) K_{f}(x, y)$, where $K_{f}(x, y)=\sum_{\gamma \in \alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F)} f\left(x^{-1} \gamma y\right)$. Then

$$
\operatorname{tr}\left(\theta_{m} r(f)\right)=\int_{\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathrm{~A})} \chi_{m}(x) K_{f}(x, x) d x
$$

Lemma 5.3. There exists $m_{f}>0$ such that if $x \in \operatorname{GL}(2, \mathbb{A}), \gamma \in \alpha^{\mathbb{Z}} \mathrm{GL}(2, F)$, ht $^{+}(x)>m_{f}, f\left(x^{-1} \gamma x\right) \neq 0$, then $\gamma \in \alpha^{\mathbb{Z}} A(F) N(F)$.

Proof. We have $\gamma x=x y, y$ in $\operatorname{supp}(f)$. Since $\mathrm{ht}^{+}(x)+\mathrm{ht}^{+}(\delta x) \leq 0$ for $\delta \in$ $\mathrm{GL}(2, F)-B(F)$, we have that $\mathrm{ht}^{+}(x)>0$. If in addition we had $\mathrm{ht}^{+}(x y)>0$, we would conclude that $\gamma \in \alpha^{\mathbb{Z}} B(F)$. The number $m_{f}=-\min \left\{\operatorname{ht}^{+}(z) ; z \in\right.$ $\mathrm{GL}(2, O) \cdot \operatorname{supp}(f)\}$ then has the property that $\mathrm{ht}^{+}(x)>m_{f}, y \in \operatorname{supp}(f)$, implies $\mathrm{ht}^{+}(x y)=\mathrm{ht}^{+}(x)+\mathrm{ht}^{+}(k y)>0$, where $x=b k$ and $k y=b^{\prime} k^{\prime}$ so that $x y=b b^{\prime} k\left(b, b^{\prime} \in B(\mathbb{A}) ; k, k^{\prime} \in \mathrm{GL}(2, \mathbb{A})\right)$.

Denote by $\xi_{m}$ the characteristic function of the set $\left\{x \in \mathrm{GL}(2, \mathbb{A}) ; \mathrm{ht}^{+}(x) \geq\right.$ $m\}$, by $A^{\prime}(F)$ the set of nonscalar diagonal matrices, and by Ell the set of elliptic matrices in GL $(2, F)$, namely those whose eigenvalues are not in $F$. Put $w=\left(\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right)$.

Lemma 5.4. If $m$ is big enough, then $\chi_{m}(y) K_{f}(x, x)$ is the sum of

$$
\begin{gathered}
T_{1, m}(x)=\chi_{m}(x) \sum_{\gamma \in \alpha^{\mathbb{Z}} \cdot F^{\times}} f(\gamma), \quad T_{2, m}(x)=\sum_{\gamma \in \alpha^{\mathbb{Z}} \cdot \mathrm{Ell}} f\left(x^{-1} \gamma x\right), \\
T_{3, m}(x)=\frac{1}{2} \sum_{\gamma \in \alpha^{Z} \cdot A^{\prime}(F)} \sum_{\delta \in A(F) \backslash \mathrm{GL}(2, F)} f\left(x^{-1} \delta^{-1} \gamma \delta x\right) \cdot\left(1-\xi_{m}(\delta x)-\xi_{m}(w \delta x)\right), \\
T_{4, m}(x)=\sum_{a \in \alpha^{\mathbb{Z}} \cdot F^{\times}} \sum_{\delta \in F^{\times} N(F) \backslash \operatorname{GL}(2, F)} f\left(x^{-1} \delta^{-1}\left(\begin{array}{cc}
a & a \\
0 & a
\end{array}\right) \delta x\right) \cdot\left(1-\xi_{m}(\delta x)\right) .
\end{gathered}
$$

Proof. $T_{1, m}(x)$ is the contribution of the elements $\gamma \in \alpha^{\mathbb{Z}} \cdot F^{\times}$in $\chi_{m}(x) K_{f}(x, x)$. We claim that the contribution of the elements $\gamma \in \alpha^{\mathbb{Z}}$. Ell in $\chi_{m}(x) K_{f}(x, x)$ is $T_{2, m}(x)$. To show this, we need to see that if $x \in \operatorname{GL}(2, \mathbb{A}), \gamma \in \alpha^{\mathbb{Z}} \cdot$ Ell and $\Phi\left(x^{-1} \gamma x\right) \neq 0$, then $\mathrm{ht}^{+}(x)<m$. Indeed, if $\operatorname{ht}(x) \geq m$ then there is some $\delta \in \mathrm{GL}(2, F)$ with $\mathrm{ht}^{+}(\delta x) \geq m$. Lemma 5.3 then implies that $\delta \gamma \delta^{-1} \in$ $\alpha^{\mathbb{Z}} A(F) N(F)$, contradicting $\gamma \in \alpha^{\mathbb{Z}}$. Ell.
Denote by $T_{3, m}^{\prime}(x)$ the contribution into $\chi_{m}(x) K_{f}(x, x)$ of the elements $\gamma$ of the form $\alpha^{j} \gamma, j \in \mathbb{Z}, \gamma \in \mathrm{GL}(2, F)$ with distinct eigenvalues in $F$. By $T_{4, m}^{\prime}(x)$
we denote the contribution of the elements $\alpha^{j} \gamma, j \in \mathbb{Z}, \gamma \in \mathrm{GL}(2, F), \gamma \notin F^{\times}$ but the eigenvalues of $\gamma$ are equal. We have

$$
T_{3, m}^{\prime}(x)=\frac{1}{2} \chi_{m}(x) \sum_{\gamma \in \alpha^{Z} \cdot A^{\prime}(F)} \sum_{\delta \in A(F) \backslash \operatorname{GL}(2, F)} f\left(x^{-1} \delta^{-1} \gamma \delta x\right)
$$

The factor $\frac{1}{2}$ appears since $\operatorname{diag}(b, a)$ is conjugate to $\operatorname{diag}(a, b)$. To show that $T_{3, m}^{\prime}(x)=T_{3, m}(x)$ it suffices to show that when $f\left(x^{-1} \delta^{-1} \gamma \delta x\right) \neq 0$,

$$
\chi_{m}(x)=1-\xi_{m}(\delta x)-\xi_{m}(w \delta x)
$$

namely if $\mathrm{ht}(x) \geq m$ then either $\mathrm{ht}^{+}(\delta x) \geq m$ or $\mathrm{ht}^{+}(w \delta x) \geq m$. So if $\operatorname{ht}(x) \geq m$, then there is some $\eta \in \mathrm{GL}(2, F)$ with $\mathrm{ht}^{+}(\eta x) \geq m$. By Lemma 5.3, $\eta \delta^{-1} \gamma \delta \eta^{-1} \in \alpha^{\mathbb{Z}} A(F) N(F)$, but this implies that $\eta \delta^{-1} \in A(F) N(F)$ or $\eta \delta^{-1} w \in A(F) N(F)$. Correspondingly, $\mathrm{ht}^{+}(\delta x)=\mathrm{ht}^{+}(\eta x) \geq m$ or $\mathrm{ht}^{+}(w \delta x)=\mathrm{ht}^{+}(\eta x) \geq m$, but both inequalities cannot hold simultaneously if $m>0$.
Now

$$
T_{4, m}^{\prime}(x)=\chi_{m}(x) \sum_{a \in \alpha^{\mathbb{Z}} \cdot F^{\times}} \sum_{\delta \in F^{\times}} \sum_{N(F) \backslash \operatorname{GL}(2, F)} f\left(x^{-1} \delta^{-1}\left(\begin{array}{cc}
a & a \\
0 & a
\end{array}\right) \delta x\right) .
$$

To show that this equals $T_{4, m}(x)$ we need to check that when

$$
f\left(x^{-1} \delta^{-1}\left(\begin{array}{cc}
a & a \\
0 & a
\end{array}\right) \delta x\right) \neq 0
$$

and $\operatorname{ht}(x) \geq m$, then $\mathrm{ht}^{+}(\delta x) \geq m$. Suppose then that $\mathrm{ht}^{+}(\eta x) \geq m$ for $\eta \in \mathrm{GL}(2, F)$. Then by Lemma 5.3 we have

$$
\eta \delta^{-1}\left(\begin{array}{cc}
a & a \\
0 & a
\end{array}\right) \delta \eta^{-1} \in \alpha^{\mathbb{Z}} A(F) N(F) .
$$

