# Centralizers and Twisted Centralizers: Application to Intertwining Operators 

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#### Abstract

The equality of the centralizer and twisted centralizer is proved based on a case-by-case analysis when the unipotent radical of a maximal parabolic subgroup is abelian. Then this result is used to determine the poles of intertwining operators.


## 1 Introduction

The purpose of this paper is to prove the equality of the centralizer and twisted centralizer (defined in Section 2.1, originally defined by Shahidi [8]), when the unipotent radical of a maximal parabolic subgroup is abelian. In that case it is known that the adjoint action of the Levi subgroup on the Lie algebra of the unipotent radical has a finite number of orbits, the union of which is an open dense subset $[4,11]$. Then it allows the treatment in [8] of determining the poles of intertwining operators.

To be more precise, let $F$ be a non-archimedean local field of characteristic zero and $\bar{F}$ its algebraic closure. Suppose $G$ is a split connected reductive algebraic group over $F, T$ a maximal split torus of $G$. Let $\Delta$ be a set of simple roots, $\theta=\Delta \backslash\{\alpha\}$, where $\alpha$ is a simple root. Let $P=M N=M_{\theta} N$ be a maximal parabolic subgroup of $G$. Denote by $\left\{n_{i}\right\}$ a set of representatives for the corresponding open orbits of $M$ in $N$ under the adjoint action of $M$ on $\mathfrak{N}=\operatorname{Lie}(N)$. Let $N^{-}$be the opposite of $N$ and suppose one can write $w_{0}^{-1} n_{i}=m_{i} n_{i}^{\prime} n_{i}^{-}$where $m_{i} \in M, n_{i}^{\prime} \in N, n_{i}^{-} \in N^{-}$and $w_{0}$ is a representative for $\widetilde{w_{0}}$, the longest element in the Weyl group of $A_{0}$ (the maximal split torus of $T$ in $G$ ) modulo that of $A_{0}$ in $M$.

Define

$$
\begin{gathered}
M_{n_{i}}=\left\{m \in M \mid \operatorname{Int}(m) \circ n_{i}=n_{i}\right\} \\
M_{m_{i}}^{t}=\left\{m \in M \mid w_{0}(m) m_{i} m^{-1}=m_{i}\right\}
\end{gathered}
$$

Observe that $M_{n_{i}} \subset M_{m_{i}}^{t}(c f$. [8]).
It is clear that each $n_{i}$ determines $m_{i}$ uniquely (as well as $n_{i}^{\prime}$ and $n_{i}^{-}$). But the converse with respect to $m_{i}$ is not true: several $n_{i}$ could have the same $m_{i}$. The primary result of this paper proves this converse if $N$ is abelian. This is the case where the number of open orbits $\left\{n_{i}\right\}$ is finite [11]. The main result of Section 3 is:

Theorem 1.1 If $N$ is abelian, then $M_{n_{i}}=M_{m_{i}}^{t}$.

[^0]Our proof of the main theorem is based on a case-by-case analysis; all the cases where $N$ can be abelian have been listed and proved. For the exceptional groups $G_{2}$, $F_{4}$ and $E_{8}$, there is no maximal parabolic subgroup $P$ such that its unipotent subgroup $N$ is abelian. So these groups are not listed nor considered.

The method we adopt to prove this theorem is an extension of Gaussian elimination. Namely, for each orbit, we find a representative for it under $\operatorname{Ad}(M)$, which is a single element from a one dimensional subgroup corresponding to a positive root in $N$ or a product of two elements from two unipotent subgroups, attached to the longest and shortest roots in $N$, respectively. Explicitly computing the Bruhat decomposition and using the uniqueness of this decomposition, we can show that $M_{n_{i}}=M_{m_{i}}^{t}$.

This result is crucial in determining the poles of intertwining operators in [8]. To be more precise, let $X(\mathbf{M})_{F}$ be the group of $F$-rational characters of $\mathbf{M}$. Denote by A the split component of the center of $\mathbf{M}$. Then $\mathbf{A} \subset \mathbf{A}_{0}$. Let

$$
\left.\mathfrak{a}=\operatorname{Hom}\left(X(\mathbf{M})_{F}\right), \mathbb{R}\right)=\operatorname{Hom}\left(X(\mathbf{A})_{F}, \mathbb{R}\right)
$$

be the real Lie algebra of $\mathbf{A}$. Set $\mathfrak{a}^{*}=X(\mathbf{M})_{F} \otimes_{\mathbb{Z}} \mathbb{R}$ and $\mathfrak{a}_{\mathbb{C}}^{*}=\mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$ to denote its real and complex dual.

For $\nu \in \mathfrak{a}_{\mathbb{C}}^{*}$ and $\sigma$ an irreducible admissible representation of $M$, let

$$
I(\nu, \sigma)=\operatorname{Ind}_{M N \uparrow G} \sigma \otimes q^{\left\langle\nu, H_{P}(\cdot)\right\rangle} \otimes 1
$$

where $H_{P}$ is the extension of the homomorphism $H_{M}: M \rightarrow \mathfrak{a}=\operatorname{Hom}\left(X(\mathbf{M})_{F}, \mathbb{R}\right)$ to $P$, extended trivially along $N$, defined by $q^{\left\langle\chi, H_{P}(m)\right\rangle}=|\chi(m)|_{F}$ for all $\chi \in X(\mathbf{M})_{F}$. Let $V(\nu, \sigma)$ be the space of $I(\nu, \sigma)$, for $h \in V(\nu, \sigma)$, let

$$
A(\nu, \sigma, w) h(g)=\int_{N_{\bar{w}}} h\left(w^{-1} n g\right) d n
$$

where $N_{\tilde{w}}=U \cap w N^{-} w^{-1}$, be the standard intertwining operator from $I(\nu, \sigma)$ into $I(w(\nu), w(\sigma))$.

Determining the reducibility of $I(\nu, \sigma)$ at $\nu=0$ is equivalent to determining the pole of $\int_{N} h\left(w_{0}^{-1} n\right) d n$ at $\nu=0$ for any $h$ in $V(\nu, \sigma)$ which is supported in $P N^{-}$, $c f$. [6-8]. For the purpose of computing the residue we may assume that there exists a Schwartz function $\phi$ on $\mathfrak{M}^{-}$, the Lie algebra of $N^{-}$, such that $h\left(\exp \left(\mathfrak{n}^{-}\right)=\right.$ $\phi\left(\mathfrak{n}^{-}\right) h(e)$, where $\mathfrak{n}^{-} \in \mathfrak{N}^{-}$. Let $n_{i}^{-}=\exp \left(\mathfrak{n}_{i}^{-}\right)$with $\mathfrak{n}_{i}^{-} \in \mathfrak{N}^{-}$. Given a representation $\sigma$, let $\psi(m)$ be among the matrix coefficients of $\sigma$, i.e, choose an arbitrary element $\tilde{v}$ in the contragredient space of $\sigma$, let $\psi(m)=\langle\sigma(m) h(e), \tilde{v}\rangle$.

With these notations and by Theorem 2.2, $M_{m_{i}}^{t} / M_{n_{i}}=1$, (not merely finite as suggested in [8]). Proposition 2.4 [8] can be refined as:

Proposition 1.2 Let $\sigma$ be an irreducible admissible representation of $M$. Then the poles of $A\left(\nu, \sigma, w_{0}\right)$ are the same as those of

$$
\sum_{\mathfrak{n}_{i} \in O_{i}} \int_{M / M_{n_{i}}} q^{\left\langle\nu, H_{M}\left(w_{0}(m) m_{i} m^{-1}\right)\right\rangle} \phi\left(\operatorname{Ad}\left(m^{-1}\right) \mathfrak{n}_{i}^{-}\right) \psi\left(w_{0}(m) m_{i} m^{-1}\right) d \dot{m}
$$

where $O_{i}$ runs through a finite number of open orbits of $\mathfrak{N}$ under $\operatorname{Ad}(M), \mathfrak{n}_{i}$ is a representative of $O_{i}$ under the correspondence that $w_{0}^{-1} n_{i}=m_{i} n_{i}^{\prime} n_{i}^{-}$with $n_{i}=\exp \left(\mathfrak{n}_{i}\right)$, $n_{i}^{-}=\exp \left(\mathfrak{n}_{i}^{-}\right)$. Furthermore dm is the measure on $M / M_{n_{i}}$ induced from $d^{*} n_{i}$.

Let $\tilde{A}$ be the center of $M$. Then there exists a function $f \in C_{c}^{\infty}(M)$ such that $\psi(m)=\int_{\tilde{A}} f(a m) \omega^{-1}(a) d a$, where $\omega$ is the central character of $\sigma$.

Define

$$
\theta: M \rightarrow M, \theta(m)=w_{0}^{-1} m w_{0}, \forall m \in M
$$

Given $f \in C_{c}^{\infty}(M)$ and $m_{0} \in M$, define the $\theta$-twisted orbit integral for $f$ at $m_{0}$ by:

$$
\phi_{\theta}\left(m_{0}, f\right)=\int_{M / M_{\theta, m_{0}}} f\left(\theta(m) m_{0} m^{-1}\right) d \dot{m}
$$

where

$$
M_{\theta, m_{0}}=M_{\theta, m_{0}}(F)=\left\{m \in M(F) \mid \theta(m) m_{0} m^{-1}=m_{0}\right\}
$$

is the $\theta$-twisted centralizer of $m_{0}$ in $M(F), d \dot{m}$ is the measure on $M / M_{\theta, m_{0}}$ induced from $d m$.

Applying Theorem 2.2, we can restate Theorem 2.5 of [8] as:
Proposition 1.3 Assume $\sigma$ is supercuspidal and $w_{0}(\sigma) \cong \sigma$. The intertwining operator $A\left(\nu, \sigma, w_{0}\right)$ has a pole at $\nu=0$ if and only if

$$
\sum_{i} \int_{Z(G) / Z(G) \cap w_{0}(\tilde{A}) \tilde{A}^{-1}} \phi_{\theta}\left(z m_{i}, f\right) \omega^{-1}(z) d z \neq 0
$$

for $f$ as above. Here $Z(G)$ is the center of $G$ and

$$
\phi_{\theta}\left(z m_{i}, f\right)=\int_{M / M_{n_{i}}} f\left(z \theta(m) m_{i} m^{-1}\right) d \dot{m}
$$

is the $\theta$-twisted orbital integral for $f$ at $z m_{i}$, where $m_{i}$ corresponds to the representatives $\left\{n_{i}\right\}$ for the open orbits in $N$ under $\operatorname{Int}(M)$ with $w_{0}^{-1} n_{i}=m_{i} n_{i}^{\prime} n_{i}^{-}$as $n_{i}$ runs through the finite number of open orbits in $N$.

## 2 Preliminaries

Let $F$ be a non-Archimedean local field of characteristic zero. Denote by $\mathcal{O}$ its ring of integers and let $\mathcal{P}$ be the unique maximal ideal of $\mathcal{O}$. Let $q$ be the number of elements in $\mathcal{O} / \mathcal{P}$ and fix a uniformizing element $\varpi$ for which $|\varpi|=q^{-1}$, where $|\cdot|=|\cdot|_{F}$ denotes an absolute value for $F$ normalized in this way. Let $\bar{F}$ be the algebraic closure of $F$.

