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## K. F. LAI <br> Tamagawa number of reductive algebraic groups

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## Numdam

# TAMAGAWA NUMBER OF REDUCTIVE ALGEBRAIC GROUPS 

K.F. Lai

## 0. Introduction

The purpose of this paper is to give a formula for the Tamagawa number of a reductive quasi-split algebraic group $G$ defined over an algebraic number field in terms of the Tamagawa number of a maximal torus of $G$ (cf. Theorem 7.1).

The Tamagawa numbers of classical groups were determined by Weil [23]. In [15] Langlands determined the Tamagawa number of all split semisimple groups. We extend the result of Langlands to quasisplit groups.

I am most grateful to R.P. Langlands for explaining his methods to me. I would like to thank M. Rapoport for sending me his paper [18] and J. Arthur for useful suggestions.

Notations:
$F=$ number field
$F_{v}=$ completion of $F$ at the place $v$
$\bar{F}=$ algebraic closure of $F$
$v \mid \infty=v$ is an infinite place
$v<\infty=v$ is a finite place
$0_{v}=0_{F_{v}}=$ ring of integers of $F_{v}(v<\infty)$
$q=$ order of residue field of $F_{v}$
$\tilde{\omega}_{v}=$ uniformizing element of $\boldsymbol{0}_{v}(v<\infty)$
$\mathbb{A}=$ adeles of $\boldsymbol{F}, \mathbb{A}_{\mathscr{S}}=$ adeles trivial outside $\mathscr{S}$
$\|_{v}=$ normalised absolute value at $v(v<\infty):\left|\tilde{\omega}_{v}\right|_{v}=q^{-1}$
$\|=$ adelic absolute value.

For an algebraic group $H$ defined over $F$, we write

$$
\begin{aligned}
H_{v} & =H\left(F_{v}\right) \\
H_{f} & =\left\{\left(h_{v}\right) \in H(\mathbb{A}) \mid h_{v}=1 \text { if } v \mid \infty\right\} \\
H_{\infty} & =\prod_{v \mid \infty} H_{v} \\
H_{\mathscr{G}} & =\left\{\left(h_{v}\right) \in H(\mathbb{A}) \mid h_{v}=1 \text { if } v \notin \mathscr{S}\right\} \\
H^{\mathscr{\varphi}} & =\left\{\left(h_{v}\right) \in H(\mathbb{A}) \mid h_{v} \in H\left(0_{v}\right) \text { if } v \notin \mathscr{S}\right\} .
\end{aligned}
$$

For a complex valued function $f(x)$, write $\bar{f}(x)$ for the complex conjugate of $f(x)$.

## 1. Quasi-split algebraic groups

1.1. Let $G$ be a connected reductive algebraic group defined over $F$. We say that $G$ is quasi-split if one of the following equivalent conditions is satisfied
(I) $G$ has a Borel subgroup $B$ defined over $F$,
(II) the centralizer in $G$ of a maximal $F$-split torus is a maximal torus of $G$,
(III) $G$ has no anisotropic roots.

In the following $G$ denotes a connected reductive quasi-split group.
1.2. Let $A$ be a maximal torus of $G$ lying in $B$ and defined over $F$, $L$ the group of characters of $A, \hat{L}=\operatorname{Hom}(L, Z), \Sigma(\hat{\Sigma})$ the set of roots (coroots) of $G$ with respect to $A, \Delta$ basis of $\Sigma$ with respect to $B$ and $\hat{\Delta}$ the elements of $\hat{\Sigma}$ corresponding to $\Delta$. There is a bijection between $\bar{F}$-isomorphism classes of triple $(G, B, A)$ and isomorphism classes of based root system $\psi_{0}(G)=(L, \Delta, \hat{L}, \hat{\Delta})$. This bijection yields a connected reductive $C$-group $\hat{G}^{0}$ with based root system $\psi_{0}\left(\hat{G}^{0}\right)=$ ( $\hat{L}, \hat{\Delta}, L, \Delta$ ). Let $\hat{A}^{0}$ (resp. $\hat{B}^{0}$ ) be the maximal torus (resp. Borel subgroup) defined by $\psi_{0}\left(\hat{G}^{0}\right)$.

Let $E$ be a Galois extension of $F$ such that $G$ splits over $E$. If $\sigma \in \operatorname{Gal}(E / F), \lambda \in L$, we denote the action of $\sigma$ on $\lambda$ by $\sigma \lambda$ where $\sigma \lambda(a)=\sigma\left(\lambda\left(\sigma^{-1} a\right)\right)$ for $a \in A$. As $G$ is quasi-split, $\sigma \Delta=\Delta$. We can define a homomorphism $\mu: \operatorname{Gal}(E / F) \rightarrow \operatorname{Aut} \psi_{0}(G)$. Since we have canonical Aut $\psi_{0}(G)=\operatorname{Aut} \psi_{0}\left(\hat{G}^{0}\right)$, we may view $\mu$ as a homomorphism of $\operatorname{Gal}(E / F)$ into Aut $\psi_{0}\left(\hat{G}^{0}\right)$. Moreover there is a split exact sequence

$$
\begin{equation*}
(1) \rightarrow \text { Int } \hat{G}^{0} \rightarrow \text { Aut } \hat{G}^{0} \rightarrow \text { Aut } \psi_{0}\left(\hat{G}^{0}\right) \rightarrow(1) \tag{1}
\end{equation*}
$$

and a splitting yields a monomorphism

$$
\text { Aut } \psi_{0}\left(\hat{G}^{0}\right) \rightarrow \operatorname{Aut}\left(\hat{G}^{0}, \hat{B}^{0}, \hat{A}^{0}\right)
$$

Together with the $\mu$ above we get a homomorphism

$$
\mu^{\prime}: \operatorname{Gal}(E / F) \rightarrow \operatorname{Aut}\left(\hat{G}^{0}, \hat{B}^{0}, \hat{A}^{0}\right)
$$

The associated group to, or $L$-group of, $G$ is then by definition the semidirect product

$$
\hat{G}=\hat{G}^{0} \rtimes \operatorname{Gal}(E / F)
$$

(See Borel [3]).
1.3. Let $Z$ be the identity component of the centre of $G$ and $G^{\prime}$ be the derived group of $G$. Then $G=Z G^{\prime}$ and $A=Z A^{\prime}$ where $A^{\prime}=$ $A \cap G^{\prime}$. Let ${ }^{0} L^{+}$be the group of rational characters of $Z$ and ${ }^{0} L^{-}$be the elements of ${ }^{0} L^{+}$which are 1 on $Z \cap A^{\prime}$. Let ${ }^{1} L^{-}$be the lattice of roots of $A^{\prime}$. (Note that there is a bijection between the roots of $(G, A)$ and ( $G^{\prime}, A^{\prime}$ ) and the corresponding Weyl groups can be identified. We shall not use a separate notation.) We denote the Weyl group of the root system by $W$. There exists a non-degenerate $W$-invariant bilinear form (., .) on ${ }^{1} L^{-} \otimes_{z} C$ such that its restriction to ${ }^{1} L^{-} \otimes_{z} R$ is positive definite. Let ${ }^{1} L$ be the lattice of rational characters of $A^{\prime}$ and

$$
{ }^{1} L^{+}=\left\{\lambda \in{ }^{1} L^{-} \bigotimes_{Z} C \left\lvert\, \frac{2(\lambda, \alpha)}{(\alpha, \alpha)} \in Z\right. \text { for all roots } \alpha\right\} .
$$

Set $L^{-}={ }^{0} L^{-} \oplus^{1} L^{-}$and $L^{+}={ }^{0} L^{+} \oplus^{1} L^{+}$. We define dual lattices by

$$
\begin{aligned}
\hat{L}^{+} & =\operatorname{Hom}\left(L^{-}, \mathbf{Z}\right)=\operatorname{Hom}\left({ }^{0} L^{-}, \mathbf{Z}\right) \oplus \operatorname{Hom}\left({ }^{1} L^{-}, Z\right)={ }^{0} \hat{L}^{+} \oplus^{1} \hat{L}^{+} \\
\hat{L} & =\operatorname{Hom}(L, Z) \\
\hat{L}^{-} & =\operatorname{Hom}\left(L^{+}, Z\right)=\operatorname{Hom}\left({ }^{0} L^{+}, Z\right) \oplus \operatorname{Hom}\left({ }^{1} L^{+}, Z\right)={ }^{0} \hat{L}^{-} \oplus^{1} \hat{L}^{-}
\end{aligned}
$$

We then have $L^{-} \subset L \subset L^{+} \subset L \otimes_{\mathrm{z}} C$ and $\hat{L}^{-} \subset \hat{L} \subset \hat{L}^{+} \subset \hat{L} \otimes_{\mathrm{z}} C$.
For the pairing $L \times \hat{L} \rightarrow C$, we use the notation $\langle\lambda, \hat{\lambda}\rangle=\hat{\lambda}(\lambda)$ where $\lambda \in L, \hat{\lambda} \in \hat{L}$ and we extend it meaningfully to the other lattices. The
form on ${ }^{1} \hat{L}^{+} \otimes C$ adjoint to the one given above on ${ }^{1} L^{-} \otimes C$ will also be denoted by (.,.), i.e. if $\mu, \nu \in{ }^{1} L^{-} \otimes C$, and if the elements $\hat{\mu}, \hat{v}$ of ${ }^{1} \hat{L}^{+} \otimes \mathbb{C}$ satisfy the equations

$$
\langle\lambda, \hat{\mu}\rangle=(\lambda, \mu) \quad \text { and } \quad\langle\lambda, \hat{\nu}\rangle=(\lambda, \nu)
$$

for all $\lambda \in{ }^{1} L^{-} \otimes C$, then $(\mu, \nu)=(\hat{\mu}, \hat{\nu})$.
Suppose $v$ is a finite place of $F$. We define a map $\nu: A\left(F_{v}\right) \rightarrow \hat{L} \otimes \mathbf{Q}$ by the condition

$$
\begin{equation*}
|\lambda(a)|_{v}=\left|\tilde{\omega}_{v}\right|_{v}^{|\lambda, \nu(a)\rangle} \tag{2}
\end{equation*}
$$

for all $\lambda \in L$ and $a \in A\left(F_{v}\right)$, where $\tilde{\omega}_{v}$ is the uniformizing element of $F_{v}$ and $\mid \|_{v}$ is the normalized valuation of $F_{v}$. For $\mu \in L \otimes C$, define $\hat{t}_{\mu} \in \hat{A}^{0}=\operatorname{Hom}\left(\hat{L}, \mathbb{C}^{*}\right)$ by

$$
\begin{equation*}
\hat{t}_{\mu}(\tilde{\lambda})=\left|\tilde{\omega}_{v}\right|_{v}^{\langle\mu, \hat{\lambda}\rangle} \tag{3}
\end{equation*}
$$

for all $\hat{\lambda} \in \hat{L}$. We sometimes write $\hat{t}$ for $\hat{\boldsymbol{t}}_{\mu}$.
We write $L_{F}$ for the lattice of $F$-rational characters of $A$. Similar notation will be used for the lattices ${ }^{0} L^{+}$etc.
1.4. Next we write down explicitly the Galois action on the derived group $\hat{G}^{\prime}$ of $\hat{G}^{0}$. Put $\hat{A}^{\prime}=\hat{A}^{0} \cap \hat{G}^{\prime}$. Let $\hat{a}$ be the Lie algebra of $\hat{A}^{\prime}$. Choose $H_{1}, \ldots, H_{r} \in \hat{a}$ so that

$$
\lambda\left(H_{i}\right)=\left\langle\alpha_{i}, \lambda\right\rangle
$$

where $\lambda \in{ }^{1} \hat{L}^{+}$and $\Delta=\left\{\alpha_{1}, \ldots, \alpha_{r}\right\}$ are the simple roots. Choose vectors $X_{ \pm \hat{\alpha}_{i}}$ belong to the $\pm \hat{\alpha}_{i}$ respectively such that

$$
\left[\boldsymbol{X}_{\hat{\alpha}_{i}}, X_{-\hat{\alpha}_{i}}\right]=\boldsymbol{H}_{i}
$$

For $\sigma \in \operatorname{Gal}(E / F), \widehat{\sigma \alpha}=\sigma \alpha$ for $\alpha \in \Delta$. If we put $\sigma\left(\hat{\alpha}_{i}\right)=\hat{\alpha}_{\sigma(i)}$, then the Galois action on the Lie algebra $\hat{\mathfrak{g}}^{\prime}$ of $\hat{G}^{\prime}$ is the unique isomorphism satisfying

$$
\sigma\left(H_{i}\right)=H_{\sigma(i)}, \quad \sigma X_{ \pm \hat{\alpha}_{i}}=X_{ \pm \hat{\alpha}_{\sigma(i)}}
$$

(see Jacobson [9] Chap. VII).
1.5. Let $\Sigma_{F}$ denote the set of $F$-roots of $G$ with respect to $A_{d}$, the
maximal $F$-split torus in $A$. As $G$ is quasi-split, each element of $\Sigma$ has a nontrivial restriction to $A_{d}$, and $\Sigma_{F}$ is equal to the set of restriction to $A_{d}$ of elements of $\Sigma$. In fact, if $G$ splits over a Galois extension $E$ of $F$, the Galois group $\operatorname{Gal}(E / F)$ acts on $\Sigma$ and each orbit restricts to an element of $\Sigma_{F}$. In each orbit choose a representative $\alpha$ and denote the corresponding orbit by $\mathcal{O}_{\alpha}$ and the element in $\Sigma_{F}$ to which the elements in $\mathscr{O}_{\alpha}$ restrict, is denoted by $\alpha_{F}$, i.e. $\alpha_{F}=\alpha \mid A_{d}$.