Hence $\eta \delta^{-1} \in A(F) N(F)$, so that ht ${ }^{+}(\delta x)=\mathrm{ht}^{+}(\eta x) \geq m$.
We conclude that $\operatorname{tr} \theta_{m} r(f)=\sum_{1 \leq i \leq 4} t_{i, m}$ with

$$
t_{i, m}=\int_{\alpha^{\mathbb{Z}} \cdot \mathrm{GL}(2, F) \backslash \mathrm{GL}(2, \mathrm{~A})} T_{i, m}(x) d x
$$

To prove the proposition it suffices to show that $t_{i, m}=S_{i}(f)+c_{i}(2 m-1)+\beta_{m}$ for all $i(1 \leq i \leq 4)$, where $c_{i}$ does not depend on $m$ and $\lim \beta_{m}=0$. It is clear that $t_{1, m} \rightarrow S_{1}(f)$ as $m \rightarrow \infty$. As $T_{2, m}(x)$ is independent of $m, t_{2, m}=S_{2}(f)$. Now

$$
\begin{gathered}
t_{3, m}=\frac{1}{2} \sum_{\gamma \in \alpha^{\mathbb{Z}} \cdot A^{\prime}(F)} \int_{\alpha^{\mathbb{Z}}(A(F) \backslash \mathrm{GL}(2, \mathbb{A})} f\left(x^{-1} \gamma x\right)\left(1-\xi_{m}(x)-\xi_{m}(w x)\right) d x \\
=\frac{1}{2} \sum_{\gamma \in \alpha^{\mathbb{Z}} \cdot A^{\prime}(F)} \int_{A(\mathbb{A}) \backslash \mathrm{GL}(2, \mathbb{A})} f\left(x^{-1} \gamma x\right) s(x) d x
\end{gathered}
$$

where

$$
\begin{gathered}
s(x)=\int_{\alpha^{\mathbb{Z}} A(F) \backslash A(\mathbb{A})}\left[1-\xi_{m}(y x)-\xi_{m}(w y x)\right] d y \\
=\operatorname{vol}\left\{y \in \alpha^{\mathbb{Z}} A(F) \backslash A(\mathbb{A}) ; \operatorname{ht}^{+}(y x)<n, \operatorname{ht}^{+}(w y x)<n\right\} .
\end{gathered}
$$

Note that for $y \in A(\mathbb{A}), \mathrm{ht}^{+}(y x)=\mathrm{ht}^{+}(y)+\mathrm{ht}^{+}(x)$ and $\mathrm{ht}^{+}(w y x)=\mathrm{ht}^{+}(w x)-$ $\mathrm{ht}^{+}(y)$. Hence

$$
s(x)=\left|\left\{y \in A(\mathbb{A}) / \alpha^{\mathbb{Z}} \cdot A(F) ; \mathrm{ht}^{+}(w x)-m<\mathrm{ht}^{+}(y)<m-\mathrm{ht}^{+}(x)\right\}\right| .
$$

This is the number of integers between $\mathrm{ht}^{+}(w x)-m$ and $m-\mathrm{ht}^{+}(x)$. So $s(x)=2 m-1-\mathrm{ht}^{+}(x)-\mathrm{ht}^{+}(w x)$.

Lemma 5.5. We have $\mathrm{ht}^{+}(x)+\mathrm{ht}^{+}(w x)=-2 r(x)$, where if $x=a\left(\begin{array}{ll}1 & y \\ 0 & 1\end{array}\right) k$, $a \in A(\mathbb{A}), k \in \mathrm{GL}(2, O)$ and $y \in \mathbb{A}$, we put $r(x)=\sum_{v} \max \left(0, \log _{q}\left|y_{v}\right|_{v}\right)$.

Proof. Note that $y$ is determined up to a change $y \mapsto b y+c, b \in O^{\times}, c \in O$, so $r(x)$ is well defined. The asserted relation does not change if $x$ is replaced by $a x k, a \in A(\mathbb{A}), k \in \operatorname{GL}(2, O)$, so we may assume $x=\left(\begin{array}{ll}1 & y \\ 0 & 1\end{array}\right) \in N(\mathbb{A})$. Then $\operatorname{ht}^{+}(x)=0$, and $\left(\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right)\left(\begin{array}{ll}1 & y \\ 0 & 1\end{array}\right)=\left(\begin{array}{cc}-\frac{1}{y} & 1 \\ 0 & y\end{array}\right)\left(\begin{array}{cc}1 & 0 \\ \frac{1}{y} & 1\end{array}\right)$ implies that ht ${ }^{+}(w x)=$ $-2 r(x)$.

Lemma 5.5 implies that

$$
t_{3, m}=S_{3}(f)+\left(m-\frac{1}{2}\right) \sum_{\gamma \in \alpha^{Z} \cdot A^{\prime}(F)} \int_{A(\mathbb{A}) \backslash \mathrm{GL}(2, \mathbb{A})} f\left(x^{-1} \gamma x\right) d x
$$

Next

$$
\begin{gathered}
t_{4, m}=\sum_{a \in \alpha^{\mathbb{Z}} \cdot F^{\times}} \int_{\alpha^{\mathbb{Z}} F^{\times} N(F) \backslash \mathrm{GL}(2, \mathbb{A})} f\left(x^{-1}\left(\begin{array}{cc}
a & a \\
0 & a
\end{array}\right) x\right)\left(1-\xi_{m}(x)\right) d x \\
=\sum_{a \in \alpha^{\mathbb{Z}} \cdot F^{\times}} \int_{\left\{x \in \alpha^{\mathbb{Z}} F^{\times} N(F) \backslash \operatorname{GL}(2, \mathbb{A}) ; \mathrm{ht}^{+}(x)<m\right\}} f\left(x^{-1}\left(\begin{array}{cc}
a & a \\
0 & a
\end{array}\right) x\right) d x .
\end{gathered}
$$

Recall that $\theta_{a, f}(t)=\int_{\alpha^{\mathbb{Z}} F^{\times} N(F) \backslash \operatorname{GL}(2, \mathbb{A})} f\left(x^{-1}\left(\begin{array}{cc}a & a \\ 0 & a\end{array}\right) x\right) t^{h^{+}(x)} d x$ is a Laurent series at $t=0$ of a rational function of $t$ with $\zeta_{F}\left(q^{-1} t\right)^{-1} \theta_{a, f}(t) \in \mathbb{C}\left[t, t^{-1}\right]$. Suppose $\theta_{a, f}(t)=\sum_{k} u_{k}(a) t^{k}$. Then $t_{4, m}=\sum_{a \in \alpha^{\mathbb{Z}} \cdot F \times} \sum_{k<m} u_{k}(a)$. Since $\zeta_{F}\left(q^{-1} t\right)$ has a simple pole at $t=1$, we have that $\theta_{a, f}(t)=\frac{\rho(a)}{1-t}+\bar{\theta}_{a, f}(t)$, with $\bar{\theta}_{a, f}(t)$ without poles on $0<|t| \leq 1$. Then

$$
\begin{gathered}
\tilde{\theta}_{a, f}(t)=\frac{1}{2}\left(\theta_{a, f}(t)+\theta_{a, f}\left(t^{-1}\right)\right)=\frac{1}{2}\left(\bar{\theta}_{a, f}(t)+\bar{\theta}_{a, f}\left(t^{-1}\right)\right)+\frac{1}{2} \rho(a), \\
\tilde{\theta}_{a, f}(1)=\bar{\theta}_{a, f}(1)+\frac{1}{2} \rho(a)=\frac{1}{2} \rho(a)+\sum_{k}\left(u_{k}(a)-\rho(a)\right)
\end{gathered}
$$

$$
=\lim _{m \rightarrow \infty}\left[\sum_{k<m} u_{k}(a)-\left(m-\frac{1}{2}\right) \rho(a)\right]
$$

Then

$$
t_{4, m}=\sum_{a \in \alpha^{Z} \cdot F^{\times}} \tilde{\theta}_{a, f}(1)+\left(m-\frac{1}{2}\right) \rho(a)+\beta_{m}, \quad \beta_{m} \rightarrow 0 \text { as } m \rightarrow \infty
$$

and $S_{4}(f)=\sum_{a \in \alpha^{Z} \cdot F^{\times}} \tilde{\theta}_{a, f}(1)$. Proposition 5.2 follows.
Note that $\beta_{m}$ is 0 for sufficiently large $m$, as will be seen below.

### 5.2 The Eisenstein contribution

Next we turn to computing $\operatorname{tr}\left(\theta_{m} r_{E}^{U}(f)\right)$ for large $m$. Put $W_{c}^{U}=C_{c}^{\infty}\left(Y_{\alpha}\right)^{U}$, $W_{M}^{U}=(1+M) W_{c}^{U}$.
Proposition 5.6. The operator $E^{*}$ maps $A_{E}^{U}$ isomorphically onto $W_{M}^{U}$.
Proof. As $A_{E}^{U}=E\left(W_{c}^{U}\right)$ and $E^{*} E=1+M$, it suffices to show that ker $E^{*} E=$ ker $E$. For $\varphi \in \operatorname{ker} E^{*} E$ we have $(E \varphi, E \varphi)=\left(E^{*} E \varphi, \varphi\right)=0$, hence $E \varphi=$ 0 .