Let $\mathbf{G}$ be a split connected reductive algebraic group over $F$. Fix an $F$-Borel subgroup $\mathbf{B}$ and write $\mathbf{B}=\mathbf{T} \mathbf{U}$, where $\mathbf{U}$ is the unipotent radical of $\mathbf{B}$ and $\mathbf{T}$ is a maximal torus there. Let $\mathbf{A}_{0}$ be the maximal split torus of $\mathbf{T}$ and let $\Delta$ be the set of simple roots of $\mathbf{A}_{0}$ in the Lie algebra of $\mathbf{U}$.

Denote by $\mathbf{P}=\mathbf{M N}$ a maximal parabolic subgroup of $\mathbf{G}$ in the sense that $\mathbf{N} \subset \mathbf{U}$. Assume $\mathbf{T} \subset \mathbf{M}$ and let $\theta=\Delta \backslash\{\alpha\}$ such that $\mathbf{M}=\mathbf{M}_{\theta}$. As usual, we use $W=$ $W\left(\mathbf{A}_{0}\right)$ to denote the Weyl group of $\mathbf{A}_{0}$ in $\mathbf{G}$. Given $\tilde{w} \in W$, we use $w$ to denote a representative for $\tilde{w}$. Particularly, let $\tilde{w}_{0}$ be the longest element in $W$ modulo the Weyl group of $\mathbf{A}_{0}$ in $\mathbf{M}$.

We use $G, P, M, N, B, T, U, A_{0}$ to denote the subgroups of $F$-rational points of the groups $\mathbf{G}, \mathbf{P}, \mathbf{M}, \mathbf{N}, \mathbf{B}, \mathbf{T}, \mathbf{U}, \mathbf{A}_{\mathbf{0}}$, respectively. We also use $\tilde{G}, \tilde{P}, \tilde{M}, \tilde{N}, \tilde{B}, \tilde{T}, \tilde{U}, \tilde{A_{0}}$ to denote the $\bar{F}$ points of $\mathbf{G}, \mathbf{P}, \mathbf{M}, \mathbf{N}, \mathbf{B}, \mathbf{T}, \mathbf{U}, \mathbf{A}_{\mathbf{0}}$, respectively.

For any $g \in \mathbf{G}$, we will use $\operatorname{Int}(g)$ to denote the inner morphism of $\mathbf{G}$ induced by $g$, i.e., for any $u \in \mathbf{G}, \operatorname{Int}(g) \circ u=g u g^{-1}$. Let $\mathfrak{g}=\operatorname{Lie}(G)$, the Lie algebra of $G$. We will use $\operatorname{Ad}(g)$ to denote the adjoint action on $\mathfrak{g}$ induced from $\operatorname{Int}(g)$.

Suppose $R$ is the root system of $G$. For each root $\beta \in R$ we choose a root vector $\mathfrak{g}_{\beta}$ in $\mathfrak{g}$. For $\beta \in R$, let $U_{\beta}$ be the one dimensional root subgroup of $\beta$ and for $x \in F$, let $U_{\beta}(x)=\exp \left(x g_{\beta}\right)$.

Let $\mathfrak{N}=\operatorname{Lie}(N)$, the Lie algebra of $N$. Then $\mathfrak{N}=\bigoplus \mathfrak{N}_{i}$, where $\mathfrak{N}_{i}$ is graded according to $\alpha$. $M$ acts on $\mathfrak{N}$ by adjoint action. In particular, each $\mathfrak{N}_{i}$ is invariant under $\operatorname{Ad}(M)$.

For each root $\beta \in R$, there is a one dimensional subtorus $H_{\beta}(F)$, dual to $\beta$, such that the subgroup generated by $H_{\beta}, U_{\beta}$ and $U_{-\beta}$ is a simply connected group of rank one which is split over $F$. So it is isomorphic to $S L_{2}(F)$. Let $\Phi_{\beta}$ be the isomorphism from $S L_{2}(F)$ to the subgroup generated by $H_{\beta}, U_{\beta}$ and $U_{-\beta}$. Then for any $\gamma \in R$ and $t \in F^{*}$,

$$
\gamma\left(\Phi_{\beta}\left(\begin{array}{cc}
t & 0 \\
0 & t^{-1}
\end{array}\right)\right)=t^{\langle\gamma, \beta\rangle}
$$

Lemma 2.1 ( [10, Proposition 8.2.3]) Let $\beta, \gamma \in R$, with $\beta \neq \gamma$. Then there exist constants $C_{\beta, \gamma ; i, j} \in \bar{F}$, such that

$$
\left(U_{\beta}(x), U_{\gamma}(y)\right)=\prod_{\substack{i \beta+j \notin R \\ i, j>0}} U_{i \beta+j \gamma}\left(C_{\beta, \gamma ; i, j} x^{i} y^{j}\right)
$$

where the order of the factors in the right side are prescribed by a fixed ordering of $R$. Actually, the constants $C_{\beta, \gamma ; i, j}$ can be normalized so that $C_{\beta, \gamma ; i, j} \in \mathbb{Z}$. Moreover, if $\gamma$ is the longer element in the two dimensional root space spanned by $\beta$ and $\gamma$. Then $C_{\beta, \gamma ; i, j}$ can be normalized such that $C_{\beta, \gamma ; 1,1}=1$ if $\beta+\gamma \in R$. (Then $C_{\gamma, \beta, 1,1}=-1$ ).

### 2.1 Centralizer and Twisted Centralizer

Let $n_{1} \in N$, suppose $w_{0}^{-1} n_{1} \in P N^{-}$, and write $w_{0}^{-1} n_{1}=p_{1} n_{1}^{-}=m_{1} n_{1}^{\prime} n_{1}^{-}$with $m_{1} \in M, n_{1}^{\prime} \in N$ and $n_{1}^{-} \in N^{-}$. Let $\operatorname{Cent}_{M}\left(n_{1}\right)=M_{n_{1}}$ be the centralizer of $n_{1}$ in $M$, i.e.,

$$
M_{n_{1}}=\left\{m \in M \mid \operatorname{Int}(m) \circ n_{1}=n_{1}\right\}
$$

and let $M_{n_{1}^{\prime}}=\operatorname{Cent}_{M}\left(n_{1}^{\prime}\right)$ and $M_{n_{1}^{-}}=\operatorname{Cent}_{M}\left(n_{1}^{-}\right)$, respectively. Let $M_{m_{1}}^{t}=\operatorname{Cent}_{m_{1}}^{t}=$ $\left\{m \in M \mid w_{0}(m) m_{1} m^{-1}=m_{1}\right\}$ be the twisted (by means of $w_{0}$ ) centralizer of $m_{1}$ in $M$. Then by the uniqueness of $P N^{-}$decomposition of $w_{0}^{-1} n_{1}$, it is not hard to see
that the groups $M_{n_{1}}, M_{n_{1}^{-}}$and $M_{n_{1}}^{\prime}$ are all equal and are all contained in $M_{m_{1}}^{t}, c f$. [8]. Let $n_{i}=\exp \left(\mathfrak{n}_{i}\right), \mathfrak{n}_{i} \in \mathfrak{N}$, and assume the set $\left\{\mathfrak{n}_{i}\right\}$ generates a dense subset of $\mathfrak{N}$ under the action of $M$.

The main result in this paper is the following:
Theorem 2.2 Let $n_{1}=\exp \left(\mathfrak{n}_{1}\right)$, where $\mathfrak{n}_{1} \in\left\{\mathfrak{n}_{i}\right\}$ is one of the generators of a dense subset of $\mathfrak{M}$ under the action of $M$. Then $M_{n_{1}}=M_{m_{1}}^{t}$.

From the above notations, we have:

$$
\begin{equation*}
w_{0}^{-1} n_{1}=m_{1} n_{1}^{\prime} n_{1}^{-} . \tag{2.1}
\end{equation*}
$$

If $m \in M_{m_{1}}^{t}$, then

$$
\begin{align*}
w_{0}^{-1} m n_{1} m^{-1} & =\left(w_{0}(m) m_{1} m^{-1}\right)\left(m n_{1}^{\prime} m^{-1}\right)\left(m n_{1}^{-} m^{-1}\right)  \tag{2.2}\\
& =m_{1}\left(m n_{1}^{\prime} m^{-1}\right)\left(m n_{1}^{-} m^{-1}\right)
\end{align*}
$$

For convenience of notation, Let

$$
n_{2}=\operatorname{Int}(m) \circ n_{1}, \quad n_{2}^{\prime}=\operatorname{Int}(m) \circ n_{1}^{\prime}, \quad n_{2}^{-}=\operatorname{Int}(m) \circ n_{1}^{-}
$$

Then equation (2.2) will be changed to:

$$
\begin{equation*}
w_{0}^{-1} n_{2}=m_{1} n_{2}^{\prime} n_{2}^{-} \tag{2.3}
\end{equation*}
$$

Multiplying the inverse of equation (2.3) by equation (2.1), we have:

$$
\begin{equation*}
n_{2}^{-1} n_{1}=\left(n_{2}^{-}\right)^{-1}\left(n_{2}^{\prime}\right)^{-1} n_{1}^{\prime} n_{1}^{-} \tag{2.4}
\end{equation*}
$$

Let

$$
\begin{gathered}
s_{1}=n_{2}^{-1} n_{1} \in N, \quad s_{1}^{-}=\left(n_{1}^{-}\right)^{-1} \in N^{-} \\
s_{2}^{-}=\left(n_{2}^{-}\right)^{-1} \in N^{-}, \quad s_{2}=\left(n_{2}^{\prime}\right)^{-1} n_{1}^{\prime} \in N .
\end{gathered}
$$

Then equation (2.4) becomes

$$
\begin{equation*}
s_{1} s_{1}^{-}=s_{2}^{-} s_{2} \tag{2.5}
\end{equation*}
$$

Let

$$
\begin{aligned}
n_{1} & =\exp \left(\mathfrak{n}_{1}\right), & n_{2} & =\exp \left(\mathfrak{n}_{2}\right) \\
s_{1} & =\exp \left(r_{1}\right), & s_{2} & =\exp \left(r_{2}\right) \\
s_{1}^{-} & =\exp \left(r_{1}^{-}\right), & s_{2}^{-} & =\exp \left(r_{2}^{-}\right)
\end{aligned}
$$

Then $\mathfrak{n}_{2}=\operatorname{Ad}(m) \circ \mathfrak{n}_{1}$ is one of the generators of a dense orbit of $\mathfrak{M}$ under $\operatorname{Ad}(M)$ since $\mathfrak{n}_{1}$ is. Similarly it is not hard to see that both $r_{1}^{-}$and $r_{2}^{-}$are generators of a dense orbit of $\mathfrak{M}^{-}$.

Our goal is to prove:
Claim Under the assumption in Theorem 2.2, we must have: $s_{1}^{-}=s_{2}^{-}$.