The Weyl group $W$ of $\Sigma$ is given by $N(A) / Z(A)$ while the rational Weyl group $W_{F}$ of $\Sigma_{F}$ is $N\left(A_{d}\right) / Z\left(A_{d}\right)$. We can identify $W_{F}$ as a subgroup of $W$.

Let ${ }_{0} \Sigma_{F}$ be the reduced $F$-root system consisting of the indivisible $F$-roots of $\Sigma_{F}$, i.e. ${ }_{0} \Sigma_{F}=\left\{\alpha_{F} \in \Sigma_{F} \left\lvert\, \frac{1}{2} \alpha_{F} \notin \Sigma_{F}\right.\right\}$. ${ }_{0} \Sigma_{F}^{+}={ }_{0} \Sigma_{F} \cap \Sigma_{F}^{+}$.

Next we define the elementary subgroup $G_{\alpha_{F}}$ of $G$ for $\alpha_{F} \in{ }_{0} \Sigma_{F}^{+}$. Let $A_{\alpha_{F}}=\left(\operatorname{ker} \alpha_{F}\right)^{0}$. Then $G_{\alpha_{F}}=Z_{G} A_{\alpha_{F}}$, i.e. we take the centralizer in $G$ of $A_{\alpha_{F}}$

It can be easily proved that $G_{\alpha_{F}}$ is connected reductive quasi-split group of semi simple $F$-rank 1.
1.6. There is a non-empty finite set $\mathscr{S}$ of places of $F$, containing all the infinite places such that the $F$-group $G$ can be regarded as defined above $\operatorname{Spec}\left(0_{\mathscr{y}}\right)$, where $0_{\mathscr{g}}$ is the ring of the elements of $F$ which are integral outside $\mathscr{S}$. Thus $G\left(0_{v}\right)$ is defined for those $v$ not in $\mathscr{S}$.

For $v \mid \infty$, let $K_{v}$ be a maximal compact subgroup of $G_{v}$ such that $G_{v}=B_{v} \cdot K_{v}$ is an Iwasawa decomposition. For $v<\infty$, let $K_{v}$ be a special open maximal compact subgroup of $G_{v}$, in the sense of Bruhat-Tits [4]. In particular, for almost all $v, K_{v}$ can be taken to be $G\left(0_{v}\right)$. Similar considerations can be given to $G_{\alpha_{F}}$ Therefore, when we consider the finite set $\left\{G, G_{\alpha_{F}}\right\}_{\alpha_{F} \in_{0} \Sigma_{F}}$ of groups taken together, except for a finite number of places, we have simultaneously

$$
\begin{align*}
G_{v} & =B_{v} G\left(0_{v}\right) \\
G_{\alpha_{F}}\left(F_{v}\right) & =B_{\alpha_{F}}\left(F_{v}\right) G_{\alpha_{F}}\left(0_{v}\right) \tag{4}
\end{align*}
$$

where $\alpha_{F} \in_{0} \Sigma_{F}$.

Let us now fix $K_{f}=\Pi_{v<\infty} K_{v}, \quad K_{\infty}=\Pi_{v \mid \infty} K_{v}, \quad K=K_{\infty} K_{f}$. Then $G(\mathbb{A})=B(\mathbb{A}) \cdot K$.
1.7. Let $X(G)$ be the lattice of rational characters on $G$. Let $L(s, G)$ be the Artin $L$-function corresponding to the $\operatorname{Gal}(E / F)$-module $X(G) \otimes Q$ and let $L_{v}(s, G)$ be its $v$-component.

Let $\chi$ be a nontrivial character on $\mathbb{A}$ trivial on $F . \chi$ defines a
nontrivial character $\chi_{v}$ of $F_{v}$ at each place $v$ of $F$. Let $\mathrm{d} x_{v}$ be the additive Haar measure on $F_{v}$ self-dual with respect to $\chi_{v}$ and let $\mathrm{d} x=\Pi_{v} \mathrm{~d} x_{v}$. For $v$ finite, the Haar measure on $F_{v}^{x}$ is chosen so that the measure of $0_{v}^{x}$ is one.

Let $\omega$ be an $F$-rational left-invariant nowhere vanishing exterior form of highest degree on $G$. For each $v, \omega$ and $\mathrm{d} x_{v}$ defines a measure $|\omega|_{v}$ on $G_{v}$ (cf. [23]). We put $\mathrm{d} g_{v}=L_{v}(1, G)|\omega|_{v}$, for finite $v$, and $\mathrm{d} g_{v}=|\omega|_{v}$ for infinite $v$. Then the Tamagawa measure $\mathrm{d} g$ on $G(\mathbb{A})$ is the Haar measure on $G(\mathbb{A})$ defined by

$$
\begin{equation*}
\mathrm{d} g=\lim _{s \rightarrow 1} \frac{1}{(s-1)^{r} L(s, G)} \prod_{v} \mathrm{~d} g_{v} \tag{5}
\end{equation*}
$$

where $r$ the rank of the lattice of $F$-rational characters $X(G)_{F}$ of $G$ (cf. [17]). This measure is independent of choice of $\chi$ and $\omega$.

Let $\chi_{1}, \ldots, \chi_{r}$ a basis of $X(G)_{F}$. Then the map $g \rightarrow$ $\left(\left|\chi_{1}(g)\right|, \ldots,\left|\chi_{r}(g)\right|\right)$ defines a homomorphism $G(\mathbb{A}) \rightarrow\left(\mathbb{R}_{+}^{x}\right)^{r}$. Let $G^{1}(\mathbb{A})$ be the kernel of this homomorphism. Also, the restriction of $\chi_{1}, \ldots, \chi_{r}$ to the split component $Z_{d}$ of the radical of $G$ defines an $F$-homomorphism $\delta$ from $Z_{d}$ to GL(1)r. This defines a homomorphism $\delta_{\infty}$ from the identity component of $Z_{d \infty}$ to GL(1) $)_{\infty}^{r}$. For each $t \in \mathbf{R}_{+}^{x}$, call $\xi(t)$ the idele $\left(\xi(t)_{v}\right)$ such that $\xi(t)_{v}=1$ for every finite place and $\xi(t)_{v}=t$ for every infinite place. Then $t \rightarrow \xi(t)$ is an isomorphism of $\mathbf{R}_{+}^{x}$ onto a subgroup $\mathrm{GL}^{+}(1)_{\infty}$ of $\mathrm{GL}(1)_{\infty}$. Let $Z_{\infty}^{+}$be the identity component of inverse image of $\mathrm{GL}^{+}(1)_{\infty}^{r}$ under $\delta_{\infty}$. Then $Z_{\infty}^{+}$is isomorphic to $\left(\mathbf{R}_{+}^{x}\right)^{r}$ and $G(\mathbb{A})=G(\mathbb{A})^{1} \times Z_{\infty}^{+}$. If we put the measure $\mathrm{d} t=\wedge_{i=1}^{r}\left(\mathrm{~d} t_{i} / t_{i}\right)$ on $\mathbf{R}_{+}^{x}$, then

$$
\begin{equation*}
\mathrm{d} g=\mathrm{d} g^{1} \times \mathrm{d} t \tag{6}
\end{equation*}
$$

defines a Haar measure on $G^{1}(\mathbb{A})$. This measure is independent of choice of $\chi_{1}, \ldots, \chi_{r}$. The Tamagawa number $\tau(G)$ is the finite number defined by

$$
\begin{equation*}
\tau(G)=\int_{G(F) \mid G^{1}(A)} \mathrm{d} g^{1}=\int_{G(F) Z_{\infty}^{+} \mid G(A)} \mathrm{d} g . \tag{7}
\end{equation*}
$$

1.8. Let $N$ be the unipotent radical of $B$. Then we can define Tamagawa measures $\mathrm{d} a$ (resp. $\mathrm{d} n$ ) on $A(\mathbb{A})$ (resp. $N(\mathbb{A})$ ) as in the case of $G$. We normalize the measure on $K_{v}$ by the condition

$$
\int_{K_{v}} \mathrm{~d} k_{v}=1 .
$$

Then we have $\mathrm{d} k=\Pi_{v} \mathrm{~d} k_{v}$ and

$$
\int_{K} \mathrm{~d} k=1
$$

Let $\rho$ be the half sum of the positive roots of $G$ with respect to $A$. To simplify notation we write $\rho$ for the quasi-character on $A(F) \backslash A(\mathbb{A})$ determined by $\rho$. Since $G(\mathbb{A})=B(\mathbb{A}) \cdot K=N(\mathbb{A}) A(\mathbb{A}) K$, there exists a positive constant $\kappa$ such that for any $f \in C_{c}(G(\mathbb{A}))$,

$$
\begin{equation*}
\int_{G(\mathrm{~A})} f(g) \mathrm{d} g=\kappa \int_{N(\mathrm{~A}) A(\mathrm{~A}) K} f(n a k) \rho^{-2}(a) \mathrm{d} n \mathrm{~d} a \mathrm{~d} k \tag{8}
\end{equation*}
$$

According to the Bruhat decomposition of $G$ we have

$$
\begin{equation*}
G_{v}=\bigcup_{w \in W_{F_{v}}} N_{v} A_{v} w N_{v} . \tag{9}
\end{equation*}
$$

But except for the Weyl group element $w_{0}$ that sends all the positive roots to negative roots, the cosets NAwN has lower dimension than that of $G$, and so $N A w N$ has measure zero. Thus if we write $g_{v}=n_{v} a_{v} w_{0} n_{v}^{\prime}$, we have

$$
\begin{equation*}
\mathrm{d} g_{v}=\rho^{-2}(a) \mathrm{d} n_{v} \overline{\mathrm{~d}}_{v} \mathrm{~d} n_{v}^{\prime} \tag{10}
\end{equation*}
$$

where $\overline{\mathrm{d} a}_{v}$ is the local measure on $A_{v}$ induced by $|\omega|_{v}$.

## 2. Eisenstein series and $M(w, \lambda)$

2.1. For our purposes it is sufficient to consider the contribution to the spectral decomposition of $\mathscr{L}^{2}\left(Z_{\infty}^{+} G(F) \backslash G(\mathbb{A}) / K\right)$ from the Borel subgroup $B$. We can define the adelic analogue of the function spaces $\mathscr{E}(V, W), \mathscr{D}(V, W)$ and $\mathscr{H}(\mathscr{D}(V, W))$ of $\S 2$ and 3 of [13] with respect to the Borel subgroup $B$, the trivial representation of $K$ and a character $\lambda$ of $Z_{\infty}^{+} A(F) \backslash A(A)$ which is trivial on the image of $B(\mathbb{A}) \cap$ $K$ in $N(\mathbb{A}) \backslash B(\mathbb{A})$.
2.2. Define $A_{\infty}^{+}\left(\right.$resp. $\left.A(\mathbb{A})^{1}\right)$ in the same way as $Z_{\infty}^{+}\left(\right.$resp. $\left.G(\mathbb{A})^{1}\right)$. Let $\left(Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})\right)^{*}$ be the set of characters of $Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})$. Fix a basis $\left\{\chi_{i}\right\}$ of $L_{F}$. Each element $\lambda=\Sigma s_{i} \chi_{i}$ of $L_{F} \otimes C$ can be considered as a character of $Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})$ via the formula

$$
\lambda(a)=\prod_{i}\left|\chi_{i}(a)\right|_{s_{i}} .
$$

In this way $L_{F} \otimes \mathbb{C}$ is identified with a subset of $\left(Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})\right)^{*}$. From now on we shall consider only those $\lambda$ in $L_{F} \otimes \mathbb{C}$.

Let $\mathscr{E}(\lambda)$ be the space of continuous functions on $N(\mathbb{A}) B(F) \backslash G(\mathbb{A}) / K$ satisfying the condition

$$
\begin{equation*}
\Phi(a g)=\lambda(a) \rho(a) \Phi(g) \tag{1}
\end{equation*}
$$

for $a \in A(\mathbb{A}), g \in G(\mathbb{A})$.
Let $\mathscr{H}(\lambda)$ be the space of functions $\Phi(\cdot, g)$, with values in $\mathscr{E}(\lambda)$, which is defined and analytic in a tube in $L_{F} \otimes C$ over a ball of radius $R$ with $R>(\rho, \rho)^{1 / 2}$ and which goes to zero at infinity faster than the inverse of any polynomial.
2.3. Let $D_{0}$ be the unitary characters of $Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})$. Then $\left(Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})\right)^{*}$ is also the union of sets of the form

$$
D_{\sigma}=\left\{\chi \in\left(Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})\right)^{*}| | \chi \mid=\sigma\right\}
$$

where $\sigma$ is a fixed character with values in $R_{+}^{x}$. We equip $D_{0}$ with the dual Haar measure via Pontrjagin duality and give $D_{\sigma}$ the measure obtained by transport of structure from $D_{0}$.

We write $\mathscr{D}$ for the space spanned by functions of the form

$$
\begin{equation*}
\phi(g)=\int_{\operatorname{Re} \lambda=\lambda_{0}} \Phi(\lambda, g)|\mathrm{d} \lambda| \tag{2}
\end{equation*}
$$

where $\Phi \in \mathscr{H}(\lambda)$ and $\lambda_{0}$ is a character with values in $R_{+}^{x}$. By means of Fourier transform we get

$$
\begin{equation*}
\Phi(\lambda, g)=\int_{Z_{\infty}^{+} A(F) \backslash A(A)} \phi(a g) \lambda^{-1}(a) \rho^{-1}(a) \mathrm{d} a . \tag{3}
\end{equation*}
$$

According to Langlands [13, 14], for $\phi \in \mathscr{D}$ the theta series

$$
\begin{equation*}
\tilde{\phi}(g)=\sum_{\gamma \in P(F) \mid G(F)} \phi(\gamma g) \tag{4}
\end{equation*}
$$

belongs to $\mathscr{L}^{2}\left(Z_{\infty}^{+} G(F) \backslash G(\mathbb{A})\right)$. Combining with (2), we get

$$
\begin{equation*}
\tilde{\phi}(g)=\int_{\operatorname{Re} \lambda=\lambda_{0}} E(g, \Phi, \lambda) \mathrm{d} \lambda \tag{5}
\end{equation*}
$$

where

$$
\begin{equation*}
E(g, \Phi, \lambda)=\sum_{\gamma \in P(F) \backslash G(F)} \Phi(\lambda, \gamma g) \tag{6}
\end{equation*}
$$

is an Eisenstein series. It converges uniformly for $g$ in compact subsets of $G(\mathbb{A})$ and $\lambda \in L_{F} \otimes \mathbb{C}$ such that $\operatorname{Re}(\lambda, \alpha)>(\rho, \alpha)$ for every positive root $\alpha$.