Definition 1. Denote by $W_{m}^{U}$ the space of $f$ in $W_{c}^{U}$ with $f(x)=0$ if $\mathrm{ht}^{+}(x)<$ $m$. Denote by $\xi_{m}$ also the operator $W_{m}^{U} \rightarrow W_{m}^{U}$ of multiplication by the characteristic function of the set $\left\{x \in Y_{\alpha} ;\right.$ ht $\left.^{+}(x) \geq m\right\}$. [If $m>0$ then $\xi_{m}$ is a left inverse to the operator $1+M: W_{m}^{U} \rightarrow W_{M}^{U}$. Indeed, if $f$ is in $W_{m}^{U}$, then $(M f)(x)=0$ already when $\mathrm{ht}^{+}(x)>-m$ since $\mathrm{ht}^{+}(w n x)+\mathrm{ht}^{+}(n x)<0$ implies $\mathrm{ht}^{+}(w n x)<m$ and so $f(w n x)=0$.] Hence $\pi^{m}=(1+M) \xi_{m}: W_{M}^{U} \rightarrow W_{M}^{U}$ satisfies $\pi^{m} \pi^{m}=\pi^{m}$, for $m>0$. Put $\pi_{m}=1-\pi^{m}$.
Proposition 5.7. For sufficiently large $m, E^{*}$ intertwines $\theta_{m}$ with $\pi_{m}$, thus $\pi_{m} E^{*}=E^{*} \theta_{m}$, namely the diagram

$$
\begin{aligned}
& A_{E}^{U} \xrightarrow{E^{*}} W_{M}^{U}
\end{aligned}
$$

is commutative.
Proof. Suppose $f \in A_{E}^{U}$ and $\left(1-\theta_{m}\right) f=0$. Then $f(x)=0$ for $x$ with ht $(x) \geq$ $m$. As $\xi_{m}(x) \neq 0$ only on $x$ with $\mathrm{ht}^{+}(x) \geq m$, we have $0=(1+M) \xi_{m} E^{*} f=$ $\left(1-\pi_{m}\right) E^{*} f$, the last equality as $1-\pi_{m}=\pi^{m}=(1+M) \xi_{m}$. For such $f$ we have $E^{*} \theta_{m} f=E^{*} f$ and $\pi_{m} E^{*} f=E^{*} f$.
If $f \in A_{E}^{U}$ and $\theta_{m} f=0$, then by Proposition 4.34 there is $\varphi \in W_{m}^{U}$ with $f=E \varphi$. Then

$$
\pi_{m} E^{*} f=\pi_{m} E^{*} E \varphi=\pi_{m}(1+M) \varphi=\pi_{m}(1+M) \xi_{m} \varphi=\pi_{m} \pi^{m} \varphi=0
$$

hence $E^{*} \theta_{m} f=\pi_{m} E^{*} f$ for such $f$.
Any $f \in A_{E}^{U}$ can be written as $f=f_{1}+f_{2}, f_{1}=\left(1-\theta_{m}\right) f, f_{2}=\theta_{m} f$, thus $\theta_{m} f_{1}=0$ and $\left(1-\theta_{m}\right) f_{2}=0$.

Definition 2. Recall that $Y_{\alpha}=\alpha^{\mathbb{Z}} A(F) N(\mathbb{A}) \backslash \mathrm{GL}(2, \mathbb{A})$. Denote by $\sigma_{c}, \sigma_{+}$, $\sigma_{M}$ the representations of $\mathrm{GL}(2, \mathbb{A})$ in the spaces

$$
W_{c}=C_{c}^{\infty}\left(Y_{\alpha}\right), \quad W_{+}=C_{+}^{\infty}\left(Y_{\alpha}\right), \quad W_{M}=(1+M) C_{c}^{\infty}\left(Y_{\alpha}\right)
$$

Consider $\sigma_{c}(f), \sigma_{+}(f), \sigma_{M}(f)$ as operators in the spaces $W_{c}^{U}, W_{+}^{U}, W_{M}^{U}$.
Corollary 5.8. We have $\operatorname{tr}\left(\theta_{m} \cdot r_{E}^{U}(f)\right)=\operatorname{tr}\left(\pi_{m} \cdot \sigma_{M}(f)\right)$.
Proof. The operator $E^{*}$ yields an isomorphism of $A_{E}^{U}=E\left(W_{c}^{U}\right)$ with $W_{M}^{U}$ intertwining $\theta_{m}$ with $\pi_{m}$.

In the proof of Proposition 4.35 we introduced a structure of $\mathbb{C}\left[z, z^{-1}\right]$-module on $W_{c}^{U}$ and $W_{+}^{U}$, as well as isomorphisms $W_{c}^{U} \simeq W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]$ and $W_{+}^{U} \simeq$ $W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}((z))$, where $W_{0}^{U}=\left\{f \in W_{c}^{U} ; f(x)=0\right.$ if ht $\left.^{+}(x) \neq 0\right\}$. Under these isomorphisms, the operator $M: W_{c}^{U} \rightarrow W_{+}^{U}$ corresponds to the operator

$$
M: W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right] \rightarrow W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}((z))
$$

which satisfies the conditions of Proposition 4.31, hence has the form

$$
(M u)(z)=P(z) u\left(z^{-1}\right) \quad \text { for } \quad u \in W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]
$$

which is viewed as a function of $z$ with values in $W_{0}^{U}$. Here $P(z)$ is a rational function in $z$ with values in Aut $W_{0}^{U}$, and $P\left(z^{-1}\right)=P(z)^{-1}$.
Now $\sigma_{c}(f)$ is an endomorphism of $W_{c}^{U}$ as a $\mathbb{C}\left[z, z^{-1}\right]$-module. The corresponding endomorphism of the module $W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]$ is determined by a function $B(z)$ in $\operatorname{End}\left(W_{0}^{U}\right) \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]$. The endomorphism of $W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}((z))$ corresponding to the operator $\sigma_{+}(f)$ is determined by the same function $B(z)$. The relation

$$
M \sigma_{c}(f)=\sigma_{+}(f) M \quad \text { becomes } \quad P(z) B\left(z^{-1}\right) u\left(z^{-1}\right)=B(z) P(z) u\left(z^{-1}\right)
$$

for any $u \in W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}\left[z, z^{-1}\right]$, thus $B\left(z^{-1}\right)=P(z)^{-1} B(z) P(z)$.
Definition 3. Under the isomorphism $W_{+}^{U} \simeq W_{0}^{U} \otimes_{\mathbb{C}} \mathbb{C}((z))$, the subspace $W_{M}^{U}=(1+M) W_{c}^{U}$ is mapped onto the subspace $L$ consisting of all rational functions of the form $u(z)+P(z) u\left(z^{-1}\right)$, with $u \in W_{0}^{U} \otimes \mathbb{C} \mathbb{C}\left[z, z^{-1}\right]$. Put

$$
L_{m}=L \cap\left(W_{0}^{U} \otimes_{\mathbb{C}} z^{-m+1} \mathbb{C}[[z]]\right) .
$$

Denote by $L^{m}$ the set of rational functions of the form

$$
u(z)+P(z) u\left(z^{-1}\right) \quad \text { with } \quad u \in W_{0}^{U} \otimes_{\mathbb{C}} z^{-m} \mathbb{C}\left[z^{-1}\right]
$$

For sufficiently large $m$ we have $L=L_{m} \oplus L^{m}$. Under the isomorphism $W_{M}^{U} \xrightarrow{\sim}$ $L$, the operator $\pi_{m}: W_{M}^{U} \rightarrow W_{M}^{U}$ corresponds to the idempotent operator $L \rightarrow L$ with kernel $L^{m}$ and image $L_{m}$. This projection will also be denoted by $\pi_{m}$. Thus $\operatorname{tr}\left(\pi_{m} \sigma_{M}(f)\right)=\operatorname{tr}\left(\pi_{m} B\right)$, where $B: L \rightarrow L$ is the operator of
multiplication by $B(z)$. On the left, $\pi_{m}$ is an operator on $W_{M}^{U}$, on the right, on $L$.
Fix $Q_{1}, Q_{2} \in M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right), k \geq 1$, such that $\operatorname{det} Q_{i} \neq 0$. Suppose the function $Q_{2}(z)^{-1} Q_{1}(z)$ is regular at $z=\infty$, thus $Q_{1}(z) \in Q_{2}(z) M\left(k, \mathbb{C}\left[\left[z^{-1}\right]\right]\right)$, and the function $Q_{1}(z)^{-1} Q_{2}(z)$ is regular at $z=0$, thus $Q_{2}(z) \in Q_{1}(z) M(k, \mathbb{C}[[z]])$. Put $R=\mathbb{C}\left[z, z^{-1}\right]^{k}$. For $m \geq 1$, put

$$
R_{m}=R \cap z^{1-m} Q_{1}(z) \mathbb{C}[[z]]^{k} \cap z^{m-1} Q_{2}(z) \mathbb{C}\left[\left[z^{-1}\right]\right]^{k}
$$