Once this has been proved, it implies $n_{1}^{-}=n_{2}^{-}$, which will lead to $n_{1}=n_{2}$ by the uniqueness of $P N^{-}$decomposition. Since $m \in M_{m_{1}}^{t}$ and $n_{2}=\operatorname{Int}(m) \circ n_{1}$, we get $m \in M_{n_{1}}$ if $m \in M_{m_{1}}^{t}$. So $M_{m_{1}}^{t} \subset M_{n_{1}}$. But we already have $M_{n_{1}} \subset M_{m_{1}}^{t}$,cf. [8]. So $M_{n_{1}}=M_{m_{1}}^{t}$ as desired.

Remark We can always assume that $s_{2}^{-} \neq 1$, since otherwise there is nothing that needs to be done. We are going to prove the claim according to the type of Dynkin diagram of $G$ since the Gaussian elimination essentially depends on the structure of the root system.

Strategy of Proof Except for some simple cases (like $A_{l}, C_{l}$ ), our proof relies on Gaussian elimination for $\mathfrak{R}$. Namely, $\mathfrak{N}$ can be generated by $\mathfrak{g}_{\beta}$ with $\beta$ a positive root in $N$, or by $\mathfrak{g}_{\beta}, \mathfrak{g}_{\gamma}$ under $\operatorname{Ad}(M)$, where $\mathfrak{g}_{\beta}, \mathfrak{g}_{\gamma}$ are root vectors attached to the shortest and longest roots in $N$. Thus by acting with a suitable $m \in M$ on both sides of equation (2.5), we can always assume that $s_{2}=U_{\beta}\left(a_{1}\right) U_{\gamma}\left(a_{2}\right)$ or $U_{\beta}\left(a_{1}\right)$.

We will multiply both sides of equation (2.5) by $U_{\beta}(x) U_{\gamma}(y)$ from the right, where $x, y$ are variables. Then the $M$-parts of $s_{1} s_{1}^{-} U_{\beta}(x) U_{\gamma}(y)$ and $s_{2}^{-} s_{2} U_{\beta}(x) U_{\gamma}(y)$ can be calculated and compared explicitly since they are in the simplest form. We can then conclude that their $M$-parts will never be equal unless $s_{1}^{-}=s_{2}^{-}$.

## 3 Proof of the Main Theorem

Now suppose $N$ is abelian, then $\operatorname{Ad}(M)$ acts on $\mathfrak{N}$ having finite number of orbits, cf. $[4,11]$.

### 3.1 Roots in Unipotent Radical

Lemma 3.1 Suppose $N$ is abelian. If

$$
\beta=c \alpha+\sum_{\alpha_{i} \neq \alpha} c_{i} \alpha_{i}
$$

is a positive root of $N$ where $\alpha_{i}$ 's are simple roots from $\theta$, then $c=1$.
Proof Using [3, Corollary of Lemma A $\S 10.2$ ], $\beta$ can be written in the form $\beta_{1}+$ $\beta_{2}+\cdots+\beta_{k}$ with $\beta_{i} \in \Delta\left(\beta_{i}\right.$ not necessary distinct $)$ such that each partial sum $\beta_{1}+\beta_{2}+\cdots+\beta_{j}$ is a root $(1 \leq j \leq k)$. Suppose $c \geq 2$, then there is $j$ such that $\beta_{j}=\alpha$ and in the remaining partial sum $\beta_{1}+\beta_{2}+\cdots+\beta_{j-1}$, there is still one $\alpha$. Let $\gamma=\beta_{1}+\beta_{2}+\cdots+\beta_{j-1}$, then $\mathfrak{g}_{\gamma}, \mathfrak{g}_{\beta_{j}} \in \mathfrak{N}$, and $\left[\mathfrak{g}_{\gamma}, \mathfrak{g}_{\beta_{j}}\right]=\mathfrak{g}_{\beta_{1}+\beta_{2}+\cdots+\beta_{j}} \neq 0$. This is a contradiction to $\mathfrak{M}$ being abelian.

If

$$
P=\sum_{\substack{\alpha_{i} \in \Delta \\ i=1}}^{k} c_{i} \alpha_{i}
$$

is a root, choose $k$ points in a plane representing each $\alpha_{i}$ and draw a line connecting $\alpha_{i}, \alpha_{j}$, if $\left\langle\alpha_{i}, \hat{\alpha_{j}}\right\rangle \neq 0$. Then the graph obtained is obviously a subgraph of the Dynkin diagram and is composed of several connected pieces. For each connected piece $C_{i}$ of this graph, we set

$$
P_{i}=\sum_{\alpha_{i} \in C_{i}} c_{i} \alpha_{i}
$$

Then

$$
P=\sum_{i}^{m} P_{i}
$$

where $m$ is the number of connected pieces. All the $C_{i}$ 's are disjoint. We call $P_{i}$ a connected piece of $P$. Call $P_{i}$ positive if each $c_{i}$ is positive, and negative if each $c_{i}$ is negative. In particular, we call $P$ a connected root if $P$ is composed of only one connected piece.

Lemma 3.2 Every positive root is connected.
Proof Let

$$
r=\sum_{i=1}^{k} P_{i}
$$

be a positive root with all $P_{i}$ 's being positive connected and disjoint with each other. Then by [3, Corollary of Lemma A §10.2], $r$ can be written as

$$
r=\sum_{i=1}^{n} \alpha_{i}
$$

such that every partial sum

$$
r_{s}=\sum_{i=1}^{s} \alpha_{i}, \quad 1 \leq s \leq n
$$

is a root. If $k>1$, then there must be one $s, s>1$, and one $i, 1 \leq i \leq k$, such that in the sum for $r_{s}$, there is only one element, say $\alpha_{j}, 1 \leq j \leq s$, which comes from $P_{i}$. Then for all $\alpha_{i}, 1 \leq i \leq s, i \neq j,\left\langle\alpha_{i}, \hat{\alpha_{j}}\right\rangle=0$ since $\alpha_{i}, \alpha_{j}$ are not in the same connected piece. So

$$
S_{\alpha_{j}}\left(r_{s}\right)=r_{s}-\left\langle r_{s}, \widehat{\alpha_{j}}\right\rangle \alpha_{j}=\sum_{i=1}^{s} \alpha_{i}-2 \alpha_{j}=\sum_{\substack{i=1 \\ i \neq j}}^{s} \alpha_{i}-\alpha_{j},
$$

where $S_{\alpha_{j}}$ is the reflection about $\alpha_{j}$ in the Weyl group of $G$. Since none of the $\alpha_{i}$ 's in the sum

$$
\sum_{i=1, i \neq j}^{s} \alpha_{i}
$$

can be $\alpha_{j}$, and all $\alpha_{i}$ are simple roots, $S_{\alpha_{j}}\left(r_{s}\right)$ is not a root. This is a contradiction to $S_{\alpha_{j}}\left(r_{s}\right)$ being a root since $r_{s}$ is a root.

### 3.2 Type $A_{l}$

Equation (2.5) implies $\exp \left(r_{1}\right) \exp \left(r_{1}^{-}\right)=\exp \left(r_{2}^{-}\right) \exp \left(r_{2}\right)$.
Since $r_{1}^{2}=r_{2}^{2}=\left(r_{1}^{-}\right)^{2}=\left(r_{2}^{-}\right)^{2}=0$, we have:

$$
\begin{equation*}
r_{1}+r_{1}^{-}+r_{1} r_{1}^{-}=r_{2}^{-}+r_{2}+r_{2}^{-} r_{2} . \tag{3.1}
\end{equation*}
$$

Choose $t \in T$ and let $\operatorname{Ad}(t)$ act on both sides of equation (3.1). We get

$$
\alpha(t) r_{1}+\alpha^{-1}(t) r_{1}^{-}+r_{1} r_{1}^{-}=\alpha^{-1}(t) r_{2}^{-}+\alpha(t) r_{2}+r_{2}^{-} r_{2}
$$

Since this is true for all $t \in T$, we must have $r_{1}=r_{2}, r_{1}^{-}=r_{2}^{-}$. Consequently, $s_{1}^{-}=s_{2}^{-}$.

### 3.3 Type $B_{l}$

In this case, we may assume that $T$ can be chosen to be the set of matrices of the form:

$$
\operatorname{diag}\left(x_{1}, x_{2}, \ldots, x_{l}, x_{1}^{-1}, x_{2}^{-1}, \ldots, x_{l}^{-1}, 1\right)
$$

since the unipotent subgroups remain unchanged in every adjoint action.
The Dynkin diagram of $G$ is:


Let $e_{i} \in \operatorname{Hom}\left(T, F^{*}\right), 1 \leq i \leq l$ such that $e_{i}(T)=x_{i}$. Then $\alpha_{i}=e_{i}-e_{i+1}, 1 \leq$ $i \leq l-1 ; \alpha_{l}=e_{l}$. The only case when $N$ can be abelian is $\alpha=\alpha_{1}$. Then the positive roots in $N$ are: $\left\{e_{1} \pm e_{i} \mid 2 \leq i \leq l\right\} \cup\left\{e_{1}\right\}$.

We choose a root vector for each positive root in $G$ as follows:

$$
\begin{aligned}
\mathfrak{g}_{e_{i}-e_{j}} & =E_{i, j}-E_{l+j, l+i}, & & 1 \leq i<j \leq l \\
\mathfrak{g}_{e_{i}+e_{j}} & =E_{i, l+j}-E_{j, l+i}, & & 1 \leq i<j \leq l \\
\mathfrak{g}_{e_{i}} & =E_{i, 2 l+1}-E_{2 l+1, l+i}, & & 1 \leq i \leq l .
\end{aligned}
$$

We also choose a root vector for each negative root in $G$ as follows:

$$
\begin{aligned}
\mathfrak{g}_{-e_{i}+e_{j}} & =E_{j, i}-E_{l+i, l+j}, & & 1 \leq i<j \leq l, \\
\mathfrak{g}_{-e_{i}-e_{j}} & =E_{l+j, i}-E_{l+i, j}, & & 1 \leq i<j \leq l, \\
\mathfrak{g}_{-e_{i}} & =E_{l+i, 2 l+1}-E_{2 l+1, i}, & & 1 \leq i \leq l,
\end{aligned}
$$

where the $E_{i, j}$ 's are elementary matrices in $M_{(2 l+1) \times(2 l+1)}$ such that its $(i, j)$ entry is 1 , all other entries are 0 .

Lemma 3.3 Given any nonzero element

$$
r=\sum_{i=2}^{l} a_{i} \mathfrak{g}_{e_{1}-e_{i}}+\sum_{i=2}^{l} b_{i} \mathfrak{g}_{e_{1}+e_{i}}+c \mathfrak{g}_{e_{1}} \in \mathfrak{N}
$$

there is an $m \in M$, such that $\operatorname{Ad}(m) \circ r=c_{0} \mathfrak{g}_{e_{1}-e_{2}}+c_{1} g_{e_{1}+e_{2}}$ with $c_{0} \neq 0$.

Proof This is [9, Lemma 4.2].
Lemma 3.4 For an element $r=c_{0} \mathfrak{g}_{e_{1}-e_{2}}+c_{1} \mathfrak{g}_{e_{1}+e_{2}} \in \mathfrak{N}$ from Lemma 3.3 with $c_{1} \neq 0$, there is $m \in \tilde{M}$ such that $\operatorname{Ad}(m) \circ r=a \mathfrak{g}_{e_{1}}$ with $a \neq 0$.