We define the constant term of the Eisenstein series $E(g, \Phi, \lambda)$ by

$$
\begin{equation*}
E_{0}(g, \Phi, \lambda)=\int_{N(F) \backslash N(A)} E(n g, \Phi, \lambda) \mathrm{d} n \tag{7}
\end{equation*}
$$

2.4. Proposition: The constant term is given by the following formula:

$$
E_{0}(g, \Phi, \lambda)=\sum_{w \in W_{F}} M(w, \lambda) \Phi(\lambda, g)
$$

where $W_{F}$ is the F-rational Weyl group of $G$ and

$$
\begin{equation*}
M(w, \lambda) \Phi(\lambda, g)=\int_{w^{-1} B(F) w \cap N(F) \backslash N(A)} \Phi(\lambda, w n g) \mathrm{d} n \tag{8}
\end{equation*}
$$

Proof: We have

$$
E_{0}(g, \Phi, \lambda)=\int_{N(F) N(A)} \sum_{B(F) \backslash G(F)} \Phi(\lambda, \gamma n g) \mathrm{d} n
$$

The proposition is immediate once we break up the sum over $B(F) \backslash G(F)$ into a sum over $W_{F}=B(F) \backslash G(F) / N(F)$ (Bruhat decomposition) and a sum over ( $\left.w^{-1} B(F) w \cap N(F)\right) \backslash N(F)$.
2.5. We can define local version of $\mathscr{E}(\lambda)$ as the space $\mathscr{E}_{v}(\lambda)$ of continuous functions $\Phi_{v}$ on $N_{v} \backslash G_{v} / K_{v}$ satisfying

$$
\Phi_{v}\left(a_{v} g_{v}\right)=\lambda\left(a_{v}\right) \rho\left(a_{v}\right) \Phi\left(g_{v}\right)
$$

(here $\rho\left(a_{v}\right)$ is to be interpreted as $\left.\left|\rho\left(a_{v}\right)\right|_{v}\right)$.
For $\Phi \in \mathscr{E}(\lambda)$, we let $\Phi_{v}$ denote its restriction to $G_{v}$. Since $\Phi$ is right invariant under $K=\Pi K_{v}$ where $K_{v}=G\left(0_{v}\right)$ almost all $v$, and
$G(A)$ is the direct limit of $G^{\varphi}$, we can write

$$
\Phi(g)=\prod \Phi_{v}\left(g_{v}\right)
$$

(Here it is understood that $\Phi(1)=1$.)
Furthermore, $M(w, \lambda)$ is a linear map from $\mathscr{E}(\lambda)$ to $\mathscr{E}\left(\lambda^{w}\right)$ where $\lambda^{w}(a)=\lambda\left(w a w^{-1}\right)$. In fact it is just multiplication by a constant to be calculated below. Moreover, $M(1, \lambda)=1$ because $\operatorname{vol}(N(F) \backslash N(\mathbb{A}))=$ 1.
2.6. Proposition: Let ${ }^{w} N=w^{-1} N w \cap N$ and $N^{w}=w^{-1} \bar{N} w \cap N$ where $\bar{N}$ is the unipotent subgroup opposite to $N$. Define a linear transform $M_{v}(w, \lambda): \mathscr{E}_{v}(\lambda) \rightarrow \mathscr{E}_{v}\left(\lambda^{w}\right)$ by

$$
\begin{equation*}
M_{v}(w, \lambda) \Phi(g)=\int_{N_{v}^{w}} \Phi(w n g) \mathrm{d} n \tag{9}
\end{equation*}
$$

for $g \in G_{v}$. Then we have

$$
\begin{equation*}
M(w, \lambda)=\prod M_{v}(w, \lambda) \tag{10}
\end{equation*}
$$

(Here one regard the $M_{v}(w, \lambda)$ as complex numbers.)

Proof: First we have $N={ }^{w} N \cdot N{ }^{w}$. So

$$
{ }^{w} N(F) \backslash N(\mathbb{A})=\left({ }^{w} N(F) \backslash{ }^{w} N(\mathbb{A})\right) \cdot N^{w}(\mathbb{A}) .
$$

It follows that, for $\Phi \in \mathscr{E}(\lambda)$

$$
\begin{aligned}
M(w, \lambda) \Phi(g) & =\int_{w_{N(F) \backslash N(A)}} \Phi(\text { wng }) \mathrm{d} n \\
& =\int_{w_{N(F))^{w}(A)}} \int_{N^{w}(A)} \Phi\left(w n_{1} w^{-1} \cdot w n_{2} g\right) \mathrm{d} n_{2} \mathrm{~d} n_{1} .
\end{aligned}
$$

The formula (10) now follows from the above and the fact that we have normalized our measure such that

$$
\int_{\left.w_{N}(F)\right)^{w} N(A)} \mathrm{d} n_{1}=1 .
$$

## 3. $M_{v}(w, \lambda)$ in the rank one case

3.1. We shall compute $M_{v}(w, \lambda)$ for those places $v$ of $F$ satisfying the following conditions:
(i) $G$ is a connected reductive quasi-split group over $F_{v}$.
(ii) $G$ splits over an unramified extension of $F_{v}$.
(iii) $G_{v}=B_{v} K_{v}$ and $K_{v}=G\left(0_{v}\right)$.
(iv) $G$ is of semisimple $F_{v}$-rank one.

Let us write $E_{v}$ for the unramified extension of $F_{v}$ over which $G$ splits and write $\tilde{\omega}$ for the uniformizing element of both $E_{v}$ and $F_{v}$. We denote by $\sigma$ the Frobenius element in $\operatorname{Gal}\left(E_{v} / F_{v}\right)$.

Under the assumption, the $F_{v}$-rational Weyl group $W_{F_{v}}=\left\{1, w_{0}\right\}$, where $w_{0}$ sends all the positive roots to negative roots. We know that

$$
M_{v}(1, \lambda)=1 .
$$

It remains to calculate $M_{v}\left(w_{0}, \lambda\right)$. As $\mathscr{C}_{v}(\lambda)$ is one dimensional it suffices to calculate

$$
\begin{equation*}
M_{v}\left(w_{0}, \lambda\right)=M_{v}\left(w_{0}, \lambda\right) \Phi(\lambda, 1)=\int_{N_{v}^{w_{0}}} \Phi\left(\lambda, w_{0} n\right) \mathrm{d} n \tag{1}
\end{equation*}
$$

where $\Phi(\lambda)$ is $\mathscr{E}(\lambda)$ is chosen to satisfy

$$
\Phi(\lambda, 1)=1
$$

$G$ has $F_{v}$-rational rank 1 also implies that $L_{F_{v}} \otimes C$ is isomorphic to $C$ and hence can be replaced by the set $\left\{\rho^{s} \mid s \in C\right\}$. Thus it suffices to consider $M\left(w_{0}, \rho^{s}\right)$. We define $\Phi\left(\rho^{s}\right)$ by:

$$
\begin{gathered}
\Phi\left(\rho^{s}, a\right)=|\rho(a)|_{v}^{s+1} \quad \text { if } a \in A_{v} \\
\Phi\left(\rho^{s}, n g k\right)=\Phi\left(\rho^{s}, g\right) \quad \text { if } n \in N_{v}, k \in K_{v} .
\end{gathered}
$$

Let us write $M(s)$ for $M\left(w_{0}, \rho^{s}\right)$. Then (1) becomes

$$
M(s)=\int_{N_{0}^{w_{0}}} \rho^{s+1}\left(w_{0} n\right) \mathrm{d} n
$$

We can further assume that $w_{0} \in K_{v}$, then changing variable by the map $n \rightarrow w_{0} n w_{0}^{-1}$, we have

$$
\begin{equation*}
M(s)=\int_{\bar{N}_{v}} \rho^{s+1}(\bar{n}) \mathrm{d} \bar{n}, \tag{2}
\end{equation*}
$$

and

$$
\left.\rho^{s}(a)=\left(|\tilde{\omega}| \tilde{F}_{v}, \nu(a)\right\rangle\right)^{s}
$$

3.2. Proposition: Let $\hat{\mathfrak{n}}$ be the subspace of the Lie algebra of $\hat{\boldsymbol{G}}$ spanned by the positive root vectors. Then

$$
\begin{equation*}
M(s)=\frac{\operatorname{det}\left(I-\left.|\tilde{\omega}|_{F_{v}} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right)}{\operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\tilde{n}}\right)} \tag{3}
\end{equation*}
$$

where $\hat{t}=\hat{t}_{s \rho}$.
Let $G^{\prime}$ be the derived subgroup of $G$. Then the unipotent radical of the Borel subgroup of $G^{\prime}$ is the same as that of the corresponding Borel subgroup $B$ of $G$. Thus we only need to compute the integral $M(s)$ for connected semisimple quasi-split groups of $F_{v}$-rank one. Henceforth, in this subsection we shall assume $G$ to be of such type.

According to Steinberg's variation of Chevalley's theme, the quasisplit form of $G$ is determined up to $F_{v}$-isomorphism by its Dynkin diagram and the twisted action of galois group (modulo inner twisting). As a result, up to central isogeny, $G$ can only be one of the following types:
(I) $G$ splits over $G_{v}$ and has a connected Dynkin diagram, i.e. $G=\mathrm{SL}_{2}$.
(II) $G$ is a twisted form of a $F_{v}$-split group whose Dynkin diagram is type $A_{2}$, i.e. $G\left(F_{v}\right)=\operatorname{SU}_{3}\left(E_{v} / F_{v}\right)=\left(\left.g \in \operatorname{SL}_{3}\left(E_{v}\right)\right|^{t} \bar{g} J g=J\right\}$ where $E_{v} / F_{v}$ is a quadratic extension; the conjugation by the nontrivial element of the Galois group $\operatorname{Gal}\left(E_{v} / F_{v}\right)$ is denoted by $x \rightarrow \bar{x} ;{ }^{t} \bar{g}$ is the conjugate-transpose of the matrix $g: J=\left(\begin{array}{cc} & 1 \\ 1 & \end{array}\right)$ is the matrix of the Hermitian form with respect to the nontrivial element of $\operatorname{Gal}\left(E_{v} / F_{v}\right)$.
(III) $G$ is a twisted form of a $F_{v}$-split group whose Dynkin diagram consists of $n$ copies of $A_{1}$, i.e. there exists an extension $E_{v} / F_{v}$ of degree $n$ and $G\left(F_{v}\right)=\mathrm{SL}_{2}\left(E_{v}\right)$.
(IV) $G$ is a twisted form of $F_{v}$-split group whose Dynkin diagram consists of $n$ copies of $A_{2}$; there exists field extensions $E_{v}, E_{v}^{\prime}$ of $F$ such that $\left[E_{v}: E_{v}^{\prime}\right]=2,\left[E_{v}: F_{v}\right]=2 n$. If $x \rightarrow \bar{x}$ is the nontrivial action of the Galois group $\operatorname{Gal}\left(E_{v} / E_{v}^{\prime}\right)$ then $G\left(F_{v}\right)=\operatorname{SU}_{3}\left(E_{v} / E_{v}^{\prime}\right)=$
$\left\{\left.g \in \mathrm{SL}_{3}\left(E_{v}\right)\right|^{\prime} \bar{g} J g=J\right\}$ where $J=\left(\begin{array}{cr} & 1 \\ 1 & \end{array}\right)$.
It is obvious that it suffices to calculate (2) up to isogeny (see for example [18] §4.3). Moreover Rapoport [18] pointed out that it is possible to avoid the calculation of (2) for the cases (III) and (IV) by proving a general lemma on the behaviour of (2) under restriction of ground field.
3.3. When $G$ is $\mathrm{SL}_{2}$, it is well known that

$$
M(s)=\frac{1-q^{-(s+1)}}{1-q^{-s}}
$$

The Lie algebra $\hat{n}$ in this case is one dimensional and it is trivial to check the formula (3). We shall omit the details.
3.4. Proposition: Let $E_{v} / F_{v}$ be an unramified quadratic extension of local fields such that 2 is a unit in $E_{v}$. Then for the quasi-split group $\mathrm{SU}_{3}\left(E_{v} / F_{v}\right)$ we have

$$
M(s)=\frac{\left(1-q^{-2(s+1)}\right)\left(1+q^{-2 s-1}\right)}{\left(1-q^{-2 s}\right)\left(1+q^{-2 s}\right)}=\frac{\operatorname{det}\left(I-\left.|\tilde{\omega}|_{F_{v}} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right)}{\operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right)}
$$

Proof: First we have

$$
\begin{aligned}
& A\left(F_{v}\right)=\left\{\left(\begin{array}{lll}
a & & \\
& b & \\
& & \bar{a}^{-1}
\end{array}\right) \left\lvert\, \begin{array}{c}
a, b \in E_{v}^{x} \\
b \bar{b}=1, a b \bar{a}^{-1}=1
\end{array}\right.\right\}, \\
& N\left(F_{v}\right)=\left\{\left.\left(\begin{array}{rrr}
1 & x & y \\
& 1 & -\bar{x} \\
& & 1
\end{array}\right) \right\rvert\, \begin{array}{c}
y+\bar{y}+x \bar{x}=0 \\
x, y \in E_{v}
\end{array}\right\}, \\
& K=\operatorname{SU}_{3}\left(0_{E_{v}}\right) .
\end{aligned}
$$

$E$ is an unramified quadratic extension of $F$, so there exists an element $c \in 0_{F_{V}}-\tilde{\omega} 0_{F_{v}}$ such that its image in $0_{F_{v}} / \tilde{\omega} 0_{F_{v}}$ is not a square and $E_{v}=F_{v}(\sqrt{c})$. Let the map $\operatorname{ord}_{F_{v}}: F_{v}^{x} \rightarrow Z$ be defined by the condition

$$
|x|_{F_{v}}=|\tilde{\omega}|_{F_{v}}^{\operatorname{rdq}_{F_{v}^{x}}^{x}} \quad \text { for } x \in F_{v}^{x}
$$

Similar condition defines $\operatorname{ord}_{E_{v}}$. Note if $x \in F_{v}$, then $|x|_{E_{v}}=|x|_{F_{v}}^{2}$ implies $\operatorname{ord}_{F_{v}} x=\operatorname{ord}_{E_{v}} x$.