Also put

$$
R_{-}^{m}=z^{-m} Q_{1}(z) \mathbb{C}\left[z^{-1}\right]^{k} \quad R_{+}^{m}=z^{m} Q_{2}(z) \mathbb{C}[z]^{k}
$$

Then $\operatorname{dim} R_{m}$ is finite.
Proposition 5.9. We have $R=R_{-}^{m} \oplus R_{m} \oplus R_{+}^{m}$,

$$
R_{m} \oplus R_{+}^{m}=R \cap z^{1-m} Q_{1}(z) \mathbb{C}[[z]]^{k}
$$

and

$$
R_{m} \oplus R_{-}^{m}=R \cap z^{m-1} Q_{2}(z) \mathbb{C}\left[\left[z^{-1}\right]\right]^{k}
$$

Proof. The natural map $\varphi: R_{-}^{m} \rightarrow X_{-}=\mathbb{C}((z))^{k} / z^{1-m} Q_{1}(z) \mathbb{C}[[z]]^{k}$ is an isomorphism (note that $\mathbb{C}((z)) / z^{1-m} \mathbb{C}[[z]] \simeq z^{-m} \mathbb{C}\left[z^{-1}\right]$ and $Q_{1}(z)$ is invertible in $\operatorname{GL}(k, \mathbb{C}((z)))$. The natural map

$$
\psi: R_{+}^{m} \rightarrow X_{+}=\mathbb{C}\left(\left(z^{-1}\right)\right)^{k} / z^{m-1} Q_{2}(z) \mathbb{C}\left[\left[z^{-1}\right]\right]^{k}
$$

is then too. The natural map $f: R / R_{m} \rightarrow X_{-} \oplus X_{+}$is injective (by definition of $R_{m}$ as the intersection of $R$ and the denominators of $X_{-}, X_{+}$) and the composition of the natural map $R_{+}^{m} \oplus R_{-}^{m} \rightarrow R / R_{m}$ with $f$ is $\varphi \oplus \psi$.

Definition 4. (1) Denote by $\mathrm{pr}_{m}: R \rightarrow R$ the projection with kernel $R_{+}^{m} \oplus R_{-}^{m}$ and image $R_{m}$. (2) If $A(z)$ is a matrix in $M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right)$, denote by $A[z]$ also the corresponding automorphism of $R=\mathbb{C}\left[z, z^{-1}\right]^{k}$. Denote by $A_{0}$ the constant term of $A(z)$.
Proposition 5.10. The trace $\operatorname{tr}\left(\operatorname{pr}_{m} \cdot A[z]\right)$ is equal to
$(2 m-1) \operatorname{tr} A_{0}-\operatorname{res}_{z=0} \operatorname{tr} A(z) Q_{1}^{\prime}(z) Q_{1}(z)^{-1} d z-\operatorname{res}_{z=\infty} \operatorname{tr} A(z) Q_{2}^{\prime}(z) Q_{2}(z)^{-1} d z$.
Proof. Define a projection $\operatorname{pr}_{+}^{m}: R \rightarrow R$ with image $R_{+}^{m}$ and kernel $R_{-}^{m}+R_{m}$, and a projection $\operatorname{pr}_{-}^{m}: R \rightarrow R$ with image $R_{-}^{m}$ and kernel $R_{+}^{m}+R_{m}$. Analogously to the decomposition $R=R_{-}^{m} \oplus R_{m} \oplus R_{+}^{m}$, consider the decomposition

$$
R=z^{-m} \mathbb{C}\left[z^{-1}\right]^{k} \oplus\left(z^{1-m} \mathbb{C}[z]^{k} \cap z^{m-1} \mathbb{C}\left[z^{-1}\right]^{k}\right) \oplus z^{m} \mathbb{C}[z]^{k}
$$

namely the case where $Q_{1}=1=Q_{2}$. Denote the associated projections by $p_{-}^{m}$, $p_{m}, p_{+}^{m}$. Since the space $z^{-m} \mathbb{C}\left[z^{-1}\right]^{k} / R_{-}^{m} \cap z^{-m} \mathbb{C}\left[z^{-1}\right]^{k}$ is finite dimensional, the operator $\mathrm{pr}_{+}^{m}-p_{+}^{m}$ has finite rank, and the operator $\mathrm{pr}_{-}^{m}-p_{-}^{m}$ has finite rank since $z^{m} \mathbb{C}[z]^{k} / R_{+}^{m} \cap z^{m} \mathbb{C}[z]^{k}$ is finite dimensional.

Lemma 5.11. We have

$$
\operatorname{tr}\left(\operatorname{pr}_{-}^{m} \cdot A[z]-p_{-}^{m} \cdot A[z]\right)=\operatorname{res}_{z=0} \operatorname{tr} A(z) Q_{1}^{\prime}(z) Q_{1}(z)^{-1} d z
$$

as well as

$$
\operatorname{tr}\left(\operatorname{pr}_{+}^{m} \cdot A[z]-p_{+}^{m} \cdot A[z]\right)=\operatorname{res}_{z=\infty} \operatorname{tr} A(z) Q_{2}^{\prime}(z) Q_{2}(z)^{-1} d z
$$

Proof. Denote by $\operatorname{Pr}_{-}^{m}: \mathbb{C}((z))^{k} \rightarrow \mathbb{C}((z))^{k}$ the projection with image $z^{-m} Q_{1}(z) \mathbb{C}\left[z^{-1}\right]^{k}$ and kernel $z^{1-m} Q_{1}(z) \mathbb{C}[[z]]^{k}$. Denote by $P_{-}^{m}: \mathbb{C}((z))^{k} \rightarrow$ $\mathbb{C}((z))^{k}$ the projection with image $z^{-m} \mathbb{C}\left[z^{-1}\right]^{k}$ and kernel $z^{1-m} \mathbb{C}[[z]]^{k}$ (thus the case of $\left.Q_{1}=1\right)$. Denote by $A((z))$ the endomorphism of $\mathbb{C}((z))^{k}$ which is defined by multiplication by $A(z)$. Then $\operatorname{Pr}_{-}^{m}=Q_{1}((z)) \cdot P_{-}^{m} \cdot Q_{1}((z))^{-1}$. Now

$$
\operatorname{Im}\left(\operatorname{Pr}_{-}^{m} \cdot A((z))-P_{-}^{m} \cdot A((z))\right) \subset \mathbb{C}\left[z, z^{-1}\right]^{k}
$$

and the restriction of the operator

$$
\operatorname{Pr}_{-}^{m} \cdot A((z))-P_{-}^{m} \cdot A((z)) \quad \text { to } \quad \mathbb{C}\left[z, z^{-1}\right]^{k} \quad\left(\subset \mathbb{C}((z))^{k}\right)
$$

is equal to $\operatorname{pr}_{-}^{m} \cdot A[z]-p_{-}^{m} \cdot A[z]$. Hence

$$
\begin{gathered}
\operatorname{tr}\left(\operatorname{pr}_{-}^{m} \cdot A[z]-p_{-}^{m} \cdot A[z]\right)=\operatorname{tr}\left(\operatorname{Pr}_{-}^{m} \cdot A((z))-P_{-}^{m} \cdot A((z))\right) \\
=\operatorname{tr}\left(Q_{1}((z)) \cdot P_{-}^{m} \cdot Q_{1}((z))^{-1} \cdot A((z))-P_{-}^{m} \cdot A((z))\right) \\
=\operatorname{tr}\left(Q_{1}((z)) \cdot P_{-}^{m} \cdot C((z))-P_{-}^{m} \cdot Q_{1}((z)) \cdot C((z))\right), \quad C(z)=Q_{1}(z)^{-1} A(z)
\end{gathered}
$$