Proof Choose $x \in \bar{F}$ such that $\frac{1}{2} c_{0} x^{2}=c_{1}$. Let $m=U_{-e_{2}}\left(\frac{1}{x}\right) U_{e_{2}}(x)$. Then $\operatorname{Ad}(m) \circ$ $r=-c_{0} x g_{e_{1}}$. Setting $a=-c_{0} x$ finishes the proof.

We start with equation (2.5). If $s_{2}=1$, then it immediately follows $s_{1}^{-}=s_{2}^{-}$, and there is nothing to do. So suppose $s_{2} \neq 1$. By the above two lemmas, applying a suitable $\operatorname{Int}(m), m \in \tilde{M}$ on both sides if necessary, we can assume $s_{2}=U_{e_{1}}(a)$ or $U_{e_{1}-e_{2}}(a)$ with $a \neq 0$. By taking a suitable finite extension of $F$, we can always assume that $m \in M$ and consequently $a \in F$. Without loss of generality, we assume $s_{2}=U_{e_{1}}(a)$.

Suppose

$$
\begin{aligned}
& s_{1}^{-}=\prod_{k=2}^{l} U_{-e_{1}-e_{k}}\left(a_{k}\right) \prod_{k=2}^{l} U_{-e_{1}+e_{k}}\left(b_{k}\right) U_{-e_{1}}\left(x_{0}\right) \\
& s_{2}^{-}=\prod_{k=2}^{l} U_{-e_{1}-e_{k}}\left(c_{k}\right) \prod_{k=2}^{l} U_{-e_{1}+e_{k}}\left(d_{k}\right) U_{-e_{1}}\left(y_{0}\right)
\end{aligned}
$$

Multiply both sides of (2.5) by $u=U_{e_{1}}(x) \in N$ on the right, where $x \in F$. Decompose both $s_{1} s_{1}^{-} u$ and $s_{2}^{-} s_{2} u$ into $P N^{-}$form, and compare their $M$ part. Their $M$ part will never be equal unless $s_{1}^{-}=s_{2}^{-}$. The reason for multiplying $u$ is to exclude the possibility of occurrence of some Weyl group elements (when $a y_{0}=-1$ ).

First we have

$$
U_{-e_{1}}\left(y_{0}\right) U_{e_{1}}(a+x)=U_{e_{1}}\left(\frac{a+x}{1+y_{0}(a+x)}\right) h_{2, x} U_{-e_{1}}\left(\frac{y_{0}}{1+y_{0}(a+x)}\right)
$$

where

$$
h_{2, x}=\Phi_{e_{1}}\left(\begin{array}{cc}
\frac{1}{1+y_{0}(a+x)} & 0 \\
0 & 1+y_{0}(a+x)
\end{array}\right) \in T
$$

Set

$$
a_{x}=\frac{a+x}{1+y_{0}(a+x)} .
$$

For any $k, 2 \leq k \leq l$, by Lemma 2.1,

$$
\begin{aligned}
& U_{-e_{1}+e_{k}}\left(d_{k}\right) U_{e_{1}}\left(a_{x}\right)=U_{e_{1}}\left(a_{x}\right) U_{e_{k}}\left(d_{k} a_{x}\right) U_{-e_{1}+e_{k}}\left(d_{k}\right), \\
& U_{-e_{1}-e_{k}}\left(c_{k}\right) U_{e_{1}}\left(a_{x}\right)=U_{e_{1}}\left(a_{x}\right) U_{-e_{k}}\left(c_{k} a_{x}\right) U_{-e_{1}-e_{k}}\left(c_{k}\right)
\end{aligned}
$$

Then by recursively applying Lemma 2.1 and using the fact that $N$ and $N^{-}$are normal in $P$ and $P^{-}$respectively, it can be calculated that the $M$ part of $s_{2}^{-} s_{2} u$ is:

$$
m_{2}=\prod_{k=2}^{l} U_{-e_{k}}\left(c_{k} a_{x}\right) \prod_{k=2}^{l} U_{e_{k}}\left(d_{k} a_{x}\right) h_{2, x}
$$

Similarly, if we set

$$
b_{x}=\frac{x}{1+x_{0} x} \quad \text { and } \quad h_{1, x}=\Phi_{e_{1}}\left(\begin{array}{cc}
\frac{1}{1+x_{0} x} & 0 \\
0 & 1+x_{0} x
\end{array}\right) \in T
$$

then the $M$ part of $s_{1} s_{1}^{-} u$ is:

$$
m_{1}=\prod_{k=2}^{l} U_{-e_{k}}\left(a_{k} b_{x}\right) \prod_{k=2}^{l} U_{e_{k}}\left(b_{k} b_{x}\right) h_{1, x}
$$

From equation (2.5), $s_{1} s_{1}^{-} u=s_{2}^{-} s_{2} u$. By the uniqueness of $M N N^{-}$decomposition, we must have $m_{1}=m_{2}$. Since $m_{1}$ and $m_{2}$ are products of unipotent groups attached to roots in $M$ in the same order, we must have $c_{k} a_{x}=a_{k} b_{x}$ and $d_{k} a_{x}=b_{k} b_{x}$ for almost all $x \in F$ and all $k, 2 \leq k \leq l$. These equations lead to:

$$
\begin{align*}
& \left(c_{k} x_{0}-a_{k} y_{0}\right) x^{2}+\left(c_{k} x_{0}+a c_{k} x_{0}-a_{k}-a a_{k} y_{0}\right) x+a c_{k}=0  \tag{3.2}\\
& \left(d_{k} x_{0}-b_{k} y_{0}\right) x^{2}+\left(d_{k} x_{0}+a d_{k} x_{0}-b_{k}-a b_{k} y_{0}\right) x+a d_{k}=0 \tag{3.3}
\end{align*}
$$

For equations (3.2) and (3.3) to have infinitely many solutions, one must have $a_{k}=$ $b_{k}=c_{k}=d_{k} \equiv 0, \forall k, 2 \leq k \leq l$, since $a \neq 0$ by assumption. Moreover, we have $h_{1, x}=h_{2, x}$ for almost all $x$, which means the equation

$$
\left(y_{0}-x_{0}\right) x+a y_{0}=0
$$

has infinitely many solutions, thus $y_{0}=0$, so $s_{2}^{-}=1$, which is a contradiction. So in order that equation (2.5) holds, we must have $s_{2}=1$, which leads to $s_{1}^{-}=s_{2}^{-}$. When $s_{2}=U_{e_{1}-e_{2}}(a)$, we can also prove that $s_{1}^{-}=s_{2}^{-}$in a similar way. That finishes the proof of the main theorem in case $G$ is of type $B_{l}$.

### 3.4 Type $C_{l}$

In this case, we may assume $T$ is the set of matrices of the form:

$$
\operatorname{diag}\left(x_{1}, x_{2}, \ldots, x_{l}, x_{1}^{-1}, x_{2}^{-1}, \ldots, x_{l}^{-1}\right)
$$

since the unipotent subgroups remain unchanged in every adjoint action.
Let $e_{i} \in \operatorname{Hom}\left(T, F^{*}\right)$ such that $e_{i}(H)=x_{i}$. Then $R=\left\{ \pm e_{i} \pm e_{j} \mid i \neq j\right\} \cup\left\{ \pm 2 e_{k}\right\}$. $N$ is abelian only in case $\alpha=2 e_{l}$. In this case, $\Delta=\left\{e_{i}-e_{i+1} \mid 1 \leq i \leq l-1\right\} \cup\left\{2 e_{l}\right\}$. The positive roots in $N$ are: $R^{+} \backslash \theta^{+}=\left\{e_{i}+e_{j} \mid i \neq j\right\} \cup\left\{2 e_{i} \mid 1 \leq i \leq l\right\}$. And $\mathfrak{N}$ is all the $2 l \times 2 l$ matrices of the form:

$$
\left(\begin{array}{ll}
0 & Y \\
0 & 0
\end{array}\right)
$$

where $Y \in M_{l}(F)$ and $Y^{t}=Y$. So for each $\mathfrak{n} \in \mathfrak{N}, \mathfrak{n}^{2}=0$, and for each $\mathfrak{n}^{-} \in$ $\mathfrak{N}^{-}, \mathfrak{n}^{-2}=0$. It can be seen that the proof for the $A_{l}$ case also applies in this case which implies $s_{1}^{-}=s_{2}^{-}$.

### 3.5 Type $D_{l}$

In this case, again $T$ may be considered to be the set of matrices of the form:

$$
\operatorname{diag}\left(x_{1}, x_{2}, \ldots, x_{l}, x_{1}^{-1}, x_{2}^{-1}, \ldots, x_{l}^{-1}\right)
$$

because the unipotent subgroups remain unchanged in every adjoint action.
Let $e_{i} \in \operatorname{Hom}\left(T, F^{*}\right)$ such that $e_{i}(H)=x_{i}$. Then $R=\left\{ \pm e_{i} \pm e_{j} \mid i \neq j\right\}$, $\Delta=\left\{e_{i}-e_{i+1} \mid 1 \leq i \leq l-1\right\} \cup\left\{e_{l-1}+e_{l}\right\}$. Let $\alpha_{i}=e_{i}-e_{i+1}$ for $1 \leq i \leq l-1$, and let $\alpha_{l}=e_{l-1}+e_{l}$. For $N$ to be abelian, $\alpha$ must be $\alpha_{1}, \alpha_{l-1}$ or $\alpha_{l}$. If $\alpha=\alpha_{l-1}$, then every element $\mathfrak{n} \in \mathfrak{N}$ has the form:

$$
\left(\begin{array}{cc}
A & Y \\
0 & -A^{t}
\end{array}\right)
$$

where

$$
A=\left(\begin{array}{cccccc}
0 & 0 & 0 & \cdots & 0 & a_{1} \\
0 & 0 & 0 & \cdots & 0 & a_{2} \\
\vdots & \vdots & \vdots & & \vdots & \vdots \\
0 & 0 & 0 & \cdots & 0 & a_{l-1} \\
0 & 0 & 0 & \cdots & 0 & 0
\end{array}\right), \quad Y=\left(\begin{array}{cc}
B & 0 \\
0 & 0
\end{array}\right)
$$

and $B \in M_{l-1}(F), B=-B^{t}$. Then it is easily checked that $\mathfrak{n}^{2}=0$ and consequently for each $\mathfrak{n}^{-} \in \mathfrak{N}^{-}, \mathfrak{n}^{-2}=0$. Again we can use the same method as in $A_{l}$ or $C_{l}$ to prove that $s_{1}^{-}=s_{2}^{-}$.

The symmetry between $\alpha_{l}$ and $\alpha_{l-1}$ takes care of the case $\alpha=\alpha_{l}$.
If $\alpha=\alpha_{1}$, then $\mathfrak{N}$ does not have the property that for each $\mathfrak{n} \in \mathfrak{M}, \mathfrak{n}^{2}=0$. In this case, the positive roots in $N$ are: $\left\{e_{1}-e_{i} \mid 1<i \leq l\right\} \cup\left\{e_{1}+e_{j} \mid 1<j \leq l\right\}$.