Next, let us determine the measure $\mathrm{d} n$ on the nilpotent group $N\left(F_{v}\right)$. Let $x, y \in E_{v}$ such that $y+\bar{y}+x \bar{x}=0$. Then we can write $y=y_{1} \sqrt{c}-\frac{x \bar{x}}{2}$ where $y_{1} \in F_{v}$. Note that $x \bar{x}=N_{E_{\nu} / F_{v}}(x)$ also belongs to $F_{v}$.

A typical element of $N\left(F_{v}\right)$ can now be written as

$$
\left(\begin{array}{rrr}
1 & x & y \\
& 1 & -\bar{x} \\
& & 1
\end{array}\right)=\left(\begin{array}{ccc}
1 & x & -\frac{x \bar{x}}{2} \\
& 1 & -\bar{x} \\
& & 1
\end{array}\right)\left(\begin{array}{lll}
1 & 0 & y_{1} \sqrt{c} \\
& 1 & 0 \\
& & 1
\end{array}\right)
$$

Thus we can write $N\left(F_{v}\right)=N_{1} N_{2}$ (as sets) and take $\mathrm{d} n$ to be the image of the product of the measure on $E_{v}$ and $F_{v}$ respectively under the maps;

$$
\begin{aligned}
& x \mapsto n_{1}=\left(\begin{array}{ccc}
1 & x & -\frac{x \bar{x}}{2} \\
& 1 & -\bar{x} \\
& & 1
\end{array}\right), \quad x \in E_{v}, \\
& y_{1} \mapsto n_{2}=\left(\begin{array}{lll}
1 & 0 & y_{1} \sqrt{c} \\
& 1 & 0 \\
& & 1
\end{array}\right), \quad y_{1} \in F_{v} .
\end{aligned}
$$

We normalize the measures on $E_{v}$ and $F_{v}$ by the condition that the volume of the respective maximal compact subrings is one.

The nontrivial element of the Weyl group corresponds to the matrix

$$
w_{0}=\left(\begin{array}{lll} 
& & 1 \\
& -1 & \\
1 & &
\end{array}\right)
$$

We have

$$
\bar{N}_{v}=\left\{\left.\left(\begin{array}{ccc}
1 & & \\
-\bar{x} & 1 & \\
y & x & 1
\end{array}\right) \right\rvert\, \begin{array}{c}
y+y+x \bar{x}=0 \\
x, y \in E_{v}
\end{array}\right\} .
$$

If $\bar{n} \in \bar{N}_{v}$, then by Iwasawa decomposition of $\operatorname{SU}_{3}\left(E_{v} / F_{v}\right)$, we get

$$
\bar{n}=\left(\begin{array}{rrr}
1 & & \\
-\bar{x} & 1 & \\
y & x & 1
\end{array}\right)=n\left(\begin{array}{lll}
\bar{a}^{-1} & & \\
& b & \\
& & a
\end{array}\right) k
$$

for some $n \in N_{v}, k \in K_{v}$.
As noted we can write $y=y_{1} \sqrt{c}-\frac{x \bar{x}}{2}$ for some $y_{1} \in F$.
Then $\operatorname{ord}_{E_{v}} y=\inf \left(\operatorname{ord}_{E_{v}} y_{1}, 2 \operatorname{ord}_{E_{v}} x\right)$ and

$$
|a|_{E_{v}}=|\tilde{\omega}|_{E_{v}}^{\inf \left(0, \operatorname{ord}_{E_{v}} x^{\prime}, \operatorname{ord}_{E_{v}} y\right)} .
$$

The zero in the "inf" is put into account for the case when both $x$ and $y$ are integral, and $\bar{n} \in K_{v}$.

Direct calculation using the definition of $\rho^{s}$ gives

$$
\rho^{s}\left(\left(\begin{array}{ccc}
a & & \\
& b & \\
& & \bar{a}^{-1}
\end{array}\right)\right)=|a|_{E_{v}}^{s}, \quad s \in C .
$$

To calculate the value of $\rho^{s+1}(\bar{n})$, we have to consider four cases:

1. $\operatorname{ord}_{E_{v}} x \geq 0$ and $\operatorname{ord}_{E} y_{1} \geq 0$
$\Rightarrow \operatorname{ord}_{E_{v}} y \geq 0$
$\Rightarrow \inf \left(0, \operatorname{ord}_{E_{v}} x, \operatorname{ord}_{E_{v}} y\right)=0$
$\Rightarrow \rho^{s+1}(\bar{n})=1$.
2. $2 \operatorname{ord}_{E_{v}} x \geq \operatorname{ord}_{E_{v}} y_{1}, \operatorname{ord}_{E_{v}} y_{1}<0, \operatorname{ord}_{E_{v}} y_{1}$ is even.
if $\operatorname{ord}_{E_{v}} x \geq 0$ then $\operatorname{ord}_{E_{v}} y_{1}<\operatorname{ord}_{E_{v}} x$.
If $\operatorname{ord}_{E_{v}} x<0$ then $\operatorname{ord}_{E_{v}} y_{1} \leq 2 \operatorname{ord}_{E_{v}} x<\operatorname{ord}_{E_{v}} x$.
Thus inf( $\left.0, \operatorname{ord}_{E_{v}} x, \operatorname{ord}_{E_{v}} y\right)=\operatorname{ord}_{E_{v}} y_{1}$ and $\rho^{s+1}(\bar{n})=\left|\bar{a}^{-1}\right|_{E_{v}}^{s+1}=q^{2(s+1) \text { ord } E_{v} y_{1}}$.

Note: if $\operatorname{ord}_{E_{v}} y_{1}=-2 m$ then

$$
\operatorname{ord}_{E_{v}} x \geq \frac{\text { ord } y_{1}}{2}=-m .
$$

3. $2 \operatorname{ord}_{E_{v}} x \geq \operatorname{ord}_{E_{v}} y_{1}<0, \operatorname{ord}_{E_{v}} y_{1}$ is odd
$\Rightarrow \inf \left(0, \operatorname{ord}_{E_{v}} x, \operatorname{ord}_{E} y\right)=\operatorname{ord}_{E_{v}} y_{2}$
$\Rightarrow \rho^{s+1}(\bar{n})=q^{2(s+1) \text { ord }_{E_{v}} y_{1}}$.

Note: if $\operatorname{ord}_{E_{v}} y_{1}=-(2 m-1), m \geq 1$ then

$$
\operatorname{ord}_{E_{v}} x \geq-m+\frac{1}{2} \text { or } \operatorname{ord}_{E_{v}} x \geq-(m-1)
$$

4. $2 \operatorname{ord}_{E_{v}} x<\operatorname{ord}_{E_{v}} y_{1}, \operatorname{ord}_{E_{v}} x<0$
$\Rightarrow \operatorname{ord}_{E_{v}} y=2 \operatorname{ord}_{E_{v}} x$
$\Rightarrow \rho^{s+1}(\bar{n})=q^{2(s+1) 2 \text { ord }_{E_{v}} x}$.
Note: if $\operatorname{ord}_{E_{v}} x=-m$ then $\operatorname{ord}_{E_{v}} y_{1}>-2 m \geq-(2 m-1)$.
Now we are ready to calculate the integral $M(s)$. We break the integral up into four pieces corresponding to the four cases above and transfer the integral over $\bar{N}\left(F_{v}\right)$ to those over $E_{v} \times F_{v}$, viz.,

$$
\begin{aligned}
M(s)= & \int_{\bar{N}\left(F_{v}\right)} \rho^{s+1}(\bar{n}) \mathrm{d} \bar{n}=\int_{\bar{N}_{1}} \int_{\bar{N}_{2}} \rho^{s+1}\left(\bar{n}_{1} \bar{n}_{2}\right) \mathrm{d} \bar{n}_{1} \mathrm{~d} \bar{n}_{2} \\
= & \int_{0_{E_{v}}} \int_{0_{F_{v}}} \mathrm{~d} x \mathrm{~d} y_{1}+\sum_{m=1}^{\infty} \int_{P_{\bar{E}_{v}}^{-2 m}-P_{F_{v}}^{-2 m}} \int_{P_{\bar{f}_{v}}^{-(2 m-1)}} q^{s(s+1)(-2 m)} \mathrm{d} x \\
& +\sum_{m=1}^{\infty} \int_{P_{E_{v}}^{-(m-1)}} \int_{P_{\bar{F}_{v}}^{(2 m-2)-} P_{\bar{F}_{v}}^{-(2 m-2)}} q^{-2(s+1)(2 m-1)} \mathrm{d} x \mathrm{~d} y_{1} \\
& +\sum_{m=1}^{\infty} \int_{P_{\bar{E}_{v}}^{-m}-P_{\bar{E}_{v}}^{(m-1)}} \int_{P_{E_{v}}^{(2 m-1)}} q^{2(s+1)(-2 m)} \mathrm{d} x \mathrm{~d} y_{1}
\end{aligned}
$$

where $P_{E_{v}}$ (resp. $P_{E_{v}}$ ) is the maximal prime ideal of $E_{v}$ (resp. $F_{v}$ ). We normalized measure on $E_{v}, F_{v}$ by $\int_{0_{E_{v}}} \mathrm{~d} x=1$ and $\int_{0_{F_{v}}} \mathrm{~d} y_{1}=1$.

Further calculation gives

$$
\begin{gathered}
\int_{0_{E_{v}}} \int_{0_{F_{v}}} \mathrm{~d} x \mathrm{~d} y_{1}=1 \\
\sum_{m=1}^{\infty} \int_{P_{\bar{F}_{v}^{m}}^{m}-P_{F_{v}}^{(2 m-1)}} q^{2(s+1)(-2 m)} \mathrm{d} x \mathrm{~d} y_{1} \\
=\sum_{m=1}^{\infty} q^{2 m}\left(q^{2 m}-q^{2 m-1}\right) q^{-4(s+1) m} \\
=\left(1-q^{-1}\right) \sum_{m=1}^{\infty}\left(q^{-4 s}\right)^{m}=\frac{\left(1-q^{-1}\right) q^{-4 s}}{1-q^{-4 s}}
\end{gathered}
$$

$$
\begin{aligned}
& \sum_{m=1}^{\infty} \int_{P_{\bar{E}_{v}}^{(m-1)}} \int_{P_{F_{v}}^{(2 m-1)-P \bar{F}_{v}^{-(2 m-2)}}} q^{-2(s+1)(2 m-1)} \mathrm{d} x \mathrm{~d} y_{1} \\
& =\sum_{m=1}^{\infty} q^{2 m-2}\left(q^{2 m-1}-q^{2 m-2}\right) q^{-2(s+1)(2 m-1)}, \\
& =\left(q^{-1}-q^{-2}\right) q^{2 s} \sum_{m=1}^{\infty}\left(q^{-4 s}\right)^{m}, \\
& =\frac{\left(q^{-1}-q^{-2}\right) q^{-2 s}}{1-q^{-4 s}} \\
& \sum_{m=1}^{\infty} \int_{P_{E_{v}^{-m}-P_{E_{v}}^{-(m-1)}} \int_{P_{F_{v}}^{-(2 m-1)}} q^{2(s+1)(-2 m)} \mathrm{d} x \mathrm{~d} y_{1}}^{=\sum_{m=1}^{\infty}\left(q^{2 m}-q^{2 m-2}\right) q^{2 m-1} q^{-4 m(s+1)},} \\
& =\left(1-q^{-2}\right) q^{-1} \sum_{m=1}^{\infty}\left(q^{-4 s}\right)^{m}, \\
& =\frac{\left(q^{-1}-q^{-3}\right) q^{-4 s}}{1-q^{-4 s}} .
\end{aligned}
$$

Adding all the terms, we have

$$
M(s)=\frac{\left(1-q^{-2 s-2}\right)\left(1+q^{-2 s-1}\right)}{\left(1-q^{-2 s}\right)\left(1+q^{-2 s}\right)}
$$

To complete the proof of the proposition, let us look at the Lie algebra $\hat{\mathfrak{g}}$ of the analytic group $\hat{G}$ associated with $G$. We can take $\hat{\mathfrak{g}}$ to be $\mathfrak{g l} l_{2}(\mathbb{C})$ and let $\hat{\Sigma}^{+}=\left\{\hat{\alpha}_{1}, \hat{\alpha}_{2}, \hat{\alpha}_{3}\right\}, \hat{\alpha}_{3}=\hat{\alpha}_{1}+\hat{\alpha}_{2}$. There exists root vectors $X_{\hat{\alpha}_{1}}, X_{\hat{\alpha}_{2}}, X_{\hat{\alpha}_{3}}$ such that

$$
\left[\boldsymbol{X}_{\hat{\alpha}_{1}}, \boldsymbol{X}_{\hat{\alpha}_{2}}\right]=\boldsymbol{X}_{\hat{\alpha}_{2}} .
$$

$\hat{\mathfrak{g}}$ has a Dynkin diagram of type $A_{2}$

the arrows indicate the action of $\sigma \in \operatorname{Gal}(E / F)$, i.e. $\sigma\left(X_{\hat{\alpha}_{1}}\right)=X_{\hat{\alpha}_{2}}$. Since this action is to be extended to a Lie algebra isomorphism, i.e. $\sigma\left[X_{\hat{\alpha}_{1}}, X_{\hat{\alpha}_{2}}\right]=\left[\sigma X_{\hat{\alpha}_{1}}, \sigma X_{\hat{\alpha}_{2}}\right]$, so $\sigma X_{\hat{\alpha}_{3}}=\left[X_{\hat{\alpha}_{2}}, X_{\hat{\alpha}_{1}}\right]=-X_{\hat{\alpha}_{3}}$.