As $\operatorname{tr} A(z) Q_{1}^{\prime}(z) Q_{1}(z)^{-1}=\operatorname{tr} C(z) Q_{1}^{\prime}(z)$, to prove the first claim of the lemma it suffices to show that

$$
\operatorname{tr}\left(Q_{1}((z)) \cdot P_{-}^{m} \cdot C((z))-P_{-}^{m} \cdot Q_{1}((z)) C((z))\right)=\operatorname{res}_{z=0} \operatorname{tr} C(z) Q_{1}^{\prime}(z) d z
$$

for any $Q_{1}(z) \in M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right), C(z) \in M(k, \mathbb{C}((z)))$. By linearity, it suffices to show this when the matrices $Q_{1}(z)$ and $C(z)$ have a single nonzero entry. Thus we may assume $k=1$, and that $Q_{1}(z)=z^{b}$. Thus we need to verify that for any formal power series $c(z)=\sum_{d} c_{d} z^{d}$ in $\mathbb{C}((z))$, we have

$$
\operatorname{tr}\left[\left(\left(\left(z^{b}\right)\right) \cdot P_{-}^{m}-P_{-}^{m} \cdot\left(\left(z^{b}\right)\right)\right) c((z))\right]=b c_{-b}
$$

where the operations here are in $\mathbb{C}((z))$. The left side is equal to

$$
\begin{gathered}
\operatorname{tr}\left[\left(\left(\left(z^{b}\right)\right) \cdot P_{-}^{m} \cdot\left(\left(z^{-b}\right)\right)-P_{-}^{m}\right) \cdot\left(\left(z^{b}\right)\right) c((z))\right]=\operatorname{tr}\left[\left(P_{-}^{m-b}-P_{-}^{m}\right) \cdot\left(\left(z^{b}\right)\right) c((z))\right] \\
\\
=\operatorname{tr}\left(\begin{array}{cccc}
c_{-b} & c_{-b+1} & \cdots & c_{-1} \\
\vdots & \vdots & \cdots & \vdots \\
c_{-b-2 b} & c_{2}-2 b & \cdots & c_{-b}
\end{array}\right)=b c_{-b} .
\end{gathered}
$$

The second claim of the lemma is similarly proven.

As $\operatorname{pr}_{m}-p_{m}=\left(1-\operatorname{pr}_{-}^{m}-\operatorname{pr}_{+}^{m}\right)-\left(1-p_{-}^{m}-p_{+}^{m}\right)=\left(p_{-}^{m}-\operatorname{pr}_{-}^{m}\right)+\left(p_{+}^{m}-\operatorname{pr}_{+}^{m}\right)$, Lemma 5.11 implies that $\operatorname{tr}\left(\operatorname{pr}_{m} \cdot A[z]-p_{m} \cdot A[z]\right)$

$$
=-\operatorname{res}_{z=0} \operatorname{tr}\left[A(z) Q_{1}^{\prime}(z) Q_{1}(z)^{-1} d z\right]-\operatorname{res}_{z=\infty} \operatorname{tr}\left[A(z) Q_{2}^{\prime}(z) A_{2}(z)^{-1} d z\right]
$$

Since $\operatorname{tr}\left(p_{m} \cdot A[z]\right)=(2 m-1) \operatorname{tr} A_{0}$, the proposition follows.
Proposition 5.12. Let $\iota: \mathbb{C}\left[z, z^{-1}\right]^{k} \rightarrow \mathbb{C}\left[z, z^{-1}\right]^{k}$ be the involution $(\iota u)(z)=$ $u\left(z^{-1}\right)$. For sufficiently large $m$ we have $2 \operatorname{tr}\left(\iota \cdot \operatorname{pr}_{m} \cdot A[z]\right)=\operatorname{tr} A(1)+\operatorname{tr} A(-1)$. Proof. Write $A(z)=\sum_{k} A_{k} z^{k}, A_{k} \in M(k, \mathbb{C})$. Then $\operatorname{tr}\left(\iota \cdot p_{m} \cdot A[z]\right)=$ $\sum_{|i|<m} \operatorname{tr} A_{2 i}$. If $m$ is big enough the right side here is equal to $\frac{1}{2}(\operatorname{tr} A(1)+$ $\operatorname{tr} A(-1))$. It remains to show that $\operatorname{tr}\left(\iota \cdot \operatorname{pr}_{m} \cdot A[z]\right)=\operatorname{tr}\left(\iota \cdot p_{m} \cdot A[z]\right)$ for large enough $m$. As $\mathrm{pr}_{m}-p_{m}=p_{+}^{m}-\operatorname{pr}_{+}^{m}+\left(p_{-}^{m}-\operatorname{pr}_{-}^{m}\right)$, it suffices to show that for large enough $m$

$$
\operatorname{tr}\left(\iota \cdot\left(p_{+}^{m}-\operatorname{pr}_{+}^{m}\right) \cdot A[z]\right)=0=\operatorname{tr}\left(\iota \cdot\left(p_{-}^{m}-\operatorname{pr}_{-}^{m}\right) \cdot A[z]\right)
$$

Note that $\operatorname{pr}_{+}^{m}=\left[z^{m}\right] \operatorname{pr}_{+}^{0}\left[z^{-m}\right]$ and $p_{+}^{m}=\left[z^{m}\right] p_{+}^{0}\left[z^{-m}\right]$, where as usual $\left[z^{m}\right]$ here means the operator of multiplication by $z^{m}$. The operators $\mathrm{pr}_{+}^{m}$ and $p_{+}^{m}$ were defined only for $m>0$, but the definition extends to $m=0$ so that the two relations above hold. Now

$$
\begin{gathered}
\operatorname{tr}\left(\iota \cdot\left(p_{+}^{m}-\operatorname{pr}_{+}^{m}\right) \cdot A[z]\right)=\operatorname{tr}\left(\iota \cdot\left[z^{m}\right]\left(p_{+}^{0}-\operatorname{pr}_{+}^{0}\right)\left[z^{-m}\right] \cdot A[z]\right) \\
=\operatorname{tr}\left(\left[z^{-m}\right] \iota \cdot\left(p_{+}^{0}-\operatorname{pr}_{+}^{0}\right)\left[z^{-m}\right] \cdot A[z]\right)=\operatorname{tr}\left(\iota \cdot\left(p_{+}^{0}-\operatorname{pr}_{+}^{0}\right)\left[z^{-m}\right] \cdot A[z]\left[z^{-m}\right]\right) \\
=\operatorname{tr}\left(\iota \cdot\left(p_{+}^{0}-\operatorname{pr}_{+}^{0}\right)\left[z^{-2 m}\right] \cdot A[z]\right)
\end{gathered}
$$

Recall that $\operatorname{dim} V$ is finite, where $V=\operatorname{im}\left[\iota\left(p_{+}^{0}-\operatorname{pr}_{+}^{0}\right)\right]$. If $m$ is big enough then

$$
\left[z^{-2 m}\right] \cdot A[z] V \subset z^{-1} \mathbb{C}\left[z^{-1}\right]^{k} \cap z^{-1} Q_{2}(z) \mathbb{C}\left[\left[z^{-1}\right]\right]^{k} \subset \operatorname{ker} p_{+}^{0} \cap \operatorname{ker} \operatorname{pr}_{+}^{0} .
$$

Hence $\operatorname{tr}\left(\iota \cdot\left(p_{+}^{0}-\operatorname{pr}_{+}^{0}\right)\left[z^{-2 m}\right] \cdot A[z]\right)$ is zero. Hence $\operatorname{tr}\left(\iota\left(p_{+}^{m}-\operatorname{pr}_{+}^{m}\right) A[z]\right)$ is zero. The proof of $\operatorname{tr}\left(\iota\left(p_{-}^{m}-\operatorname{pr}_{-}^{m}\right) A[z]\right)=0$ for large $m$ is analogous.

Definition 5. Fix $P \in \operatorname{GL}(k, \mathbb{C}(z))$ such that $P(z)$ is regular at $z=0$ and $P(z)^{-1}$ is regular at $z=\infty$. Put $S=\mathbb{C}\left[z, z^{-1}\right]^{k}+P \cdot \mathbb{C}\left[z, z^{-1}\right]^{k}$,
$S_{m}=S \cap z^{1-m} \mathbb{C}[[z]]^{k} \cap z^{m-1} P \cdot \mathbb{C}[[z]]^{k}, \quad S^{m}=z^{-m} \mathbb{C}\left[z^{-1}\right]^{k}+z^{m} P \cdot \mathbb{C}[z]^{k}$.
Fix $B$ in $M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right)$ such that $P^{-1} B P$ lies in $M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right)$. Then $B S \subset$ $S$. We denote by $[B]$ or $B[z]$ the operator $S \rightarrow S$ of multiplication by $B(z)$.
Proposition 5.13. We have $S=S_{m} \oplus S^{m}$. Denote by $\mathrm{ps}_{m}: S \rightarrow S$ the projection with image $S_{m}$ and kernel $S^{m}$. Then

$$
\begin{aligned}
& \operatorname{tr}\left(\mathrm{ps}_{m} \cdot[B]\right)=(2 m-1) \operatorname{tr} B_{0} \\
& -\operatorname{res}_{z=\infty} \operatorname{tr}\left[B(z) P^{\prime}(z) P(z)^{-1}\right] d z+\operatorname{tr}\left([B] ; S / \mathbb{C}\left[z, z^{-1}\right]^{k}\right)
\end{aligned}
$$

Here $B_{0}$ is the constant term of $B=B(z)$, and $\operatorname{tr}\left([B] ; S / \mathbb{C}\left[z, z^{-1}\right]^{k}\right)$ denotes the trace of the endomorphism of $S / \mathbb{C}\left[z, z^{-1}\right]^{k}$ induced by multiplication by $B(z)$.