We choose a root vector for each positive root in $G$ as follows:

$$
\begin{array}{cl}
\mathfrak{g}_{e_{i}-e_{j}}=E_{i, j}-E_{l+j, l+i}, & 1 \leq i<j \leq l \\
\mathfrak{g}_{e_{i}+e_{j}}=E_{i, l+j}-E_{j, l+i}, & 1 \leq i<j \leq l
\end{array}
$$

We also choose a root vector for each negative root in $G$ as follows:

$$
\begin{array}{ll}
\mathfrak{g}_{-e_{i}+e_{j}}=E_{j, i}-E_{l+i, l+j}, & 1 \leq i<j \leq l, \\
\mathfrak{g}_{-e_{i}-e_{j}}=E_{l+j, i}-E_{l+i, j}, & 1 \leq i<j \leq l,
\end{array}
$$

where the $E_{j, k}$ 's are elementary matrix in $M_{2 l \times 2 l}$. Then $\left\{\mathfrak{g}_{e_{1} \pm e_{i}} \mid 1<i \leq l\right\}$ is a basis for $\mathfrak{N}$.

Theorem 3.5 (Gaussian Elimination) For any nonzero $r \in \mathfrak{N}$, there exist $m \in M$ and $k_{0}, k_{1} \in F$, with $k_{0} \neq 0$, such that $\operatorname{Ad}(m) \circ r=k_{0} \mathfrak{g}_{e_{1}-e_{2}}+k_{1} \mathfrak{g}_{e_{1}+e_{2}}$.

Proof Suppose

$$
r=\sum_{i=1}^{l-1} a_{i} \mathfrak{g}_{e_{1}-e_{i+1}}+\sum_{i=1}^{l-1} a_{i}^{\prime} \mathfrak{g}_{e_{1}+e_{i+1}}
$$

We first prove that by applying a suitable $m^{\prime} \in M$ on $r$ if necessary, we can always assume that $a_{1} \neq 0$.

Assume $a_{1}=0$. Let

$$
m^{\prime}= \begin{cases}U_{-e_{2}+e_{i+1}}(1) & \exists i, 2 \leq i \leq l-1, \text { such that } a_{i} \neq 0 \\ U_{-e_{2}-e_{i+1}}(1) & \exists i, 2 \leq i \leq l-1, \text { such that } a_{i}^{\prime} \neq 0 \\ s_{e_{2}} & \text { otherwise }\end{cases}
$$

where $s_{e_{2}}$ is a representative of the Weyl group element $S_{e_{2}}$, which is the reflection about $e_{2}$.

By applying the formula $\operatorname{Ad}\left(\exp \left(x g_{\beta}\right)\right)=e^{\operatorname{ad}\left(x g_{\beta}\right)}$ for each root $\beta \in R$, it is easily checked that the coefficient of $\mathfrak{g}_{e_{1}-e_{2}}$ in $\operatorname{Ad}\left(m^{\prime}\right) \circ r$ is nonzero.

Let $k_{0}=a_{1}$, and

$$
m=\left[\prod_{i=3}^{l} \exp \left(\frac{a_{i-1}^{\prime}}{k_{0}} \mathfrak{g}_{e_{2}+e_{i}}\right)\right] \cdot\left[\prod_{i=3}^{l} \exp \left(\frac{a_{i-1}}{k_{0}} \mathfrak{g}_{e_{2}-e_{i}}\right)\right]
$$

Then a direct calculation shows that

$$
\operatorname{Ad}(m) \circ r=k_{0} \mathfrak{g}_{e_{1}-e_{2}}+\left(a_{l-1}^{\prime}+\sum_{i=3}^{l} \frac{a_{i-1} \cdot a_{i-1}^{\prime}}{k_{0}}\right) \mathfrak{g}_{e_{1}+e_{2}}
$$

Let $k_{1}$ denote the coefficient of $\mathfrak{g}_{e_{1}+e_{2}}$ from the right-hand side of the above equation, then $\operatorname{Ad}(m) \circ r=k_{0} \mathfrak{g}_{e_{1}-e_{2}}+k_{1} \mathfrak{g}_{e_{1}+e_{2}}$ as desired.

Considering equation (2.5), if $s_{2} \neq 0$, by Theorem 3.5, applying an $m \in M$ on both sides, we can assume $s_{2}=U_{e_{1}-e_{2}}\left(k_{0}\right) \cdot U_{e_{1}+e_{2}}\left(k_{1}\right)$ with $k_{0} \neq 0$.

Suppose

$$
\begin{aligned}
& s_{1}^{-}=\prod_{i=2}^{l} U_{-e_{1}-e_{i}}\left(a_{i}\right) \prod_{i=2}^{l} U_{-e_{1}+e_{i}}\left(b_{i}\right) \\
& s_{2}^{-}=\prod_{i=2}^{l} U_{-e_{1}-e_{i}}\left(c_{i}\right) \prod_{i=2}^{l} U_{-e_{1}+e_{i}}\left(d_{i}\right)
\end{aligned}
$$

We will adopt the strategy we have used in the case of $B_{l}$ : multiply both sides of (2.5) by $u=U_{e_{1}-e_{2}}(x) U_{e_{1}+e_{2}}(y) \in N$ on the right, where $x, y$ are variables in $F$. Decompose both $s_{1} s_{1}^{-} u$ and $s_{2}^{-} s_{2} u$ into $P N^{-}$form and compare their $M$ parts.

Now let us consider the $P N^{-}$decomposition of $s_{1} s_{1}^{-} u$ and $s_{2}^{-} s_{2} u$. For $s_{1}^{-} U_{e_{1}-e_{2}}(x)$, first we have:

$$
\begin{equation*}
U_{-e_{1}+e_{2}}\left(b_{2}\right) U_{e_{1}-e_{2}}(x)=U_{e_{1}-e_{2}}\left(x^{\prime}\right) h_{0, x} U_{-e_{1}+e_{2}}\left(\frac{b_{2}}{1+b_{2} x}\right) \tag{3.4}
\end{equation*}
$$

where

$$
x^{\prime}=\frac{x}{1+b_{2} x} \quad \text { and } \quad h_{0, x}=\Phi_{e_{1}-e_{2}}\left(\begin{array}{cc}
\frac{1}{1+b_{2} x} & 0 \\
0 & 1+b_{2} x
\end{array}\right) \in T
$$

For each $i, 3 \leq i \leq l$, by applying Lemma 2.1, we get:

$$
\begin{align*}
U_{-e_{1}+e_{i}}\left(b_{i}\right) U_{e_{1}-e_{2}}\left(x^{\prime}\right) & =U_{e_{1}-e_{2}}\left(x^{\prime}\right) U_{-e_{2}+e_{i}}\left(b_{i} x^{\prime}\right) U_{-e_{1}+e_{i}}\left(b_{i}\right)  \tag{3.5}\\
U_{-e_{1}-e_{i}}\left(a_{i}\right) U_{e_{1}-e_{2}}\left(x^{\prime}\right) & =U_{e_{1}-e_{2}}\left(x^{\prime}\right) U_{-e_{2}-e_{i}}\left(a_{i} x^{\prime}\right) U_{-e_{1}-e_{i}}\left(a_{i}\right) \tag{3.6}
\end{align*}
$$

And $U_{e_{1}-e_{2}}$ commutes with $U_{-e_{1}-e_{2}}$.
From equations (3.5), (3.6) and using the fact that both $N$ and $N^{-}$are normal in $P$ and $P^{-}$, respectively, we reach the following:

$$
\begin{equation*}
s_{1}^{-} U_{e_{1}-e_{2}}(x)=U_{e_{1}-e_{2}}\left(x^{\prime}\right) \prod_{i=3}^{l} U_{-e_{2}-e_{i}}\left(a_{i} x^{\prime}\right) \prod_{i=3}^{l} U_{-e_{2}+e_{i}}\left(b_{i} x^{\prime}\right) h_{0, x_{1}}^{-\prime} \tag{3.7}
\end{equation*}
$$

for a suitable $s_{1}^{-\prime} \in N^{-}$. When $s_{2}^{-} s_{2} u=s_{2}^{-} U_{e_{1}-e_{2}}\left(k_{0}+x\right) U_{e_{1}+e_{2}}\left(k_{1}+y\right)$, a similar calculation shows that

$$
\begin{equation*}
s_{2}^{-} U_{e_{1}-e_{2}}\left(k_{0}+x\right)=U_{e_{1}-e_{2}}\left(k_{0, x}\right) \prod_{i=3}^{l} U_{-e_{2}-e_{i}}\left(c_{i} k_{0, x}\right) \prod_{i=3}^{l} U_{-e_{2}+e_{i}}\left(d_{i} k_{0, x}\right) h_{0, x}^{\prime} s_{2}^{-\prime} \tag{3.8}
\end{equation*}
$$

for a suitable $s_{2}^{-\prime} \in N^{-}$, where

$$
k_{0, x}=\frac{k_{0}+x}{1+d_{2}\left(k_{0}+x\right)} \quad \text { and } \quad h_{0, x}^{\prime}=\Phi_{e_{1}-e_{2}}\left(\begin{array}{cc}
\frac{1}{1+d_{2}\left(k_{0}+x\right)} & 0 \\
0 & 1+d_{2}\left(k_{0}+x\right)
\end{array}\right) \in T
$$

## Suppose

$$
\begin{aligned}
s_{1}^{\prime \prime} & =\prod_{i=2}^{l} U_{-e_{1}-e_{i}}\left(a_{i}^{\prime}\right) \prod_{i=2}^{l} U_{-e_{1}+e_{i}}\left(b_{i}^{\prime}\right), \\
s_{2}^{-\prime} & =\prod_{i=2}^{l} U_{-e_{1}-e_{i}}\left(c_{i}^{\prime}\right) \prod_{i=2}^{l} U_{-e_{1}+e_{i}}\left(d_{i}^{\prime}\right) .
\end{aligned}
$$

Then with a similar calculation as above, by applying Lemma 2.1 recursively, we get:

$$
\begin{equation*}
\left(s_{1}^{-}\right)^{\prime} U_{e_{1}+e_{2}}(y)=U_{e_{1}+e_{2}}\left(y^{\prime}\right) \prod_{i=3}^{l} U_{e_{2}-e_{i}}\left(-a_{i}^{\prime} y^{\prime}\right) \prod_{i=3}^{l} U_{e_{2}+e_{i}}\left(-b_{i}^{\prime} y^{\prime}\right) h_{1, y} s_{1}^{-\prime \prime} \tag{3.9}
\end{equation*}
$$

with a suitable $s_{1}^{-\prime \prime} \in N^{-}$, where

$$
y^{\prime}=\frac{y}{1+a_{2}^{\prime} y} \quad \text { and } \quad h_{1, y}=\Phi_{e_{1}+e_{2}}\left(\begin{array}{cc}
\frac{1}{1+a_{2}^{\prime} y} & 0 \\
0 & 1+a_{2}^{\prime} y
\end{array}\right) \in T
$$