Also, we have
or

$$
\begin{aligned}
(\operatorname{Ad} \hat{t}) X_{\hat{\alpha}} & =\hat{\alpha}(\hat{t}) X_{\hat{\alpha}}=|\tilde{\omega}|_{F_{v}}^{s \rho, \hat{\alpha}\rangle} X_{\hat{\alpha}} \\
& =|\tilde{\omega}|_{F_{v}}^{s} X_{\hat{\alpha}} \quad \text { if } \hat{\alpha}=\hat{\alpha}_{1} \text { or } \hat{\alpha}_{2}, \\
& =|\tilde{\omega}|_{F_{v}}^{s} X_{\hat{\alpha}} \quad \text { if } \hat{\alpha}=\hat{\alpha}_{3},
\end{aligned}
$$

because $\langle\rho, \hat{\alpha}\rangle=\frac{2(\rho, \alpha)}{(\alpha, \alpha)}=1$ if $\alpha$ simple and

$$
\left\langle\rho, \hat{\alpha}_{3}\right\rangle=\left\langle\rho, \alpha_{1}\right\rangle+\left\langle\rho, \alpha_{2}\right\rangle=2 .
$$

We take $\hat{\mathfrak{n}}=\mathbf{C} X_{\hat{\alpha}_{1}}+\mathbf{C} X_{\hat{\alpha}_{2}}+\mathbf{C} \boldsymbol{X}_{\hat{\alpha}_{3}}$. Then

$$
\begin{aligned}
& \operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\tilde{n}}\right) \\
& =\operatorname{det}\left(I-\left(\begin{array}{ccc}
0 & |\tilde{\omega}|_{F_{v}}^{s} & 0 \\
|\tilde{\omega}|_{F_{v}}^{s} & 0 & 0 \\
0 & 0 & -|\tilde{\omega}|_{F_{v}}^{2 s}
\end{array}\right)\right), \\
& =\left(1-q^{-2 s}\right)\left(1+q^{-2 s}\right)
\end{aligned}
$$

and

$$
\begin{aligned}
& \operatorname{det}\left(I-\left.|\tilde{\omega}|_{F_{v}} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right) \\
& =\left(1-q^{-2 s-s}\right)\left(1+q^{-2 s-1}\right) .
\end{aligned}
$$

This completes the proof of the proposition.
3.5. Let us now consider the case (III). $G$ is a connected semisimple quasi-split algebraic group defined over $F_{V}$ splits over an unramified extension $E_{v} / F_{v}$ of degree $n$.

The absolute Dynkin diagram of $G$ consists of $n$ copies of $A_{1}$, and the action of the Frobenius $\sigma$ in $\operatorname{Gal}\left(E_{v} / F_{v}\right)$ is the cyclic permutation as indicated


The action has only one orbit; $G$ is of $F$-rank 1 and $G\left(F_{v}\right)=\mathrm{SL}_{2}\left(E_{v}\right)$. The integral that we are interested in becomes $M(s)=\int_{\bar{N}_{v}} \rho^{s+1}(\bar{n}) \mathrm{d} \bar{n}$
where

$$
\begin{aligned}
\bar{N}_{v} & =\left\{\left.\left(\begin{array}{ll}
1 & 0 \\
x & 1
\end{array}\right) \right\rvert\, x \in E_{v}\right\}, \\
A_{v} & =\left\{\left.\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) \right\rvert\, a \in E_{v}^{x}\right\},
\end{aligned}
$$

and

$$
|a|_{E_{v}}^{s+1}=\rho^{s+1}\left(\left(\begin{array}{ll}
a & \\
& a^{-1}
\end{array}\right)\right) .
$$

So by §3.3

$$
\begin{aligned}
M(s) & =\frac{1-q_{E_{v}}^{-(s+1)}}{1-q_{E_{v}}^{-s}} \\
& =\frac{1-q^{-n(s+1)}}{1-q^{-n s}}
\end{aligned}
$$

 corresponding to the positive root $\hat{\alpha}_{i}$ of the $\boldsymbol{i}^{\text {th }}$ copy of $\mathfrak{b l}_{2}$ in the product. Then

$$
(\operatorname{Ad} \hat{t}) X_{\hat{\alpha}_{i}}=\hat{\alpha}_{i}(\hat{t}) X_{\hat{\alpha}_{i}}=\left.|\tilde{\omega}|\right|_{F_{v}} ^{s\left(\rho, \hat{c}_{i}\right\rangle} X_{\hat{\alpha}_{i}}=q^{-s} X_{\hat{\alpha}_{i}}
$$

because $\rho=\frac{1}{2} \Sigma \alpha_{i}$ and as the diagram is disconnected $\left\langle\alpha_{j}, \hat{\alpha}_{i}\right\rangle=0$ if $i \neq j$, and $\left\langle\frac{\alpha_{i}}{2}, \hat{\alpha}_{i}\right\rangle=1$. So,

$$
\begin{aligned}
& \operatorname{det}\left(I-\left.\operatorname{Ad} \hat{t}\right|_{\hat{n}}\right)=\left[\begin{array}{ccccc}
1 & & & & \\
-q^{-s} & 1 \cdot \ddots & & & \\
& -q^{-s} \cdot \ddots & \ddots & \\
& & \ddots \ddots & \\
& & & -q^{-s} & 1
\end{array}\right], \\
&=1-q^{-n s} .
\end{aligned}
$$

Similarly,

$$
\operatorname{det}\left(I-\left.|\tilde{\omega}|_{F_{v}} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right)=1-q^{-n(s+1)},
$$

and we are done.
3.6. Finally, let us look at the last case IV. Here $G$ is a $F_{v}$-form of a split group with a Dynkin diagram consisting of $n$ copies of $A_{2} . G$ is defined over $F_{v}$ splits over an unramified extension $E_{v}$ of degree $2 n$; there exists a field $E_{v}^{\prime}$ in $E_{v} / F_{v}$ such that $\left[E_{v}^{\prime}: F_{v}\right]=n$; the non-trivial element of $\operatorname{Gal}\left(E_{v} / E_{v}^{\prime}\right)\left(\subset \operatorname{Gal}\left(E_{v} / F_{v}\right)\right)$ give rise to the twisting; the action of this element is shown in the diagram


This determines a special unitary group $\mathrm{SU}_{3}\left(E_{v} / E_{v}^{\prime}\right)$ with respect to the form

$$
J=\left(\begin{array}{lll} 
& & 1 \\
& 1 & \\
1 & &
\end{array}\right)
$$

such that

$$
G(F) \approx \mathrm{SU}_{3}\left(E_{v} / E_{v}^{\prime}\right)=\left\{\left.g \in \mathrm{SL}_{3}\left(E_{v}\right)\right|^{t} \bar{g} g=J\right\}
$$

Thus, using the result in $\S 3.4$, we get

$$
M(s)=\frac{\left(1-q^{-2 n(s+1)}\right)\left(1+q^{-n(2 s+1)}\right)}{\left(1-q^{-2 n s}\right)\left(1+q^{-2 n s}\right)}
$$

(Note: modulus of $E_{v}=q^{2 n}$.)
To establish the formula

$$
M(s)=\frac{\operatorname{det}\left(I-\left.|\tilde{\omega}|_{F_{v}} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right)}{\operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right)}
$$

we shall evaluate the determinants directly.
Let us denote the simple root system $\Delta$ by $\left\{\alpha_{1}, \beta_{1} ; \ldots ; \alpha_{n}, \beta_{n}\right\}$. We calculate

$$
\begin{aligned}
(\operatorname{Ad} \hat{t}) X_{\hat{\alpha}_{i}} & =\hat{\alpha}_{i}(\hat{t}) X_{\hat{\alpha}_{i}}=|\tilde{\omega}|_{F_{b}}^{s\left(\rho, \hat{\alpha}_{i}\right\rangle} X_{\hat{\alpha}_{i}} \\
& =q^{-s} X_{\hat{\alpha}_{i}} .
\end{aligned}
$$

Here $\rho=\frac{1}{2} \sum_{i=1}^{n}\left(\alpha_{i}+\beta_{i}+\left(\alpha_{i}+\beta_{i}\right)\right)$,

$$
\left\langle\rho, \hat{\alpha}_{i}\right\rangle=\sum_{j=1}^{n}\left\langle\rho_{j}, \hat{\alpha}_{i}\right\rangle \quad \text { where } \rho_{j}=\alpha_{j}+\beta_{j}
$$

because $\boldsymbol{i} \neq \mathbf{j}$

$$
\left\langle\rho_{j}, \hat{\alpha}_{i}\right\rangle=0,
$$

and

$$
\left\langle\rho_{i}, \hat{\alpha}_{i}\right\rangle=1
$$

Similarly

$$
(\operatorname{Ad} \hat{t}) X_{\hat{\beta}_{i}}=q^{-s} X_{\hat{\beta}_{i}},
$$

and

$$
(\operatorname{Ad} \hat{t}) X_{\hat{\alpha}_{i}+\hat{\beta}_{i}}=q^{-2 s} X_{\hat{\alpha}_{i}+\hat{\beta}_{i}} .
$$

Next we write down the effect of the Galois action as indicated by the arrows in the above diagram. For $1 \leq i \leq n-1$,

$$
\begin{gathered}
\boldsymbol{\sigma} \boldsymbol{X}_{\hat{\alpha}_{i}}=\boldsymbol{X}_{\hat{\alpha}_{i+1},}, \\
\boldsymbol{\sigma} \boldsymbol{X}_{\hat{\beta}_{i}}=\boldsymbol{X}_{\hat{\beta}_{i+1},}, \\
\boldsymbol{\sigma} \boldsymbol{X}_{\hat{\alpha}_{i}+\hat{\beta}_{i}}=\boldsymbol{\sigma}\left[\boldsymbol{X}_{\hat{\alpha}_{i}}, X_{\hat{\beta}_{i}}\right]=\left[\boldsymbol{\sigma} \boldsymbol{X}_{\hat{\alpha}_{i}}, \boldsymbol{\sigma} \boldsymbol{X}_{\hat{\beta}_{i}}\right] \\
=\left[X_{\hat{\alpha}_{i+1}, 1}, X_{\hat{\beta}_{i+1}}\right]=\boldsymbol{X}_{\hat{\alpha}_{i+1}+\hat{\beta}_{i+1}},
\end{gathered}
$$

and

$$
\begin{gathered}
\boldsymbol{\sigma} \boldsymbol{X}_{\hat{\alpha}_{n}}=\boldsymbol{X}_{\hat{\beta}_{1}}, \\
\boldsymbol{\sigma} \boldsymbol{X}_{\hat{\beta}_{n}}=\boldsymbol{X}_{\hat{\alpha}}, \\
\boldsymbol{\sigma} \boldsymbol{X}_{\hat{\alpha}_{n}+\hat{\beta}_{n}}=\left[\boldsymbol{\sigma} \boldsymbol{X}_{\hat{\alpha}_{n}}, \boldsymbol{\sigma} \boldsymbol{X}_{\hat{\boldsymbol{\beta}}_{n}}\right]=\left[\boldsymbol{X}_{\hat{\boldsymbol{\beta}}_{1}}, X_{\hat{\alpha}_{1}}\right]=-\boldsymbol{X}_{\hat{\alpha}_{1}+\hat{\beta}_{1}} .
\end{gathered}
$$

If we take the basis of $\mathfrak{n}$ to be $X_{\hat{\alpha}_{1}}, X_{\hat{\beta}_{1}}, X_{\hat{\alpha}_{1}+\hat{\beta}_{1}}, \ldots, X_{\hat{\alpha}_{n}}, X_{\hat{\beta}_{n}}, X_{\hat{\alpha}_{n}+\hat{\beta}_{n}}$ (in that order), then it is trivial to show that

$$
\begin{aligned}
& \operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right) \\
& =\left(1-q^{-2 n s}\right)\left(1+q^{-2 n s}\right),
\end{aligned}
$$

and

$$
\begin{gathered}
\operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}}\right) \\
=\left(1-q^{-2 n(s+1)}\right)\left(1+q^{-n(2 s+1)}\right) .
\end{gathered}
$$

Thus the required formula is proved. With this we complete the proof of Proposition 3.2.