Proof. The space $S$ is a $k$-dimensional free $\mathbb{C}\left[z, z^{-1}\right]$-submodule of $\mathbb{C}(z)^{k}$. Hence there exists a matrix $D$ in $\operatorname{GL}(k, \mathbb{C}(z))$ such that $S=D \cdot \mathbb{C}\left[z, z^{-1}\right]^{k}$. Since $S$ contains $\mathbb{C}\left[z, z^{-1}\right]^{k}, D^{-1}$ lies in $M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right)$. Since $S$ contains $P \cdot \mathbb{C}\left[z, z^{-1}\right]^{k}$ we deduce that $D^{-1} P \in M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right)$. Put $Q_{1}=D^{-1}, Q_{2}=D^{-1} P$. The function $Q_{1}(z)^{-1} Q_{2}(z)=P(z)$ is regular at $z=0$. The function $Q_{2}(z)^{-1} Q_{1}(z)$ is regular at $z=\infty$. Under the isomorphism $S \stackrel{\sim}{\rightarrow} \mathbb{C}\left[z, z^{-1}\right]^{k}, u \mapsto D^{-1} u$, the subspaces $S_{m}$ and $S^{m}$ correspond to the subspaces $R_{m}$ and $R^{m}$ of Proposition 5.9. The multiplication $[B]: S \rightarrow S$ corresponds to $[A]: \mathbb{C}\left[z, z^{-1}\right]^{k} \rightarrow \mathbb{C}\left[z, z^{-1}\right]^{k}$, $A=D^{-1} B D$. Then Proposition 5.10 implies the first part of the proposition, as well as the equality

$$
\begin{aligned}
\operatorname{tr}\left(\mathrm{ps}_{m} \cdot B[z]\right)= & (2 m-1) \operatorname{tr} A_{0}-\operatorname{res}_{z=0} \operatorname{tr} A(z) Q_{1}^{\prime}(z) Q_{1}(z)^{-1} d z \\
& -\operatorname{res}_{z=\infty} \operatorname{tr} A(z) Q_{2}^{\prime}(z) Q_{2}(z)^{-1} d z
\end{aligned}
$$

Here $A_{0}$ is the constant term of $A(z)$. We have

$$
\begin{gathered}
\operatorname{tr}\left(A Q_{1}^{\prime} Q_{1}^{-1}\right)=-\operatorname{tr}\left(D^{-1} B D^{\prime}\right)=-\operatorname{tr}\left(B D^{\prime} D^{-1}\right) \\
\operatorname{tr}\left(A Q_{2}^{\prime} Q_{2}^{-1}\right)=-\operatorname{tr}\left(D^{-1} B P^{\prime} P^{-1} D-D^{-1} B D^{\prime}\right)=\operatorname{tr}\left(B P^{\prime} P^{-1}\right)-\operatorname{tr}\left(B D^{\prime} D^{-1}\right) .
\end{gathered}
$$

As $A=D^{-1} B D$, we have $\operatorname{tr} A=\operatorname{tr} B$, and $\operatorname{tr} A_{0}=\operatorname{tr} B_{0}$. Hence

$$
\left.\begin{array}{rl}
\operatorname{tr}\left(\mathrm{ps}_{m} \cdot B[z]\right)=(2 m-1) \operatorname{tr} B_{0}-\operatorname{res}_{z=\infty} & \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z \\
+\operatorname{res}_{z=0} \operatorname{tr} B(z) D^{\prime}(z) D(z)^{-1} d z+\operatorname{res}_{z=\infty} \operatorname{tr} B(z) D^{\prime}(z) D(z)^{-1} d z \\
& +(2 m-1)
\end{array}\right) \operatorname{tr} B_{0}-\operatorname{res}_{z=\infty} \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z .
$$

Lemma 5.14. Suppose $T \in \operatorname{GL}(k, \mathbb{C}((z))), C \in M(k, \mathbb{C}[[z]])$ and $T^{-1} C T \in$ $M(k, \mathbb{C}[[z]])$. Then $\operatorname{res}_{z=0} \operatorname{tr} C(z) T^{\prime}(z) T(z)^{-1}=a-b$, where a denotes the trace of the operator multiplication by $C$ in the space $\left(\mathbb{C}[[z]]^{k}+T \mathbb{C}[[z]]^{k}\right) / T \mathbb{C}[[z]]^{k}$, while $b$ denotes the trace of multiplication by $C$ in the space

$$
\left(\mathbb{C}[[z]]^{k}+T \mathbb{C}[[z]]^{k}\right) / \mathbb{C}[[z]]^{k} .
$$

Proof. Both sides of the asserted equality do not change if $(T, C)$ is replaced by ( $U T V, U C U^{-1}$ ) where $U, V \in \mathrm{GL}(k, \mathbb{C}[[z]])$. We may then assume that $T$ is a diagonal matrix, hence that $k=1$. When $k=1$ both sides of the asserted relation are simply $m C(0)$, where $m$ is the multiplicity of zero of $T(z)$ at $z=0$.

It follows from the lemma that $-\operatorname{res}_{z=\zeta} \operatorname{tr}\left(B(z) D^{\prime}(z) D(z)^{-1}\right) d z$ is just the trace of the operator of multiplication by $B(z)$ on the $\zeta$ component of the module $S / \mathbb{C}\left[z, z^{-1}\right]^{k}$. This, and the equality just before the lemma, implies the proposition.

Suppose we have $P\left(z^{-1}\right)=P(z)^{-1}$. Replace the assumption $P(z)^{-1} B(z) P(z)$ $\in M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right)$ in Proposition 5.13 with the stronger assumption

$$
P(z)^{-1} B(z) P(z)=B\left(z^{-1}\right)
$$

Recall that $L$ is the space of all rational functions of the form $u(z)+P(z) u\left(z^{-1}\right)$ with $u \in \mathbb{C}\left[z, z^{-1}\right]^{m}$. In view of the stronger assumption, $L$ is invariant under multiplication by $B$.

Definition 6. Denote by $B_{L}$ the operator of multiplication by $B$ on $L$. Put $L_{m}=L \cap z^{1-m} \mathbb{C}[[z]]^{k}$. Denote by $L^{m}$ the set of rational functions of the form $u(z)+P(z) u\left(z^{-1}\right)$ with $u(z) \in z^{-m} \mathbb{C}\left[z^{-1}\right]^{k}$.

Proposition 5.15. The space $L_{m}$ is finite dimensional, and $L=L_{m} \oplus L^{m}$. Denote by $\pi_{m}: L \rightarrow L$ the projection with image $L_{m}$ and kernel $L^{m}$. Suppose the function $P(z)$ is regular at $z= \pm 1$. Then for large enough $m$ we have that $\operatorname{tr}\left(\pi_{m} B_{L}\right)$ equals

$$
\begin{aligned}
(m & \left.-\frac{1}{2}\right) \operatorname{tr} B_{0}-\frac{1}{2} \operatorname{res}_{z=\infty} \operatorname{tr}\left(B(z) P^{\prime}(z) P(z)^{-1}\right) d z \\
& +\frac{c}{2}+\frac{1}{4}[\operatorname{tr}(B(1) P(1))+\operatorname{tr}(B(-1) P(-1))]
\end{aligned}
$$

Here $B_{0}$ is the constant term of $B(z)$, and $c$ is the trace of the operator of multiplication by $B(z)$ in the space $\left(\mathbb{C}\left[z, z^{-1}\right]^{k}+P(z) \mathbb{C}\left[z, z^{-1}\right]^{k}\right) / \mathbb{C}\left[z, z^{-1}\right]^{k}$.