While

$$
\begin{equation*}
\left(s_{2}^{-}\right)^{\prime} U_{e_{1}+e_{2}}\left(k_{1}+y\right)=U_{e_{1}+e_{2}}\left(k_{1, y}\right) \prod_{i=3}^{l} U_{e_{2}-e_{i}}\left(-c_{i}^{\prime} k_{1, y}\right) \prod_{i=3}^{l} U_{e_{2}+e_{i}}\left(-d_{i}^{\prime} k_{1, y}\right) h_{1, y}^{\prime} s_{2}^{-\prime \prime} \tag{3.10}
\end{equation*}
$$

with a suitable $s_{2}^{-\prime \prime} \in N^{-}$, where

$$
k_{1, y}=\frac{k_{1}+y}{1+c_{2}^{\prime}\left(k_{1}+y\right)} \quad \text { and } \quad h_{1, y}^{\prime}=\Phi_{e_{1}+e_{2}}\left(\begin{array}{cc}
\frac{1}{1+c_{2}^{\prime}\left(k_{1}+y\right)} & 0 \\
0 & 1+c_{2}^{\prime}\left(k_{1}+y\right)
\end{array}\right) \in T
$$

Thus from (3.7), (3.9), the $M$-part of $s_{1} s_{1}^{-} u$ is:

$$
\begin{aligned}
m_{1}= & \prod_{i=3}^{l} U_{-e_{2}-e_{i}}\left(a_{i} x^{\prime}\right) \prod_{i=3}^{l} U_{-e_{2}+e_{i}}\left(b_{i} x^{\prime}\right) h_{0, x} \prod_{i=3}^{l} U_{e_{2}-e_{i}}\left(-a_{i}^{\prime} y^{\prime}\right) \\
& \times \prod_{i=3}^{l} U_{e_{2}+e_{i}}\left(-b_{i}^{\prime} y^{\prime}\right) h_{1, y}
\end{aligned}
$$

While from (3.8) (3.10), the $M$-part of $s_{2}^{-} s_{2} u$ is:

$$
\begin{aligned}
m_{2}= & \prod_{i=3}^{l} U_{-e_{2}-e_{i}}\left(c_{i} k_{0, x}\right) \prod_{i=3}^{l} U_{-e_{2}+e_{i}}\left(d_{i} k_{0, x}\right) h_{0, x}^{\prime} \prod_{i=3}^{l} U_{e_{2}-e_{i}}\left(-c_{i}^{\prime} k_{1, y}\right) \\
& \quad \times \prod_{i=3}^{l} U_{e_{2}+e_{i}}\left(-d_{i}^{\prime} k_{1, y}\right) h_{1, y}^{\prime}
\end{aligned}
$$

Because both $m_{1}$ and $m_{2}$ are products of one dimensional unipotent subgroups of different root vectors in the same order, for $m_{1}=m_{2}$ to be true, then $a_{i} x^{\prime}=c_{i} k_{0, x}$ and $b_{i} x^{\prime}=d_{i} k_{0, x}$ must hold for all $i, 3 \leq i \leq l$, and for almost all $x \in F$. These equations lead to:

$$
\begin{align*}
\left(c_{i} b_{2}-a_{i} d_{2}\right) x^{2}+\left(c_{i} k_{0}+c_{i} b_{2}-a_{i} d_{2} k_{0}-a_{i}\right) x+c_{i} k_{0} & =0  \tag{3.11}\\
\left(d_{i} b_{2}-b_{i} d_{2}\right) x^{2}+\left(d_{i} k_{0}+d_{i} b_{2}-b_{i} d_{2} k_{0}-b_{i}\right) x+d_{i} k_{0} & =0 \tag{3.12}
\end{align*}
$$

For equations (3.11) and (3.12) to have infinitely many solutions, since $k_{0} \neq 0$, we must have $a_{i}=b_{i}=c_{i}=d_{i} \equiv 0$ for all $i, 3 \leq i \leq l$. Moreover, we must have $h_{0, x} h_{1, y}=h_{0, x}^{\prime} h_{1, y}^{\prime}$, which implies $h_{0, x}=h_{0, x}^{\prime}$ for almost all $x \in F$, since $h_{0, x}\left(h_{0, x}^{\prime}\right), h_{1, y}\left(h_{1, y}^{\prime}\right)$ are dual to $e_{1}-e_{2}, e_{1}+e_{2}$, respectively. So $\left(d_{2}-b_{2}\right) x+d_{2} k_{0}=0$ has infinitely many solutions in $F$, and consequently $d_{2}=b_{2}=0$.

So $s_{1}^{-}=U_{-e_{1}-e_{2}}\left(a_{2}\right), s_{2}^{-}=U_{-e_{1}-e_{2}}\left(c_{2}\right)$. And it can be easily calculated that $m_{1}=h_{1, y}, m_{2}=h_{1, y}^{\prime}$ with $a_{2}^{\prime}=a_{2}$ in $h_{1, y}$ and $c_{2}^{\prime}=c_{2}$ in $h_{1, y}^{\prime}$. Thus for $m_{1}=m_{2}$ to be true for almost all $y \in F$, we must have $\left(c_{2}-a_{2}\right) y+c_{2} k_{1}=0$. Since $s_{2}^{-} \neq$ $0, c_{2} \neq 0$, so we must have $k_{1}=0$ and $a_{2}=c_{2}$. So $s_{2}^{-} s_{2}=U_{-e_{1}-e_{2}}\left(c_{2}\right) U_{e_{1}-e_{2}}\left(k_{0}\right)=$ $U_{e_{1}-e_{2}}\left(k_{0}\right) U_{-e_{1}-e_{2}}\left(c_{2}\right)=s_{2} s_{2}^{-}=s_{1} s_{1}^{-}$.

By the uniqueness of Bruhat decomposition, $s_{1}^{-}=s_{2}^{-}$. That finishes the proof of the main theorem for the case $G$ is of type $D_{l}$.

### 3.6 Type $E_{6}$



In this case, $N$ is abelian only when $\alpha=\alpha_{1}$ or $\alpha_{6}$ by Lemma 2.1. Since $\alpha_{1}$ is symmetric to $\alpha_{6}$ on the Dynkin diagram, we need only prove the claim when $\alpha=\alpha_{1}$.

Let

$$
\begin{aligned}
& \theta_{1}=\left\{\alpha_{1} ; \alpha_{1}+\alpha_{2} ; \alpha_{1}+\alpha_{2}+\alpha_{3} ; \quad \alpha_{1}+\alpha_{2}+\alpha_{3}+\alpha_{5} ; \quad \alpha_{1}+\alpha_{2}+\alpha_{3}+\alpha_{4} ;\right. \\
& \alpha_{1}+\alpha_{2}+\alpha_{3}+\alpha_{5}+\alpha_{6} ; \quad \alpha_{1}+\alpha_{2}+\alpha_{3}+\alpha_{4}+\alpha_{5} ; \\
& \alpha_{1}+\alpha_{2}+\alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ; \quad \alpha_{1}+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5} ; \\
& \left.\alpha_{1}+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ; \quad \alpha_{1}+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+\alpha_{6} ;\right\} \\
& \theta_{2}=\left\{\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5} ; \quad \alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ;\right. \\
& \alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+\alpha_{6} ; \alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+\alpha_{6} ; \\
& \left.\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+2 \alpha_{4}+2 \alpha_{5}+\alpha_{6} .\right\}
\end{aligned}
$$

Then the positive roots in $N$ are $R^{+} \backslash \theta^{+}=\theta_{1} \cup \theta_{2}$. Notice that for each root $\beta \in$ $\theta_{1}, \beta-\alpha_{1}$ is still a root, while for $\beta \in \theta_{2}, \beta-\alpha_{1}$ is not a root. Also notice that the coefficient of $\alpha_{2}$ of roots in $\theta_{1}$ is 1 , while the coefficient of $\alpha_{2}$ of roots in $\theta_{2}$ is 2 .

Let $\beta_{1}, \beta_{2}, \ldots, \beta_{11}$ be the roots in $\theta_{1}$ according to the order listed in $\theta_{1}$, and let $\gamma_{1}, \gamma_{2}, \ldots, \gamma_{5}$ be the roots in $\theta_{2}$ accordingly. Let $\tau_{i}=\beta_{i}-\alpha_{1}, i=1, \ldots, 11 ; \nu_{i}=$ $\gamma_{i}-\alpha_{1}, i=1, \ldots, 5,\left(\nu_{i}\right.$ is not a root) .

The roots of $N$ are divided into these two sets because each element in $U_{\beta}$ with $\beta \in \theta_{1}, \beta \neq \alpha_{1}$ can be eliminated by an element in $U_{\alpha_{1}}$ and each element in $U_{\gamma_{i}}$ with $i \neq 5$ can be eliminated by an element in $U_{\gamma_{5}}$. Elements in $U_{\gamma_{i}}$ cannot be eliminated directly by elements in $U_{\alpha_{1}}$ since $\gamma_{i}-\alpha_{1}$ is not a root.

We will define an order on $R$ : suppose $\beta, \gamma \in R$ and

$$
\beta-\gamma=\sum_{i=1}^{6} c_{i} \alpha_{i}
$$

If

$$
\sum_{i=1}^{6} c_{i}>0
$$

then $\beta \succ \gamma$. If

$$
\sum_{i=1}^{6} c_{i}=0
$$

and if the first nonzero coefficient is $>0$, then $\beta \succ \gamma$, otherwise $\beta \prec \gamma$. In particular, if $\beta \in R$ is a positive root, then $\beta \succ 0$. It is easily verified that this order is well defined and we have $\beta_{i} \prec \beta_{j}$ if $1 \leq i<j \leq 11$ and $\gamma_{i} \prec \gamma_{j}$ if $1 \leq i<j \leq 5$.

Let

$$
N_{1}=\left\{\prod_{i=1}^{11} U_{\beta_{i}}\right\} \in N, \quad N_{2}=\left\{\prod_{i=1}^{5} U_{\gamma_{i}}\right\} \in N
$$

be the subgroups (because $N$ is abelian) of $N$ consisting of the unipotent subgroups of roots in $\theta_{1}, \theta_{2}$, respectively. We will prove that $N_{1}$ can be generated by $U_{\beta_{1}}=U_{\alpha_{1}}$ and $N_{2}$ can be generated by $U_{\gamma_{5}}$ under the adjoint action of $M$.

For each pair of roots $\beta, \gamma \in R$, by Lemma 2.1 we know that

$$
\begin{equation*}
U_{\gamma}(x) U_{\beta}(y) U_{\gamma}(-x)=\prod_{i, j>0} U_{i \gamma+j \beta \in R}\left(C_{\gamma, \beta, i, j} x^{i} y^{j}\right) U_{\beta}(y) \tag{3.13}
\end{equation*}
$$

Suppose the structure constants are normalized as in Lemma 2.1.

Lemma 3.6 For each $u \in N$, if $u=u_{1} u_{2}, u_{i} \in N_{i}, i=1$, 2, with $u_{1} \neq 1$. Then there exists $m \in M$ such that

$$
\operatorname{Int}(m) \circ u=\prod_{i=1}^{11} U_{\beta_{i}}\left(x_{i}^{\prime}\right) \prod_{i=1}^{5} U_{\gamma_{i}}\left(y_{i}^{\prime}\right)
$$

with $x_{1}^{\prime} \neq 0$.