## 4. Reduction to rank one

To determine the local factor $M_{v}(w, \lambda)$ for almost all $v$ for $G$ of arbitrary $F$-rank, we use the method of reduction to $F$-rank one which was first studied by Bhanu-Murti [1] and was extended by Gindikin and Karpelevich [6]. This method has also been used in Langlands' Euler Product (Yale, 1971) and in the thesis of Jacquet (Paris) and Lai (Yale). Here we shall follow Shiffmann [19].
4.1. We want to calculate the integral (9) of $\S 2$. For $\lambda \in L_{F} \otimes \mathbb{C}$, $\mathscr{E}(\lambda) \neq 0$ and so $\mathscr{E}_{v}(\lambda) \neq 0$ for all $v$. We have $W_{F} \subset W_{F}$. We can consider $w$ as an element of $W_{F_{v}}$ and do the rest of the calculation over $F_{V}$. Moreover for almost all $v, \mathscr{E}_{v}(\lambda)$ is one dimensional. It is sufficient to evaluate the integral for the following function in $\mathscr{E}_{v}(\lambda)$ :

$$
\begin{equation*}
\Phi\left(g_{v}\right)=\left|\lambda\left(a_{v}\right) \rho\left(a_{v}\right)\right|_{v} \tag{1}
\end{equation*}
$$

where $g_{v}=n_{v} a_{v} k_{v} \in G_{v}$. The linear transformation $M_{v}(w, \lambda)$ is just multiplication by the following constant which we also denoted by $M_{v}(w, \lambda)$ :

$$
M_{v}(w, \lambda)=\int_{N_{v}^{w}} \Phi(w n) \mathrm{d} n .
$$

Changing the variable by $n \rightarrow w^{-1} n w$ and writing $\bar{N}^{w}=w N^{w} w^{-1}=$ $w N w^{-1} \cap N$, we have

$$
\begin{equation*}
M_{v}(w, \lambda)=\int_{\bar{N}_{v}^{w}} \Phi(n w) \mathrm{d} n . \tag{2}
\end{equation*}
$$

Recall that the length $\ell(w)$ of $w$ is the smallest integer $g$ of such that there exists $g$ simple $F_{v}$-roots $\beta_{1}, \ldots, \beta_{g}$ with

$$
\begin{equation*}
w=s_{\beta_{1}}, \ldots, s_{\beta_{8}} \tag{3}
\end{equation*}
$$

( $s_{\alpha_{j}}$ is the symmetry with respect to $\alpha_{j}$ ). Moreover the $F_{v}$-roots $\alpha_{j}=s_{\beta_{\ell(w)}} \ldots s_{\beta_{j+1}}\left(\beta_{j}\right) j=1, \ldots, \ell(w)$ are positive and if we write

$$
{ }_{0} \Sigma_{F_{v}}^{+}(w)=\left\{\left.\alpha \in \in_{0} \Sigma_{F_{v}}^{+}\right|^{w} \alpha<0\right\}
$$

then

$$
{ }_{0} \Sigma_{F_{v}}^{+}(w)=\left\{\alpha_{1}, \ldots, \alpha_{\ell(w)}\right\} .
$$

We quote the following lemma from Schiffmann ([19], Prop. 1.3).
4.2. Lemma: Let $w, w^{\prime}, w^{\prime \prime}$ be three elements of $w_{F}$ such that $w=$ $w^{\prime} w^{\prime \prime}$ with $\ell(w)=\ell\left(w^{\prime}\right)+\ell\left(w^{\prime \prime}\right)$. Then the map (4) $\left(n^{\prime}, n^{\prime \prime}\right) \rightarrow$ $n^{\prime}\left(w^{\prime} n^{\prime \prime} w^{\prime-1}\right)$ defines a variety isomorphism $\bar{N}^{w^{\prime}} \times \bar{N}^{w^{\prime \prime}} \rightarrow \bar{N}^{w}$.
4.3. Using the above lemma, and assuming the integrals involve converges, we have

$$
\begin{aligned}
M_{v}(w, \lambda) & =\int_{\bar{N}_{v}^{w^{\prime}} \times \bar{N}_{v}^{w^{\prime \prime}}} \Phi\left(n^{\prime} w^{\prime} n^{\prime \prime} w^{\prime-1} w\right) \mathrm{d} n^{\prime} \mathrm{d} n^{\prime \prime}, \\
& =\int_{\bar{N}_{v}^{w^{\prime \prime}}} M_{v}\left(w^{\prime}, \lambda\right) \Phi\left(n^{\prime \prime} w^{\prime \prime}\right) \mathrm{d} n^{\prime \prime},
\end{aligned}
$$

and so

$$
\begin{equation*}
M_{v}(w, \lambda)=M_{v}\left(w^{\prime}, \lambda^{w^{\prime \prime}}\right) M_{v}\left(w^{\prime \prime}, \lambda\right) . \tag{5}
\end{equation*}
$$

If we write $w$ as a product of symmetries (as in (3)) then formula (5) allows us to reduce the calculation to the case $\ell(w)=1$, i.e. the $F$-rank one case, and in this case the convergence follows from the explicit formula given in §3. To summarize we have
4.4. Proposition: Let $N_{\alpha}=G_{\alpha} \cap N$ for $\alpha \in{ }_{0} \Sigma_{F}^{+}$and $\bar{N}_{\alpha}$ the unipotent subgroup of $G_{\alpha}$ opposite to $N_{\alpha}$. Then the integral (2) converges for $\lambda \in L_{F} \otimes C$ with $\operatorname{Re}(\langle\lambda, \hat{\alpha}\rangle)>0$ for all $\alpha \in{ }_{0} \Sigma_{F}^{+}(w)$,

$$
\begin{equation*}
M_{v}(w, \lambda)=\prod_{\alpha \in E_{0} \Sigma_{F}^{+}(w)} \int_{\bar{N}_{\alpha}\left(F_{v}\right)} \Phi_{\alpha}(\bar{n}) \mathrm{d} \bar{n} \tag{6}
\end{equation*}
$$

where $\Phi_{\alpha}$ is the restriction of $\Phi$ to $G_{\alpha}$.
4.5. As each $G_{\alpha}$ has $F_{v}$-rank one we can apply Proposition 3.2 to get

$$
\begin{equation*}
\int_{\bar{N}_{\alpha}\left(F_{v}\right)} \Phi_{\alpha}(\bar{n}) \mathrm{d} \bar{n}=\frac{\operatorname{det}\left(I-\left.|\tilde{\omega}|_{v} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}_{\alpha}}\right)}{\operatorname{det}\left(I-\sigma \operatorname{Ad~} \hat{t} \mid \hat{n}_{\alpha}\right)} \tag{7}
\end{equation*}
$$

Let $\hat{\mathfrak{n}}$ be the nilpotent subalgebra of $\hat{\mathfrak{g}}$ spanned by $\hat{\mathfrak{g}}_{\alpha}$ for $\alpha \in$ ${ }_{0} \Sigma_{F_{v}}^{+}(w)$. The action of $\sigma \operatorname{Ad} \hat{t}$ on $\hat{\mathfrak{n}}^{w}$ preserves the subspaces $\hat{\mathrm{n}}_{\alpha}$. Hence

$$
\begin{equation*}
\frac{\operatorname{det}\left(I-\left.|\tilde{\omega}|_{v} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}^{w}}\right)}{\operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}^{w}}\right)}=\prod_{\alpha \in_{0} \Sigma_{F_{v}}(w)} \frac{\operatorname{det}\left(I-\left.|\tilde{\omega}|_{v} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{a}_{\alpha}}\right)}{\operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\dot{n}_{\alpha}}\right)} \tag{8}
\end{equation*}
$$

The following proposition follows immediately from (6), (7) and (8).
4.6. Proposition: For almost all v, we have

$$
\begin{equation*}
M_{v}(w, \lambda)=\frac{\operatorname{det}\left(I-\left.|\tilde{\omega}|_{v} \sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}^{w}}\right)}{\operatorname{det}\left(I-\left.\sigma \operatorname{Ad} \hat{t}\right|_{\hat{n}^{w}}\right)} \tag{9}
\end{equation*}
$$

where $\sigma$ is the Frobenius and $\hat{t}=\hat{t}_{\lambda}$.

## 5. Value of the local factor at one

5.1. Let $\mathscr{S}$ be a finite set of places of $F$ containing all the infinite place of $F$, all the ramified places of $F$ and all the places at which the conditions (i) to (iii) of $\S 3.1$ are not satisfied. Let us write

$$
\boldsymbol{M}_{\mathscr{Y}}(s)=\prod_{v \in \mathscr{Y}} \boldsymbol{M}_{v}\left(w_{0}, \rho^{s}\right)
$$

where $s \in \mathbb{C}$ and $w_{0} \in W_{F}$ sends all positive roots to negative roots. Then $M_{\varphi}(1)$ can be considered as a linear map $E_{\varphi}(\rho) \rightarrow E_{\varphi}\left(\rho^{-1}\right)$ and

$$
\begin{equation*}
M_{\varphi}(1) \Phi(g)=\int_{N_{\mathscr{Y}}} \Phi\left(w_{0} n g\right) \mathrm{d} n \tag{1}
\end{equation*}
$$

for $\Phi \in \mathscr{E}_{\mathscr{y}}(\rho), g \in G_{\varphi}$. Now $G_{\mathscr{y}}=B_{\mathscr{}} K_{\mathscr{\varphi}}$ implies that $\mathscr{C}_{\varphi}(\rho)$ is one dimensional and $M_{\mathscr{G}}(1)$ is just multiplication by a constant which we also
denoted by $M_{\mathscr{C}}(1)$. We have

$$
\begin{equation*}
M_{\varphi}(1)=\int_{N_{\varphi}} \rho^{2}\left(w_{0} n\right) \mathrm{d} n . \tag{2}
\end{equation*}
$$

5.2. Let $L(s, G)$ be the Artin $L$-function of the Galois action on the rational characters of $G, L_{v}(s, G)$ be the local factor at $v$ of $L(s, G)$ and

$$
\mu_{G}=\lim _{s \rightarrow 1}(s-1)^{r_{G}} L(s, G)
$$

where $r_{G}$ is the rank of $X(G)_{F}$. Similar definitions are made with $A$ replacing $G$.

Proposition: For $\mathscr{S}$ sufficiently large we have

$$
\begin{equation*}
M_{\mathscr{C}}(1)=\kappa \frac{\mu_{G}}{\mu_{A}} \prod_{v \in \mathscr{Y}} \frac{L_{v}(1, A)}{L_{v}(1, G)} \prod_{v \notin \mathscr{Y}} \operatorname{vol} K_{v} \tag{3}
\end{equation*}
$$

where the vol $K_{v}$ is calculated by the local measure $d g_{v}$.
Proof: Let $h$ be an integrable function on $N_{\mathscr{\varphi}}+A_{\varphi \cdot}$. Let $f$ be a function on $G(\mathbb{A})$ which vanishes at $g$ except if $g_{v} \in K_{v}$ for all $v \notin \mathscr{S}$ and if the latter condition is satisfied, we have

$$
f(g)=f\left(g_{\varphi}\right)=h(n, a)
$$

for $g=n a k$. First of all we have

$$
\begin{equation*}
\int_{G(A)} f(g) \mathrm{d} g=\kappa \int_{N_{g} \times A_{g}} h\left(n_{2}, a_{2}\right) \rho^{-2}\left(a_{2}\right) \mathrm{d} n_{2} \mathrm{~d} a_{2} . \tag{4}
\end{equation*}
$$

On the other hand, suppose that $g_{g}$ lies in the large cell $N_{\varphi} S_{\varphi} w_{0} N_{\varphi}$ of the Bruhat decomposition: $g_{\varphi}=n_{2} a_{2} w_{0} n_{1}$ where $a_{2} \in A_{\varphi}$ and $n_{1}, n_{2} \in N_{\mathscr{C}}$ and if we write $w_{0} n_{1}=n\left(n_{1}\right) a\left(n_{1}\right) k$ with $n\left(n_{1}\right) \in N_{\mathscr{C}}$ and $a\left(n_{1}\right) \in A_{\varphi}$, then $g_{\varphi}=n_{2} a_{2} n\left(n_{1}\right) a_{2}^{-1} a_{2} a\left(n_{1}\right) k$ and

$$
\begin{gather*}
\int_{G(A)} f(g) \mathrm{dg}  \tag{5}\\
=\prod_{v \notin \mathscr{Y}} \operatorname{vol}\left(K_{v}\right) \int_{N_{g} A_{g} N_{\mathscr{C}}} h\left(n_{2} a_{2} n\left(n_{1}\right) a_{2}^{-1}, a_{2} a\left(n_{1}\right)\right) \rho^{-2}\left(a_{2}\right) \mathrm{d} n_{2} \overline{\mathrm{~d}}_{2} \mathrm{~d} n_{1} .
\end{gather*}
$$

After changing the measures, the integral in the above formula becomes

$$
\int_{N_{\mathscr{G}} A_{\mathscr{H}} \mathrm{N}_{\mathscr{Y}}} \rho^{2}\left(a\left(n_{1}\right)\right) h\left(n_{2}, a_{2}\right) \rho^{-2}\left(a_{2}\right) \mathrm{d} n_{2} \overline{\mathrm{~d}}_{1} \mathrm{~d} n_{1} .
$$

Substitute this and

$$
\mathrm{d} a_{2}=\left(\prod_{v \in \mathscr{Y}} L_{v}(1, A)\right) \overline{\mathrm{d} a}_{2}
$$

into (5). Comparing the result with (4), we obtain (3) by noting that the choice of $h$ is arbitrary.
5.3. Corollary: For $v \notin \mathscr{S}$, if we write

$$
M_{v}(1)=M_{v}\left(w_{0}, \rho\right)=\int_{N_{v}} \rho^{2}\left(w_{0} n\right) \mathrm{d} n
$$

then

$$
\begin{equation*}
M_{v}(1)=\operatorname{vol}\left(K_{v}\right) \cdot L_{v}(1, A) / L_{v}(1, G) . \tag{6}
\end{equation*}
$$

Proof: Apply the proposition to $\mathscr{S}^{\prime}=\mathscr{S} \cup\{v\}$. The corollary then follows immediate form

$$
M_{Y^{\prime}}(1)=M_{v}(1) M_{\mathscr{C}}(1) .
$$

5.4. Remark: We have followed Rapoport [18] in the proof of corollary 5.3. An alternative approach is given in my thesis (Yale 1974) in which (6) is deduced from (9) of $\S 4$ by calculating directly $\operatorname{vol}\left(K_{v}\right)$ via reduction $\bmod v$.