Proof. Let $S, S_{m}, S^{m}, \mathrm{ps}_{m}, B$ be as in Proposition 5.13. From $P\left(z^{-1}\right)=$ $P(z)^{-1}$ it follows that if $u \in S$ then $\tilde{u}$, given by $\tilde{u}(z)=P(z) u\left(z^{-1}\right)$, is also in $S$. Define $\tau: S \rightarrow S$ by $\tau(u)=\tilde{u}$. Then $\tau^{2}=1, L=\{u \in S ; \tau(u)=u\}$, $L_{m}=S_{m} \cap L, L^{m}=S^{m} \cap L$, and

$$
\operatorname{tr}\left(\pi_{m} B_{L}\right)=\frac{1}{2} \operatorname{tr}\left(\mathrm{ps}_{m} \cdot B[z]\right)+\frac{1}{2} \operatorname{tr}\left(\tau \cdot \mathrm{ps}_{m} \cdot B[z]\right)
$$

The finite dimensionality of $S_{m}$ and Proposition 5.13 then imply that the space $L_{m}$ is finite dimensional, and $L=L_{m} \oplus L^{m}$. To deduce the last claim of the proposition from Proposition 5.13, it remains to show that

$$
\operatorname{tr}\left(\tau \cdot \mathrm{ps}_{m} \cdot[B]\right)=\frac{1}{2}(\operatorname{tr}(B(1) P(1))+\operatorname{tr}(B(-1) P(-1)))
$$

for large enough $m$.
Let $D, Q_{1}, Q_{2}$ be as in Proposition 5.13. Then under the isomorphism $S \xrightarrow{\sim} \mathbb{C}\left[z, z^{-1}\right]^{k}, u \mapsto D^{-1} u$, the operator $\mathrm{ps}_{m}: S \rightarrow S$ translates into the operator $\mathrm{pr}_{m}$ (of Proposition 5.9), and multiplication by $B: S \rightarrow S$ translates into multiplication by

$$
A=D^{-1} B D, \quad \mathbb{C}\left[z, z^{-1}\right]^{k} \rightarrow \mathbb{C}\left[z, z^{-1}\right]^{k}
$$

The map $\tau: S \rightarrow S$ translates into

$$
[C] \iota: \mathbb{C}\left[z, z^{-1}\right]^{k} \rightarrow \mathbb{C}\left[z, z^{-1}\right]^{k}, \quad(\iota u)(z)=u\left(z^{-1}\right), \quad C(z)=D(z)^{-1} P(z) D\left(z^{-1}\right) .
$$

Hence

$$
\operatorname{tr}\left(\tau \cdot \mathrm{ps}_{m} \cdot B[z]\right)=\operatorname{tr}\left(C[z] \iota \cdot \operatorname{pr}_{m} \cdot A[z]\right)=\operatorname{tr}\left(\iota \operatorname{pr}_{m} A[z] C[z]\right)
$$

which - by Proposition 5.9 - is

$$
\frac{1}{2}(\operatorname{tr} A(1) C(1)+\operatorname{tr} A(-1) C(-1))=\frac{1}{2} \operatorname{tr}(B(1) P(1)+\operatorname{tr} B(-1) P(-1)) ;
$$

note that $D(z)$ is regular at $z= \pm 1$, since so is $P(z)$.
If $F \in M(k, \mathbb{C})$ and $Y \subset \mathbb{C}^{k}$ is an $F$-invariant subspace, write $\operatorname{tr}(F, Y)$ for the trace of $F$ on $Y$.

Proposition 5.16. Fix $P(z) \in \mathrm{GL}(k, \mathbb{C}(z))$ with $P\left(z^{-1}\right)=P(z)^{-1}$. Suppose that the function $P(z)$ is regular on $|z|=1$ and at $z=0$, and that it has order 1 at all its poles $\zeta_{1}, \ldots, \zeta_{s}$ inside $\{z \in \mathbb{C} ; 0<|z|<1\}$. Denote by $Y_{i}$ the image of the operator $\lim _{z \rightarrow \zeta_{i}}\left(z-\zeta_{i}\right) P(z)$ acting on $\mathbb{C}^{k}$. Fix $B(z) \in M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right)$ and suppose $B_{1}(z)=P(z)^{-1} B(z) P(z) \in M\left(k, \mathbb{C}\left[z, z^{-1}\right]\right)$. Then

$$
\begin{aligned}
\operatorname{tr}\left(\mathrm{ps}_{m} \cdot[B]\right)= & (2 m-1) \operatorname{tr} B_{0}+\frac{1}{2 \pi i} \int_{|z|=1} \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z \\
& +\sum_{1 \leq i \leq s} \operatorname{tr}\left(B\left(\zeta_{i}\right)+B_{1}\left(\zeta_{i}^{-1}\right), Y_{i}\right)
\end{aligned}
$$

with $B_{0}$ being the constant term of $B(z)$.
If in addition $B_{1}(z)=B\left(z^{-1}\right)$ then

$$
\begin{aligned}
& \operatorname{tr}\left(\pi_{m} B_{L}\right)=\left(m-\frac{1}{2}\right) \operatorname{tr} B_{0}+\frac{1}{4 \pi i} \int_{|z|=1} \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z \\
& \quad+\sum_{1 \leq i \leq s} \operatorname{tr}\left(B\left(\zeta_{i}\right), Y_{i}\right)+\frac{1}{4}[\operatorname{tr}(B(1) P(1))+\operatorname{tr}(B(-1) P(-1))]
\end{aligned}
$$

Note that the subspace $Y_{i}$ of $\mathbb{C}^{k}$ is invariant under $B\left(\zeta_{i}\right)$ and $B_{1}\left(\zeta_{i}^{-1}\right)$.
Proof. In view of Propositions 5.13 and 5.15 it suffices to verify that

$$
\begin{aligned}
& \frac{1}{2 \pi i} \oint_{|z|=1} \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z+\sum_{1 \leq i \leq s} \operatorname{tr}\left(B\left(\zeta_{i}\right)+B_{1}\left(\zeta_{i}^{-1}\right), Y_{i}\right) \\
& \quad=\operatorname{tr}\left([B], S / \mathbb{C}\left[z, z^{-1}\right]^{k}\right)-\operatorname{res}_{z=\infty} \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z,
\end{aligned}
$$

where

$$
S=\mathbb{C}\left[z, z^{-1}\right]^{k}+P(z) \mathbb{C}\left[z, z^{-1}\right]^{k} .
$$

For any $\zeta \neq 0$ in $\mathbb{C}$ denote by $M_{\zeta}$ and $N_{\zeta}$ the $\zeta$-components of the $\mathbb{C}\left[z, z^{-1}\right]$ modules $S / \mathbb{C}\left[z, z^{-1}\right]^{k}$ and $S / P(z) \mathbb{C}\left[z, z^{-1}\right]^{k}$, respectively. From Cauchy's formula and Lemma 5.14, it follows that

$$
\begin{aligned}
& \frac{1}{2 \pi i} \oint_{|z|=1} \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z=\sum_{1<|\zeta|<\infty} \operatorname{tr}\left([B], M_{\zeta}\right) \\
& -\sum_{1<|\zeta|<\infty} \operatorname{tr}\left([B], N_{\zeta}\right)-\operatorname{res}_{z=\infty} \operatorname{tr}\left(B(z) P^{\prime}(z) P(z)^{-1}\right) d z
\end{aligned}
$$

On the other hand, $\operatorname{tr}\left([B], S / \mathbb{C}\left[z, z^{-1}\right]^{k}\right)=\sum_{\zeta \in \mathbb{C}^{\times}} \operatorname{tr}\left([B], M_{\zeta}\right)$. Hence the required identity follows from

$$
\begin{aligned}
\sum_{0<|\zeta|<1} \operatorname{tr}\left([B], M_{\zeta}\right) & =\sum_{1 \leq i \leq s} \operatorname{tr}\left(B\left(\zeta_{i}\right), Y_{i}\right), \\
\sum_{1<|\zeta|<\infty} \operatorname{tr}\left([B], N_{\zeta}\right) & =\sum_{1 \leq i \leq s} \operatorname{tr}\left(B_{1}\left(\zeta_{i}^{-1}\right), Y_{i}\right) .
\end{aligned}
$$

If $P(z)$ is regular at $\zeta$ then $M_{\zeta}=0$. At each $\zeta_{i}, P(z)$ has a pole of order one. Hence there exists isomorphisms $M_{\zeta_{i}} \tilde{\rightarrow} Y_{i}$ which translate the operator $[B]: M_{\zeta_{i}} \rightarrow M_{\zeta_{i}}$ to the operator of multiplication by $B\left(\zeta_{i}\right)$ on $Y_{i}$. This implies the first identity.
For the second identity, for any $\zeta \in \mathbb{C}^{\times}$, denote by $\bar{N}_{\zeta}$ the $\zeta$-component of the module $\left(\mathbb{C}\left[z, z^{-1}\right]^{k}+P(z)^{-1} \mathbb{C}\left[z, z^{-1}\right]^{k}\right) / \mathbb{C}\left[z, z^{-1}\right]^{k}$. Multiplication by $P(z)^{-1}$ induces an isomorphism $N_{\zeta} \stackrel{\sim}{\rightarrow} \bar{N}_{\zeta}$. Under this isomorphism, multiplication by $B: N_{\zeta} \rightarrow N_{\zeta}$ translates into multiplication by $B_{1}: \bar{N}_{\zeta} \rightarrow \bar{N}_{\zeta}$, hence $\operatorname{tr}\left([B], N_{\zeta}\right)=\operatorname{tr}\left(\left[B_{1}\right], \bar{N}_{\zeta}\right)$. From $P(z)^{-1}=P\left(z^{-1}\right)$ we deduce that $\bar{N}_{\zeta}=0$ if $P(z)$ is regular at $z=\zeta^{-1}$, and that $\operatorname{tr}\left(\left[B_{1}\right], \bar{N}_{\zeta_{i}^{-1}}\right)=\operatorname{tr}\left(B_{1}\left(\zeta_{i}^{-1}\right), Y_{i}\right)$. This implies the second identity, hence the proposition.