Proof Suppose

$$
u=u_{1} u_{2}=\prod_{i=1}^{11} U_{\beta_{i}}\left(x_{i}\right) \prod_{i=1}^{5} U_{\gamma_{i}}\left(y_{i}\right)
$$

If $x_{1} \neq 0$, then there is nothing we need to do. Otherwise, let $k$ be the smallest $i$ such that $x_{i} \neq 0$. Notice such $i$ exists since $u_{1} \neq 1$, and by the assumption,

$$
u_{1}=\prod_{i=k}^{11} U_{\beta_{i}}\left(x_{i}\right)
$$

Let $m=U_{-\tau_{k}}(1)$. For any pair $\{i, j\}$ of positive integers, $i \beta_{k}+j\left(-\tau_{k}\right)$ is a root only when $i=j=1$, and $\beta_{k}+\left(-\tau_{k}\right)=\beta_{1}$. So we apply equation (3.13):

$$
\operatorname{Int}(m) \circ U_{\beta_{k}}\left(x_{k}\right)=U_{\beta_{1}}\left(x_{k}\right) U_{\beta_{k}}\left(x_{k}\right)
$$

since $C_{-\tau_{k}, \beta_{k}, 1,1}$ is normalized to be 1 .
For any $n>k, n \leq 11$, there is no pair $\{i, j\}$ of positive integers such that $i \beta_{n}+$ $j\left(-\tau_{k}\right)$ is a root. To verify this, we need only to check the coefficients of $\alpha_{1}$ and $\alpha_{2}$ in $i \beta_{n}+j\left(-\tau_{k}\right)$. Namely, since $N$ is abelian, the coefficient of $\alpha_{1}$ in any root in $N$ must be 1 , so $i=1$. Meanwhile the coefficient of $\alpha_{2}$ of $i \beta_{n}+j\left(-\tau_{k}\right)$ is $1-j \leq 0$, so $j$ must be 1 , too, and if this is the case, the coefficient of $\alpha_{2}$ in $\beta_{n}-\tau_{k}$ is 0 . Then $\beta_{n}-\tau_{k}=\beta_{1}$, since $\beta_{1}$ is the only root in $N$ that has coefficient of $\alpha_{2}$ equal to 0 . But $\beta_{n} \succ \beta_{k}=\beta_{1}+\tau_{k}$, this is a contradiction. So by Lemma $2.1 \operatorname{Int}(m)$ fixes $U_{\beta_{n}}$.

Also for each $n$ with $1 \leq n \leq 5, i \gamma_{n}+j\left(-\tau_{k}\right)$ can possibly be a root only when $i=j=1$. (Since $N$ is abelian, $i$ must be 1 and we can exclude the possibility $j=2$ since $\gamma_{n}+2\left(-\tau_{k}\right)$ would not be connected by just applying Lemma 3.2.) If $\gamma_{n}-\tau_{k}$ is a root, then $\gamma_{n}-\tau_{k} \succ \alpha_{1}=\beta_{1}$. So by Lemma 2.1

$$
\operatorname{Int}(m) \circ U_{\gamma_{n}} \subset \prod_{\beta \succ \beta_{1}} U_{\beta}
$$

With these facts,

$$
\operatorname{Int}(m) \circ u \in U_{\beta_{1}}\left(x_{k}\right) \prod_{\beta \succ \beta_{1}} U_{\beta}
$$

Lemma 3.7 For each $u_{2} \neq 1 \in N_{2}$, there exists an $m \in M$ such that $\operatorname{Int}(m)$ fixes $U_{\beta_{1}}$ and

$$
\operatorname{Int}(m) \circ u_{2}=\prod_{i=1}^{5} U_{\gamma_{i}}\left(y_{i}\right), \quad \text { with } y_{5} \neq 0
$$

Proof Suppose

$$
u_{2}=\prod_{i=1}^{5} U_{\gamma_{i}}\left(x_{i}\right)
$$

If $x_{5} \neq 0$, then nothing needs to be done. Otherwise, let $k$ be the smallest $i$ such that $x_{i} \neq 0$. So $x_{i} \neq 0$ only when $k \leq i \leq 4$. Let $\gamma=\gamma_{5}-\gamma_{k}$, and $m=U_{\gamma}(1)$.

For each pair $\{i, j\}$ of positive integers, $i \gamma+j \gamma_{k}$ can be a root only when $i=j=1$, since otherwise $i \gamma+j \gamma_{k} \succ \gamma_{5}$, and $\gamma_{5}$ is the longest element in $R$ such that its $\alpha_{1}$ part is nonzero. Moreover, in this case $\gamma+\gamma_{k}=\gamma_{5}$. So by applying Lemma 2.1, we have: $\operatorname{Int}(m) \circ U_{\gamma_{k}}\left(x_{k}\right)=U_{\gamma_{k}}\left(x_{k}\right) U_{\gamma_{5}}\left(C_{\gamma, \gamma_{k}, 1,1} x_{k}\right)$, where $C_{\gamma, \gamma_{k}, 1,1}$ is a structure constant, so is nonzero.

For all other $q$ with $k<q \leq 4, i \gamma+j \gamma_{q}$ could not be a root since $i \gamma+j \gamma_{q} \succ \gamma_{5}$ for any positive integers $i, j$. So $\operatorname{Int}(m)$ fixes all these $U_{\gamma_{q}}$.

With these two facts, it is easily calculated that $\operatorname{Int}(m) \circ u=u U_{\gamma_{5}}\left(C_{\gamma, \gamma_{k}, 1,1} x_{k}\right)$. Now, set $y_{5}=C_{\gamma, \gamma_{k}, 1,1} x_{k}$. Then $y_{5} \neq 0$ as we have shown. Because $\gamma \subset \operatorname{span}\left\{\alpha_{3}, \alpha_{4}, \alpha_{5}\right.$, $\left.\alpha_{6}\right\}$, for each pair $\{i, j\}$ of positive integers, $i \gamma+j \beta_{1}$ cannot be a root by Lemma 3.2. So $\operatorname{Int}(m)$ fixes $U_{\beta_{1}}$ by Lemma 2.1.

Theorem 3.8 (Gaussian Elimination) For each $u \neq 1 \in N$, there exists $m \in M$ such that $\operatorname{Int}(m) \circ u=U_{\beta_{1}}\left(k_{0}\right) U_{\gamma_{5}}\left(k_{1}\right)$.

Proof We can write $u$ as

$$
u=\prod_{\beta_{i} \in \theta_{1}} U_{\beta_{i}}\left(x_{i}\right) \prod_{\gamma_{i} \in \theta_{2}} U_{\gamma_{i}}\left(x_{i}^{\prime}\right)=u_{1} u_{2}, \quad u_{1} \in N_{1}, u_{2} \in N_{2} .
$$

If $u_{1}=1$, then just set $m_{1}=1$. If $u_{1} \neq 1$, by applying Lemma3.6 and a suitable $\operatorname{Int}\left(m^{\prime}\right)$, if necessary, we can assume $x_{1} \neq 0$.

Let

$$
m_{1}=\prod_{i=2}^{11} U_{\tau_{i}}\left(\frac{x_{i}}{x_{1}}\right)
$$

Then

$$
\operatorname{Int}\left(m_{1}\right) \circ u=\left[\prod_{\beta_{i} \in \theta_{1}} \operatorname{Int}\left(m_{1}\right) \circ U_{\beta_{i}}\left(x_{i}\right)\right] \cdot\left[\prod_{\gamma_{i} \in \theta_{2}} \operatorname{Int}\left(m_{1}\right) \circ U_{\gamma_{i}}\left(x_{i}^{\prime}\right)\right]
$$

For each fixed $k$, with $2 \leq k \leq 11, i \beta_{1}+j \tau_{k}$ is a root for $i, j>0$ only when $i=j=1$, and in this case $\beta_{1}+\tau_{k}=\beta_{k}$. So by applying Lemma 2.1,

$$
\operatorname{Int}\left(U_{\tau_{k}}\left(\frac{x_{k}}{x_{1}}\right)\right) \circ U_{\beta_{1}}\left(x_{1}\right)=U_{\beta_{1}}\left(x_{1}\right) \cdot U_{\beta_{k}}\left(-x_{k}\right)
$$

For each $q, 2 \leq q \leq 11, q \neq k$, and each pair of positive integers $\{i, j\}, i \beta_{q}+j \tau_{k}$ can possibly be a root only when $i=j=1$. And in this case $\beta_{q}+\tau_{k} \in \theta_{2}$ if it is a root, since the coefficient of $\alpha_{2}$ in $\beta_{q}+\tau_{k}$ is 2 . So

$$
\operatorname{Int}\left(U_{\tau_{k}}\left(\frac{x_{k}}{x_{1}}\right)\right) \circ U_{\beta_{q}}\left(x_{q}\right)=U_{\beta_{q}}\left(x_{q}\right) \cdot n_{q} \quad \text { with } n_{q} \in N_{2}
$$

For each pair of positive integers $\{i, j\}$, none of $i \beta_{k}+j \tau_{k}$ can be a root. So also by Lemma 2.1,

$$
\operatorname{Int}\left(U_{\tau_{k}}\left(\frac{x_{k}}{x_{1}}\right)\right) \quad \text { fixes } U_{\beta_{k}} .
$$

With these facts, one can conclude from

$$
\operatorname{Int}\left(m_{1}\right)=\prod_{i=2}^{11} \operatorname{Int}\left(U_{\tau_{i}}\left(\frac{x_{i}}{x_{1}}\right)\right)
$$

that

$$
\begin{aligned}
& \operatorname{Int}\left(m_{1}\right) \circ U_{\beta_{1}}\left(x_{1}\right)=U_{\beta_{1}}\left(x_{1}\right) \prod_{i=2}^{11} U_{\beta_{i}}\left(-x_{i}\right) \cdot n_{1} \quad \text { with } n_{1} \in N_{2} \\
& \operatorname{Int}\left(m_{1}\right) \circ U_{\beta_{i}}\left(x_{i}\right)=U_{\beta_{i}}\left(x_{i}\right) \cdot n_{i}^{\prime} \quad \text { with } n_{i}^{\prime} \in N_{2}, \forall i, 2 \leq i \leq 11
\end{aligned}
$$

By the last two equations, one can get

$$
\operatorname{Int}\left(m_{1}\right) \circ\left(u_{1}\right)=U_{\beta_{1}}\left(x_{1}\right) \cdot n^{\prime} \quad \text { where } n^{\prime}=n_{1} \cdot \prod_{i=2}^{11} n_{i}^{\prime} \in N_{2}
$$

For each $\gamma \in \theta_{2}$, none of $i \tau_{k}+j \gamma$ is a root for any pair of positive integers $\{i, j\}$, since in the decomposition of $i \tau_{k}+j \gamma_{i}$ as a summation of simple roots, the coefficient of $\alpha_{2}$ will be $i+2 j \geq 3$, which is not possible. So $\operatorname{Int}\left(m_{1}\right) \circ u_{2}=u_{2}$.