## 6. The constant functions

We calculate in this section the projection of $\mathscr{E}$ into the subspace of constant functions in $\left.\mathscr{L}^{2}\left(Z_{\infty}^{+} G(F)\right) \backslash G(\mathbb{A})\right)$.
6.1. Let $\mathscr{L}$ be the closed subspace of $\mathscr{L}^{2}\left(Z_{\infty}^{+} G(F)\right) \backslash G(\mathbb{A})$ generated by $\tilde{\phi}$ for $\phi \in \mathscr{D}$. Write $\mathscr{H}$ for the union of $\mathscr{H}(\lambda)$ for all $\lambda$ in $L_{F} \otimes C$. Suppose that $f$ is a complex valued function defined, bounded and
analytic in a tube in $L_{F} \otimes C$ over a ball of radius $R$ with centre at zero and $R>(\rho, \rho)^{1 / 2}$. Assume also that $f\left({ }^{( } \lambda\right)=f(\lambda)$ for all $w \in W_{F}$. Then

$$
\Phi \rightarrow \Psi=f \Phi
$$

where $\Psi(\lambda, g)=f(\lambda) \Phi(\lambda, g)$, defines a linear map on $\mathscr{H}$ and induces a bounded linear operator

$$
\Lambda(f): \tilde{\phi} \rightarrow \tilde{\psi}
$$

on $\mathscr{L}$. If $a>(\rho, \rho)$ and $f(\lambda)=(a-(\lambda, \lambda))^{-1}$, then $\Lambda(f)$ is self-adjoint. We define

$$
\mathscr{A}=a-\Lambda(f)^{-1} .
$$

It is an unbounded self-adjoint operator on $\mathscr{L}(\mathscr{A}$ is introduced in Langlands [14] §6 and [15]). It is obvious that if $\Psi(\lambda, g)=$ $(\lambda, \lambda) \Phi(\lambda, g)$ then $\mathscr{A} \bar{\phi}=\tilde{\psi}$. The following two lemmas and the corollary are easy to prove.
6.2. Lemma: Let (,) be the inner product on $\mathscr{L}^{2}\left(Z_{x}^{+} G(F)\right) \backslash G(A)$ and 1 be the constant function. For $\tilde{\phi} \in \mathscr{L}$, we have

$$
\begin{equation*}
(\tilde{\phi}, 1)=\kappa \Phi(\rho, 1) . \tag{1}
\end{equation*}
$$

6.3. Lemma: For $\tilde{\phi} \in \mathscr{L}$ and $\mathscr{A}$ as defined above we have

$$
\begin{equation*}
(\mathscr{A} \tilde{\phi}, 1)=(\rho, \rho)(\tilde{\phi}, 1) . \tag{2}
\end{equation*}
$$

6.4. Corollary: $\mathscr{A} 1=(\rho, \rho) 1$.
6.5. For $z \in \mathbf{C}$, let $R(z, \mathscr{A})=(z-\mathscr{A})^{-1}$ be the resolvent of $\mathscr{A}$. For $\lambda_{0} \in L_{F} \otimes R$ if $\operatorname{Re} z>\left(\lambda_{0}, \lambda_{0}\right)$, then it is easy to show that

$$
\begin{equation*}
(R(z, \mathscr{A}) \tilde{\phi}, \tilde{\psi})=\kappa \sum_{w \in w_{F}} \int_{|\lambda|=\lambda_{0}} \frac{M(w, \lambda) \Phi(\lambda) \bar{\Psi}\left(-{ }^{w} \bar{\lambda}\right)}{z-(\lambda, \lambda)} \mathrm{d} \lambda . \tag{3}
\end{equation*}
$$

Let $E(x),-\infty<x<\infty$ be a right continuous spectral resolution of the self-adjoint operator $\mathscr{A}$. It is obvious that $(\rho, \rho)$ belongs to the point spectrum of $\mathscr{A}$ and corollary 6.4 implies that the constant functions are in the range of the projection $E((\rho, \rho))-E((\rho, \rho)-0)=$ $E$ (say). Suppose $a>(\rho, \rho)>b$, and $a-b$ is small, then $(E \tilde{\phi}, \tilde{\psi})$ is
given by Stieljes inversion,

$$
\begin{align*}
\frac{1}{2}\{(E(a) \tilde{\phi}, \tilde{\psi})+ & (E(a-0) \tilde{\phi}, \tilde{\psi})\}-\frac{1}{2}\{(E(b) \tilde{\phi}, \tilde{\psi})+(E(b-0) \tilde{\phi}, \tilde{\psi})\}  \tag{4}\\
= & \lim _{\epsilon \rightarrow 0^{+}} \frac{1}{2 \pi i} \int_{C(a, b, c, \epsilon)}(R(z, \mathscr{A}) \tilde{\phi}, \tilde{\psi}) \mathrm{d} z
\end{align*}
$$

where $C(a, b, c, \epsilon)$ is the following contour:

6.6. Next we want to determine the dual measure for the Fourier transform on $A$.

We have put on $A(\mathbb{A})$ the Tamagawa measure $\mathrm{d} a$ which can be written as $\mathrm{d} a=\mathrm{d} a^{1} \mathrm{~d}$ corresponding to the decomposition $A(\mathbb{A})=$ $A^{1}(\mathbb{A}) A_{\infty}^{+}$. In $\S 2.3$ we put a measure on $\left(Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})\right)^{*}$ via Pontryagin duality. But

$$
\left(Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})\right)^{*}=\left(A(F) \backslash A(\mathbb{A})^{1}\right)^{*} \times \operatorname{Hom}\left(Z_{\infty}^{+} \backslash A_{\infty}^{+}, \mathbb{C}^{*}\right)
$$

and $\left(A(F) \backslash A(\mathbb{A})^{1}\right)^{*}$ is discrete, $\operatorname{Hom}\left(Z_{\infty}^{+} \backslash A_{\infty}^{+}, C^{*}\right)$ is a vector space over C. Thus we can give $\left(Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})\right)^{*}$ the structure of a complex manifold; as such, it has a natural measure which gives the measure 1 to the identity element of the Pontryagin dual of the compact abelian group $A(F) \backslash A^{1}(\mathbb{A})$; while the dual measure to $\mathrm{d} a^{1}$ gives the measure $1 / \operatorname{vol}\left(A(F) \backslash A^{1}(\mathbb{A})\right)$ to the identity element.

The measure on $A_{\infty}^{+}\left(\right.$resp. $\left.A_{\infty}^{+}, Z_{\infty}^{+}\right)$is fixed by identifying it with a power of $\mathbf{R}^{\times}$by means of a basis of the lattice $L_{F}$ (resp. ${ }^{1} L_{F},{ }^{0} L_{F}^{+}$). Since $A_{\infty}^{+}=Z_{\infty}^{+} A_{\infty}^{\prime+}$, we see that the dual measure to da gives the measure $1 / f$ to the identity element of ${ }^{1} L_{F}$, where $f=$ $\left[{ }^{1} L_{F} \bigoplus^{0} L_{F}^{+}: L_{F}\right] /\left[{ }^{0} L_{F}^{+}:{ }^{0} L_{F}^{-}\right]$.

Now $A_{\infty}^{\prime}$ is identified with ${ }^{1} \hat{L}_{F} \otimes R$. Let $\left\{\mu_{j}\right\}$ be a basis of ${ }^{1} L_{F}$ and let
$\left\{\hat{\mu}_{k}\right\}$ be a dual basis in ${ }^{1} \hat{L}_{F} \otimes R$ defined by $\left\langle\mu_{j}, \hat{\mu}_{k}\right\rangle=\delta_{j k}$. Take the Euclidean measure $d \lambda$ on ${ }^{1} L_{F} \otimes R$ to be the one induced by identification of ${ }^{1} L_{F} \otimes R$ with $R^{r}$ via the basis $\left\{\mu_{j}\right\}$, where $r$ is the rank of ${ }^{1} L_{F}$. Suppose we change the basis of ${ }^{1} L_{F} \otimes R$, namely, we use the Euclidean measure $\mathrm{d} \lambda^{+}$with respect to ${ }^{1} L_{F}^{+} \otimes R$. Then $\mathrm{d} \lambda^{+}=e \mathrm{~d} \lambda$ where $e=\left[{ }^{1} L_{F}^{+}:{ }^{1} L_{F}\right]$. Choose a basis $\left\{\mu_{j}^{+}\right\}$of ${ }^{1} L_{F}^{+}$such that $\left\langle\mu_{j}^{+}, \hat{\alpha}_{k}\right\rangle=$ $\delta_{j k}$, where $\left\{\alpha_{k}\right\}$ is the set of simple $F$-roots. Let

$$
\lambda: \mathbb{C}^{r} \rightarrow{ }^{1} L_{F} \otimes \mathbb{C}
$$

be the isomorphism defined by

$$
\begin{equation*}
\left\langle\lambda\left(s_{1}, \ldots, s_{r}\right), \hat{\alpha}_{k}\right\rangle=s_{k}, \quad 1 \leq k \leq r . \tag{5}
\end{equation*}
$$

That is we identify ${ }^{1} L_{F} \otimes \mathbb{C}$ with $\mathbb{C}^{r}$ via the basis $\left\{\mu_{j}^{+}\right\}$. Then $e \mathrm{~d} \lambda=$ $\mathrm{d} s_{1}, \ldots, \mathrm{~d} s_{r}$. Finally we remark that for Fourier inversion in Euclidean space, the dual measure to ${ }^{1} \hat{L}_{F} \otimes R \approx R^{r}$ is $(2 \pi i)^{-r}$ times the measure on ${ }^{1} L_{F} \otimes R$.

To summarize we have the following lemma.
6.7. Lemma: The measure induced on ${ }^{1} L_{F} \otimes C$ by that of $\left(Z_{\infty}^{+} A(F) \backslash A(\mathbb{A})\right)^{*}$ is

$$
\begin{equation*}
\mathrm{d} s_{1} \ldots \mathrm{~d} s_{r} / c \operatorname{vol}\left(A(F) \backslash A^{1}(\mathbb{A})\right)(2 \pi i)^{r} \tag{6}
\end{equation*}
$$

where

$$
c=e f=\left[L_{F}^{+}: L_{F}\right] /\left[{ }^{0} L_{F}^{+}:{ }^{0} L_{F}^{-}\right]
$$

6.8. Remark: In the remainder of this section we essentially reproduce Langlands [15] in adelic form. We follow Rapoport [18] in the proofs of lemma 6.9 and $\mathbf{6 . 1 0}$.
6.9. Lemma: All the local factors $M_{v}(w, \lambda(s))$ are holomorphic in $s$ in an open half space of $\mathbf{C}^{r}$ containing the point $(1, \ldots, 1)$.

Proof: Rewriting the formula (6) of §4 as

$$
\begin{equation*}
M_{v}(w, \lambda(s))=\prod_{\alpha \in \in_{0}^{+} \Sigma_{F}(w)} M_{v}^{G_{\alpha}}(\langle\lambda(s), \hat{\alpha}\rangle) \tag{7}
\end{equation*}
$$

we see that it is sufficient to consider the $F$-rank 1 case. And in this case, if $\phi$ is a locally constant function with compact support on $F_{v}$, then the integral of $\phi(\rho(a(\bar{n})))$ over $\bar{N}_{V}$ exists.

Thus there exists a non-negative measure $\mathrm{d} \mu$ on $F_{v}$ such that

$$
\int_{\bar{N}_{v}} \phi\left(\rho(a(\bar{n})) \mathrm{d} \bar{n}=\int_{F_{v}} \phi(t) \mathrm{d} \mu\right.
$$

for all reasonable functions $\phi$ on $F_{v}$. In particular, for $\phi: t \rightarrow|t|^{s+1}$ $(\operatorname{Re} s>t)$, we get

$$
M_{v}(s)=\int_{F_{v}}|t|^{s+1} \mathrm{~d} \mu
$$

That is $M_{v}(s)$ is the Mellin transform of a non-negative measure and is continuous at 1 (§5). 6.9 now results from a variant of Landau's lemma.
6.10. Lemma: $M(w, \lambda(s))$ is meromorphic in $s$. There exists a positive number $\epsilon$ such that the only singularities of $M(w, \lambda)$ in the region $1-\epsilon<\operatorname{Re} s_{i}<1+\epsilon \quad(i=1, \ldots, r)$ are simple poles in the hyperplane $s_{i}=1$ for $i$ corresponding to a simple positive root in ${ }_{0} \Sigma_{F}^{+}(w)$.