### 5.3 Spectral terms

To deduce the trace formula from Proposition 5.16, we use properties of the function $M\left(\mu_{1}, \mu_{2}, t\right)$.
Recall that we have the projection $\pi_{m}: L \rightarrow L$ with kernel $L^{m}$ and image $L_{m}$, and $B_{L}$ denotes the operator of multiplication by $B(z)$ on $L$. The operator $P(z)$ is the restriction to the subspace of $U$-invariant vectors of the operator $M$ on the space $I_{0}=\oplus I_{0}\left(\mu_{1}, \mu_{2}\right)\left(\mu_{1}, \mu_{2}\right.$ range over the characters of $\left.\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}\right)$, which maps $I_{0}\left(\mu_{1}, \mu_{2}\right)$ to $I_{0}\left(\mu_{2}, \mu_{1}\right)$ via $M\left(\mu_{1}, \mu_{2}, z\right)$.

Proposition 5.17. There exists $a_{f} \in \mathbb{C}$ such that for sufficiently large $m$,

$$
\operatorname{tr}\left(\pi_{m} B_{L}\right)=\left(m-\frac{1}{2}\right) a_{f}-\sum_{5 \leq i \leq 8} S_{i}(f)
$$

Proof. By Proposition 4.29 the function $P(z)$ has two poles in the domain $|z| \leq 1$, namely at $z= \pm q^{-1 / 2}$, each of order 1 . We have $P\left(z^{-1}\right)=P(z)^{-1}$ and $P(z)^{-1} B(z) P(z)=B\left(z^{-1}\right)$. Hence the final claim of Proposition 5.16 applies and implies that for large enough $m$,

$$
\begin{aligned}
\operatorname{tr}\left(\pi_{m}[B]\right) & =\left(m-\frac{1}{2}\right) \operatorname{tr} B_{0}+\frac{1}{4 \pi i} \oint_{|z|=1} \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z+\operatorname{tr}\left(B\left(q^{-1 / 2}\right), Y_{+}\right) \\
& +\operatorname{tr}\left(B\left(-q^{-1 / 2}\right), Y_{-}\right)+\frac{1}{4}[\operatorname{tr}(B(1) P(1))+\operatorname{tr}(B(-1) P(-1))] .
\end{aligned}
$$

Here $B_{0}$ is the constant term of $B(z)$ and the image of the operator $\lim _{z \rightarrow \pm q^{-1 / 2}}\left(z \mp q^{-1 / 2}\right) P(z)$ is denoted by $Y_{ \pm}$. The proposition follows once we show that

$$
\begin{align*}
& \oint_{|z|=1} \operatorname{tr} B(z) P^{\prime}(z) P(z)^{-1} d z=-4 \pi i\left(S_{5}(f)+S_{6}(f)\right),  \tag{1}\\
& \operatorname{tr}\left(B\left(q^{-1 / 2}\right), Y_{+}\right)+\operatorname{tr}\left(B\left(-q^{-1 / 2}\right), Y_{-}\right)=-S_{8}(f),  \tag{2}\\
& \operatorname{tr}(B(1) P(1))+\operatorname{tr}(B(-1) P(-1))=-4 S_{7}(f) . \tag{3}
\end{align*}
$$

Denote by $r(z)$ the representation of $\mathrm{GL}(2, \mathbb{A})$ by right translation in $I(z)=$ $\otimes_{\mu_{1}, \mu_{2}} I\left(\mu_{1} \nu_{z^{-1}}, \mu_{2} \nu_{z}\right)$. Here $\mu_{1}, \mu_{2}$ are characters of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$. Let $r(z, f)$ be the convolution operator defined by $r(z)$ and the compactly supported function $f$ in $C_{c}^{\infty}(\mathrm{GL}(2, \mathbb{A}))$. Identify, as usual, $I(z)$ to the space $I_{0}$, and consider $r(z, f)$ as an operator in $I_{0}$. From Proposition 4.36, $B(z)$ coincides with the restriction of $r(z, f)$ to $I_{0}^{U}$. Also, $P(z)$ coincides with the restriction of $M(z)$ to $I_{0}^{U}$. Hence the integral on the left of (1) equals

$$
\begin{gathered}
\oint_{|z|=1} \operatorname{tr} r(z, f) M^{\prime}(z) M(z)^{-1} d z \\
=\sum_{\mu_{1}, \mu_{2}} \oint_{|z|=1} \operatorname{tr} I\left(\mu_{2} \nu_{z}-1, \mu_{1} \nu_{z}, f\right) M^{\prime}\left(\mu_{1}, \mu_{2}, z\right) M\left(\mu_{1}, \mu_{2}, z\right)^{-1} d z \\
=\sum_{\mu_{1}, \mu_{2}} \oint_{|z|=1} \operatorname{tr} M\left(\mu_{1}, \mu_{2}, z\right)^{-1} I\left(\mu_{2} \nu_{z^{-1}}, \mu_{1} \nu_{z}, f\right) M^{\prime}\left(\mu_{1}, \mu_{2}, z\right) d z \\
=\sum_{\mu_{1}, \mu_{2}} \oint_{|z|=1} \operatorname{tr} I\left(\mu_{1} \nu_{z}, \mu_{2} \nu_{z^{-1}}, f\right) M\left(\mu_{1}, \mu_{2}, z\right)^{-1} M^{\prime}\left(\mu_{1}, \mu_{2}, z\right) d z
\end{gathered}
$$

Then (1) follows from Proposition 4.9.
For (2), it follows from Proposition 4.29 that $Y_{+}=L^{U}$, with $L=\oplus L_{\mu}, L_{\mu} \subset$ $I(\mu, \mu)$ being generated by the function $x \mapsto \mu(x)$. The operator $r\left(q^{-1 / 2}, f\right)$ acts in $L_{\mu}$ as the operator of multiplication by $\int_{\mathrm{GL}(2, \mathbb{A})} f(x) \mu(\operatorname{det} x) d x$. Hence

$$
\operatorname{tr}\left(B\left(q^{-1 / 2}\right), Y_{+}\right)=\operatorname{tr}\left(r\left(q^{-1 / 2}, f\right), L\right)=\sum_{\mu} \int_{\mathrm{GL}(2, \mathbb{A})} f(x) \mu(\operatorname{det} x) d x
$$

where $\mu$ ranges over the set of characters of $\mathbb{A}^{\times} / F^{\times} \cdot \alpha^{\mathbb{Z}}$. Similarly

$$
\begin{aligned}
& \operatorname{tr}\left(B\left(q^{-1 / 2}\right), Y_{-}\right)=\operatorname{tr}\left(r\left(-q^{-1 / 2}, f\right), L\right) \\
= & \sum_{\mu} \int_{\mathrm{GL}(2, \mathbb{A})} f(x) \mu(\operatorname{det} x) \nu_{-1}(\operatorname{det} x) d x .
\end{aligned}
$$

Every character of $\mathbb{A}^{\times}$which is trivial on $F^{\times} \cdot \alpha^{2 \mathbb{Z}}$ is either trivial on $F^{\times} \cdot \alpha^{\mathbb{Z}}$ or its product with $\nu_{-1}$ is, so (2) follows.
For (3) note that

$$
\operatorname{tr} B(1) P(1)=\operatorname{tr} r(1, f) M(1)=\sum_{\mu} \operatorname{tr} I(\mu, \mu, f) M(\mu, \mu, 1)=-\sum_{\mu} \operatorname{tr} I(\mu, \mu, f)
$$

by Proposition 4.30. Similarly $\operatorname{tr} B(-1) P(-1)=-\sum_{\mu} \operatorname{tr} I\left(\mu \nu_{-1}, \mu \nu_{-1}, f\right)$.
This completes the proof of the trace formula.

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