Proof: By the preceding lemma, we can leave out a finite number of factors $M_{v}(w, \lambda)$ from $M(w, \lambda)$. In the relative rank 1 case, up to a finite number of factors, there are four cases:
(I) $\quad M(s)=\frac{\zeta_{F}(s)}{\zeta_{F}(s+1)}$

$$
\begin{equation*}
M(s)=\zeta_{F}(2 s) \prod_{v<\infty} \frac{\left(1-\left|\tilde{\omega}_{v}\right|_{F_{v}}^{2 s+1)}\right)\left(1+\left|\tilde{\omega}_{v}\right|_{F_{v}}^{2 s+1}\right)}{\left(1+\left|\tilde{\omega}_{v}\right|_{F_{v}}^{2 s}\right)} \tag{II}
\end{equation*}
$$

(III) $\quad M(s)=\frac{\zeta_{E}(s)}{\zeta_{E}(s+1)}$

$$
\begin{equation*}
M(s)=\zeta_{E}(2 s) \prod_{v<\infty} \frac{\left(1-\left|\tilde{\omega}_{v}\right|_{E_{v}}^{2(s+1)}\right)\left(1+\left|\tilde{\omega}_{v}\right|_{E_{v}}^{2 s+1}\right)}{\left(1+\left|\tilde{\omega}_{v}\right|_{E_{v}}^{2 s}\right)} \tag{IV}
\end{equation*}
$$

where $\zeta_{F}$ (resp. $\zeta_{E}$ ) is the Dedekind zeta function of $F$ (resp. $E$ ). It is clear that in the cases (I) and (III) $M(s)$ has a simple pole at $s=1$ and
in cases (II) and (IV) $M(s)$ is holomorphic in an open half-space of $\mathbf{C}$ containing 1. The higher rank case now follows immediately from (7).
6.11. Proposition: For $\Phi, \Psi \in \mathscr{H}$, we have

$$
\begin{equation*}
(E \tilde{\phi}, \tilde{\psi})=\frac{\kappa \mu_{A}}{\mu_{G} c \tau(A)} \lim _{s \rightarrow 1} \frac{L(s, G)}{L(s, A)} M\left(w_{0}, s \rho\right) \Phi(s \rho) \bar{\Psi}(\bar{s} \rho) \tag{8}
\end{equation*}
$$

where $w_{0} \in W_{F}$ is the unique element which sends all the positive roots to negative roots.

First we introduce some functions:

$$
\begin{gathered}
f_{r}(w ; s)=M(w, \lambda(s)) \Phi(\lambda(s)) \bar{\Psi}\left(-{ }^{w} \overline{\lambda(s)}\right) \\
f_{q}\left(w ; s_{1}, \ldots, s_{q}\right)=\operatorname{Res}_{s_{q+1}=1} f_{q+1}\left(w ; s_{1}, \ldots, s_{q+1}\right) \text { for } 0 \leq q \leq \tau-1 \\
Q_{r}(s)=(\lambda(s), \lambda(s)) \\
Q_{q}\left(s_{1}, \ldots, s_{q}\right)=Q_{r}\left(s_{1}, \ldots, s_{q}, 1, \ldots, 1\right) .
\end{gathered}
$$

We also write $s^{q}$ for $\left(s_{1}, \ldots, s_{q}\right)$.
6.12. Lemma: (i) For $0 \leq q \leq r$, the functions $f_{q}\left(w, s^{q}\right)$ are meromorphic in all the $s^{q}$-spaces. In the region

$$
\left\{s^{q} \mid \operatorname{Re} s_{i}>1,1 \leq i \leq q\right\}
$$

$f_{q}\left(w, s^{q}\right)$ is holomorphic, goes to zero faster than the inverse of all polynomials as the imaginary past of $s^{q}$ goes to infinity and the real part stays in a compact subset of this region.
(ii) There exists a positive number $\epsilon$ such that the only singularities of $f_{q}\left(w ; s^{q}\right)$ in the region

$$
\left\{s^{q} \mid 1-\epsilon<\operatorname{Re} s_{i}<1+\epsilon ; i=1, \ldots, q\right\}
$$

are simple poles lying the hyperplane $s_{i}=1$.

Proof: (i) is just a restatement of the corresponding property of property of $M(w, \lambda)$ which is a consequence of the global theory of Eisenstein series (cf. [14]). (ii) follows from lemma 6.10.
6.13. It follows from $\S 6.4$ and 6.5 that
(9)

$$
c \operatorname{vol}\left(A(F) \backslash A^{1}(\mathbb{A})\right)(E \tilde{\phi}, \tilde{\psi})
$$

$$
\left.=\lambda \sum_{w \in W_{F} \in \rightarrow 0^{+}} \lim _{2 \pi i} \frac{1}{C(a, b, c, \epsilon)} \int_{(2 \pi i)^{r}} \int_{\operatorname{Re} s=s_{0}} \frac{f_{r}(w ; s)}{z-Q_{r}(s)} \mathrm{d} s_{1} \ldots \mathrm{~d} s_{r}\right\} \mathrm{d} z
$$

provided each of these limits exists. We shall show by induction that there exists the limit

$$
\begin{array}{r}
\lim _{\epsilon \rightarrow 0^{+}} \frac{1}{2 \pi i} \int_{C(a, b, c, c)} \mathrm{d} z\left\{\frac{1}{(2 \pi i)^{q}}\right.  \tag{10}\\
\left.\times \int_{\operatorname{Re}^{q}=s_{0}^{q}} \frac{f_{q}\left(w ; s^{q}\right)}{z-Q_{q}\left(s^{q}\right)} \mathrm{d} s_{1} \ldots \mathrm{~d} s_{q}\right\}
\end{array}
$$

if $s_{0}^{q}=s\left(s_{0,1}, \ldots, s_{0, q}\right)$ with $s_{0, i}>1,1 \leq i \leq q$. Note that analyticity implies that expression is independent of the actual value of $s_{0}^{q}$, provided its coordinates are strictly greater than one.

Take two small positive real numbers $u$, and $v$ such that $u$ is much smaller than $v$. Set $s_{0}^{q}=(1+u, \ldots, 1+u, 1+v)$ and $s_{0}^{q-1}=$ $(1+u, \ldots, 1+u)$. Then $Q_{q}(1+u, \ldots, 1+u, 1-v)<(\rho, \rho)$. Pick $b$ such that $Q(1+u, \ldots, 1+u, 1-v)<b<(\rho, \rho)$. Then, we can find a constant $\tau$ such that if either

$$
\left\{\begin{array}{l}
\operatorname{Re} s_{i}=1+u, \quad 1 \leq i \leq q-1 \\
\operatorname{Re} s_{q}=1-v
\end{array}\right.
$$

or

$$
\left\{\begin{array}{l}
\operatorname{Re} s_{i}=1+u, \quad 1 \leq i \leq q-1 \\
1-v \leq \operatorname{Re} s_{q} \leq 1+v \\
\left|\operatorname{Im} s_{q}\right| \geq \tau
\end{array}\right.
$$

then

$$
\operatorname{Re} Q_{q}\left(s^{q}\right)<b-\frac{1}{\tau}
$$

We integrate

$$
\frac{1}{(2 \pi i)^{q}} \int_{\text {Re } s^{q}=s_{0}^{q}} \frac{f_{q}\left(w ; s^{q}\right)}{z-Q_{q}\left(s^{q}\right)} \mathrm{d} s_{1} \ldots \mathrm{~d} s_{q}
$$

first with respect to $s_{q}$; we change the contour $\operatorname{Re} s_{q}=s_{0, q}$ to

The result is


$$
\frac{1}{(2 \pi i)^{q-1}} \int_{\operatorname{Re} s^{q-1}=s_{0}^{q-1}} \frac{f_{q-1}\left(w ; s^{q-1}\right)}{z-Q_{q-1}\left(s^{q-1}\right)} \mathrm{d} s_{1} \ldots \mathrm{~d} s_{q-1}
$$

plus

$$
\frac{1}{(2 \pi i)^{q}} \int_{\operatorname{Re} s^{q-1}=s_{0}^{q-1}}\left\{\int_{C^{q}(\nu, \tau)} \frac{f_{q}\left(w ; s^{q-1}\right)}{z-Q_{q}\left(s^{q}\right)} \mathrm{d} s_{q}\right\} \mathrm{d} s_{1} \ldots \mathrm{~d} s_{q-1} .
$$

For $s^{q-1}$ fixed and $s_{q}$ in $C^{q}(\nu, \tau)$, the image in the $Z$-plane of $C=C^{q}(\nu, \tau)$ under $Q_{q}$ is given in the following diagram


It follows that for $\operatorname{Re} s^{q-1}=s_{0}^{q-1}$ and $s_{q} \in C$ the function $1 /\left(z-Q_{q}\left(s^{q}\right)\right)$ is holomorphic in a region containing $C(a, b, c, \epsilon)$ such that

$$
\lim _{\epsilon \rightarrow 0^{+}} \int_{C(a, b, c, \epsilon)} \frac{\mathrm{d} z}{z-Q_{q}\left(s^{q}\right)}=0
$$

and (10) becomes

$$
\begin{gathered}
\lim _{\epsilon \rightarrow 0^{+}} \frac{1}{2 \pi i} \int_{C(a, b, c, \epsilon)} \mathrm{d} z \frac{1}{(2 \pi i)^{q-1}} \\
\times \int_{\operatorname{Re} s^{q-1}=s_{0}^{q-1}} \frac{f_{q-1}\left(w ; s^{q-1}\right)}{z-Q_{q-1}\left(s^{q-1}\right)} \mathrm{d} s_{1} \ldots \mathrm{~d} s_{q-1} .
\end{gathered}
$$

Finally, we get, for $q=0$

$$
\lim _{\epsilon \rightarrow 0^{+}} \frac{f_{0}(w)}{2 \pi i} \int_{C(a, b, c, c)} \frac{\mathrm{d} z}{z-(\rho, \rho)}=f_{0}(w) .
$$

But it follows from lemma 6.10 that $f_{0}(w)$ is zero unless $w=w_{0}$ and $w_{0}$ takes $\rho$ to $-\rho$. We have

$$
f_{0}\left(w_{0}\right)=\lim _{s \rightarrow 1}(s-1)^{r} M\left(w_{0}, s \rho\right) \Phi(s \rho) \bar{\Psi}\left(w_{0}(-\bar{s} \rho)\right) .
$$

Hence

$$
(E \tilde{\phi}, \tilde{\psi})=\frac{\kappa \lim _{s \rightarrow 1}(s-1)^{r} M\left(w_{0}, s \rho\right) \Phi(s \rho) \bar{\Psi}(\bar{s} \rho)}{c \operatorname{vol}\left(A(F) \backslash A^{1}(\mathbb{A})\right)}
$$

and (8) now follows from Ono's formula for Tamagawa number of the torus $A$ (cf. [17]).

Using the formula

$$
M\left(w_{0}, \lambda\right)=M_{\mathscr{C}}\left(w_{0}, \lambda\right) \prod_{v \notin \mathscr{G}} M_{v}\left(w_{0}, \lambda\right)
$$

and the result in $\S 5$ for the values of $M$, we see immediately that

$$
\begin{equation*}
(E \tilde{\phi}, \tilde{\psi})=\kappa^{2}(c \tau(A))^{-1} \Phi(\rho) \bar{\Psi}(\rho) . \tag{9}
\end{equation*}
$$

## 7. Computation of Tamagawa number

7.1. Theorem: Let $G$ be a connected reductive quasi-split group defined over an algebraic number field $F$. Let $A$ be a maximal torus of $G$ defined over $F$ lying inside the Borel subgroup of $G$ defined over $F$. Then

$$
\tau(G)=c \tau(A)
$$

where $\tau(G)($ resp. $\tau(A))$ denotes the Tamagawa number of $G$ (resp. $A)$, and $c=\left[L_{F}^{+}: L_{F}\right] /\left[{ }^{0} L_{F}^{+}:{ }^{0} L_{F}^{-}\right]$.

Proof: In the Hilbert space $\mathscr{L}^{2}\left(Z_{\infty}^{+} G(F) \backslash G(\mathbb{A})\right)$ we have

$$
\begin{equation*}
(\tilde{\phi}, 1)(1, \tilde{\psi})=(1,1)(\mathscr{P} \tilde{\phi}, \mathscr{P} \tilde{\psi}) \tag{1}
\end{equation*}
$$

According the last formula of §6, the dimension of the image of $E$ is at most one. As we have already pointed out that the constant functions are in the image of $E$, we get $E=\mathscr{P}$ and so

$$
(\mathscr{P} \tilde{\phi}, \mathscr{P} \tilde{\psi})=\kappa^{2}(c \tau(A))^{-1} \Phi(\rho) \bar{\Psi}(\rho)
$$

Since $(\tilde{\phi}, 1)=\kappa \Phi(\rho),(1, \tilde{\psi})=\kappa \bar{\Psi}(\rho)$ and $\tau(G)=(1,1)$ the theorem is proved.
7.2. Weil conjectured that the Tamagawa number of a semi-simple simply-connected connected algebraic group is one [17]. This conjecture holds for all classical groups ( ${\neq{ }^{3} D_{4},{ }^{6} D_{4} \text { ) (Tamagawa, Weil, }}_{\text {) }}$ Mars), for some exceptional groups (Mars, Demazure) and for Chevalley groups (Langlands), but it is not yet completely solved. We shall show that the Weil conjecture is true for simply-connected connected semi-simple quasi-split group $G$. This in fact follows immediately from our formula

$$
\tau(G)=c \tau(A)
$$

where $A$ is a maximal torus of $G$.
First, we observe that $G$ is simply-connected implies $L_{F}^{+}=L_{F}$, i.e. $c=1$; and the representation of the Galois group in the lattice of weights in a direct sum of permutation representation. Thus by duality theory of algebraic tori, we have

$$
A \approx \prod_{i=1}^{n} R_{E j F}\left(G_{m}\right)
$$

where $E_{i}$ are finite separable extension of $F$ which is the field of definition of $G$, and $G_{m}$ is the 1 -dimensional multiplicative group. Now we have (by Ono [17])

$$
\tau_{F}(A)=\prod_{i=1}^{n} \tau_{F}\left(R_{E_{i} / F}\left(G_{m}\right)\right)=\prod_{i=1}^{n} \tau_{E_{i}}\left(G_{m}\right)=1,
$$

because $\tau\left(G_{m}\right)=1$ (which follows from the value of the residue of zeta function $\zeta_{E}$ at 1).

Thus by the formula of the preceeding subsection $\tau(G)=c \tau(A)=1$ for a simply-connected semi-simple quasi-split connected algebraic group.

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