Now we have $\operatorname{Int}\left(m_{1}\right) \circ u=\operatorname{Int}\left(m_{1}\right) \circ\left(u_{1} u_{2}\right)=U_{\beta_{1}}\left(x_{1}\right) n^{\prime} u_{2}$. Suppose

$$
n^{\prime} u_{2}=\prod_{i=1}^{5} U_{\gamma_{i}}\left(y_{i}^{\prime}\right)
$$

If $n^{\prime} u_{2}=1$, i.e., $y_{i}=0$ for $1 \leq i \leq 5$, then we are done. Otherwise, let $m_{2}$ be the element in $m$ that comes from Lemma 3.7. Then

$$
\operatorname{Int}\left(m_{2} m_{1}\right) \circ u=U_{\beta_{1}} \cdot \prod_{i=1}^{5} U_{\gamma_{i}}\left(y_{i}\right), \quad \text { with } y_{5} \neq 0
$$

Now let

$$
m_{3}=\prod_{i=1}^{4} U_{\gamma_{i}-\gamma_{5}}\left(-\frac{y_{i}}{y_{5}}\right)
$$

Then by Lemma 2.1, for any fixed $i$,

$$
\operatorname{Int}\left(U_{\gamma_{i}-\gamma_{5}}\left(-\frac{y_{i}}{y_{5}}\right)\right) \circ U_{\gamma_{5}}\left(y_{5}\right)=U_{\gamma_{5}}\left(y_{5}\right) \cdot U_{\gamma_{i}}\left(-y_{i}\right)
$$

It can be easily shown, by checking the coefficients of $\alpha$ and $\alpha_{4}$, that for any pair $\{j, k\}$ of positive integers, and any $q$, with $1 \leq q \leq 4$, none of $j\left(\gamma_{i}-\gamma_{5}\right)+k \gamma_{q}$ can be a root. (Namely, for $j\left(\gamma_{i}-\gamma_{5}\right)+k \gamma_{q}$ to be a root, $k$ must be 1 since the coefficient of $\alpha$ in $j\left(\gamma_{i}-\gamma_{5}\right)+k \gamma_{q}$ is $k$. Then the coefficient of $\alpha_{2}$ in $j\left(\gamma_{i}-\gamma_{5}\right)+k \gamma_{q}$ is 2 , so $j\left(\gamma_{i}-\gamma_{5}\right)+k \gamma_{q} \in \theta_{2}$ if it is a root. Then the coefficient of $\alpha_{4}$ in $j\left(\gamma_{i}-\gamma_{5}\right)+k \gamma_{q}$ is
$1-j \leq 0$ which is not possible since every root in $\theta_{2}$ has its coefficient of $\alpha_{4}$ equal to 1.) So again by Lemma 2.1, $\operatorname{Int}\left(U_{\gamma_{i}-\gamma_{5}}\right)$ fixes all other $U_{\gamma_{q}}\left(y_{q}\right)$. Thus

$$
\begin{aligned}
& \operatorname{Int}\left(m_{3}\right) \circ U_{\gamma_{5}}\left(y_{5}\right)=U_{\gamma_{5}}\left(y_{5}\right) \prod_{i=1}^{4} U_{\gamma_{i}}\left(-y_{i}\right) \\
& \operatorname{Int}\left(m_{3}\right) \circ\left(\prod_{i=1}^{4} U_{\gamma_{i}}\left(y_{i}\right)\right)=\prod_{i=1}^{4} U_{\gamma_{i}}\left(y_{i}\right)
\end{aligned}
$$

So

$$
\operatorname{Int}\left(m_{3}\right) \circ\left(\prod_{i=1}^{5} U_{\gamma_{i}}\left(y_{i}\right)\right)=U_{\gamma_{5}}\left(y_{5}\right)
$$

Moreover, for each $i, \operatorname{Int}\left(U_{\gamma_{i}-\gamma_{5}}\right)$ fixes $U_{\beta_{1}}$ since, from the proof of Lemma 3.7, $\gamma_{i}-\gamma_{5} \subset \operatorname{span}\left\{\alpha_{3}, \alpha_{4}, \alpha_{5}, \alpha_{6}\right\}$. Consequently, $\operatorname{Int}\left(m_{3}\right)$ fixes $U_{\beta_{1}}\left(x_{1}\right)$. Now let $m=m_{3} m_{2} m_{1}$. Then $\operatorname{Int}(m) \circ u=U_{\beta_{1}}\left(x_{1}\right) U_{\gamma_{5}}\left(y_{5}\right)$. Setting $k_{0}=x_{1}, k_{1}=y_{5}$ proves the theorem.

Returning to equation (2.5), by the above lemma and applying $\operatorname{Int}(m)$ on both sides, we can assume $s_{2}=U_{\beta_{1}}\left(k_{0}\right) U_{\gamma_{5}}\left(k_{1}\right)$. Since without loss of generality we can always assume $s_{2} \neq 1$ (otherwise nothing needs to be proved), we assume $k_{0} \neq 0$.

Now suppose

$$
\begin{aligned}
& s_{1}^{-}=\prod_{i=1}^{11} U_{-\beta_{i}}\left(a_{i}\right) \cdot \prod_{i=1}^{5} U_{-\gamma_{i}}\left(b_{i}\right) \\
& s_{2}^{-}=\prod_{i=1}^{11} U_{-\beta_{i}}\left(c_{i}\right) \cdot \prod_{i=1}^{5} U_{-\gamma_{i}}\left(d_{i}\right)
\end{aligned}
$$

Multiply both sides of (2.5) by $u=U_{\beta_{1}}(x) U_{\gamma_{5}}(y)$ on the right, where $x, y$ are variables in $F$. We will decompose both $s_{1} s_{1}^{-} u$ and $s_{2}^{-} s_{2} u$ into $P N^{-}$form, and compare their $M$ parts.

First for $s_{1} s_{1}^{-} U_{\beta_{1}}(x)$, we have:

$$
\begin{equation*}
U_{-\beta_{1}}\left(a_{1}\right) U_{\beta_{1}}(x)=U_{\beta_{1}}\left(x^{\prime}\right) h_{x} U_{-\beta_{1}}\left(\frac{a_{1}}{1+a_{1} x}\right) \tag{3.14}
\end{equation*}
$$

where

$$
x^{\prime}=\frac{x}{1+a_{1} x}, \quad \text { and } \quad h_{x}=\Phi_{\beta_{1}}\left(\begin{array}{cc}
\frac{1}{1+a_{1} x} & 0 \\
0 & 1+a_{1} x
\end{array}\right) \in T
$$

For each $k$ with $2 \leq k \leq 11$, by Lemma 2.1 we have

$$
\begin{equation*}
U_{-\beta_{k}}\left(a_{k}\right) U_{\beta_{1}}\left(x^{\prime}\right)=U_{\beta_{1}}\left(x^{\prime}\right) U_{-\tau_{k}}\left(-a_{k} x^{\prime}\right) U_{-\beta_{k}}\left(a_{k}\right) \tag{3.15}
\end{equation*}
$$

For any $k$ with $1 \leq k \leq 5$, and any pair $\{i, j\}$ of positive integers, none of $i \beta_{1}+j\left(-\gamma_{k}\right)$ can be a root. So by Lemma 2.1, $U_{\beta_{1}}$ commutes with $U_{-\gamma_{k}}$ for all $k$.

Since $N^{-}$is normal in $P^{-}=M N^{-}$, from equations (3.14), (3.15) and the above fact, the $P N^{-}$decomposition of $s_{1}^{-} U_{\beta_{1}}(x)$ is as follows:

$$
\begin{equation*}
s_{1}^{-} U_{\beta_{1}}(x)=U_{\beta_{1}}\left(x^{\prime}\right)\left[\prod_{i=2}^{11} U_{-\tau_{i}}\left(-a_{i} x^{\prime}\right)\right] h_{x} \cdot\left(s_{1}^{-}\right)^{\prime} \tag{3.16}
\end{equation*}
$$

with a suitable $\left(s_{1}^{-}\right)^{\prime} \in N^{-}$. Then for $s_{2}^{-} s_{2} u=s_{2}^{-} U_{\beta_{1}}\left(k_{0}+x\right) U_{\gamma_{5}}\left(k_{1}+y\right)$. Similarly, the $P N^{-}$decomposition of $s_{2}^{-} U_{\beta_{1}}\left(k_{0}+x\right)$ is:

$$
\begin{equation*}
s_{2}^{-} U_{\beta_{1}}\left(k_{0}+x\right)=U_{\beta_{1}}\left(k_{x}\right)\left[\prod_{i=2}^{11} U_{-\tau_{i}}\left(-k_{x} c_{i}\right)\right] h_{x}^{\prime}\left(s_{2}^{-}\right)^{\prime} \tag{3.17}
\end{equation*}
$$

with a suitable $\left(s_{2}^{-}\right)^{\prime} \in N^{-}$, where

$$
k_{x}=\frac{k_{0}+x}{1+c_{1}\left(k_{0}+x\right)}, \quad \text { and } \quad h_{x}^{\prime}=\Phi_{\beta_{1}}\left(\begin{array}{cc}
\frac{1}{1+c_{1}\left(k_{0}+x\right)} & 0 \\
0 & 1+c_{1}\left(k_{0}+x\right)
\end{array}\right) \in T
$$

For convenience of notation, we will set $U_{\nu_{i}} \equiv 1$ if $\nu_{i}$ is not a root. Suppose

$$
\left(s_{1}^{-}\right)^{\prime}=\prod_{i=1}^{11} U_{-\beta_{i}}\left(a_{i}^{\prime}\right) \cdot \prod_{i=1}^{5} U_{-\gamma_{i}}\left(b_{i}^{\prime}\right), \quad\left(s_{2}^{-}\right)^{\prime}=\prod_{i=1}^{11} U_{-\beta_{i}}\left(c_{i}^{\prime}\right) \cdot \prod_{i=1}^{5} U_{-\gamma_{i}}\left(d_{i}^{\prime}\right)
$$

Then with a similar discussion on roots and applying Lemma 2.1, following a similar process of calculation, we get:

$$
\begin{equation*}
\left(s_{1}^{-}\right)^{\prime} U_{\gamma_{5}}(y)=U_{\gamma_{5}}\left(y^{\prime}\right)\left[\prod_{i=1}^{11} U_{\gamma_{5}-\beta_{i}}\left(a_{i}^{\prime} y^{\prime}\right)\right]\left[\prod_{i=1}^{4} U_{\gamma_{5}-\gamma_{i}}\left(b_{i}^{\prime} y^{\prime}\right)\right] h_{y}\left(s_{1}^{-}\right)^{\prime \prime} \tag{3.18}
\end{equation*}
$$

with a suitable $\left(s_{1}^{-}\right)^{\prime \prime} \in N^{-}$, where

$$
y^{\prime}=\frac{y}{1+b_{5}^{\prime} y} \quad \text { and } \quad h_{y}=\Phi_{\gamma_{5}}\left(\begin{array}{cc}
\frac{1}{1+b_{5}^{\prime} y} & 0 \\
0 & 1+b_{5}^{\prime} y
\end{array}\right) \in T
$$

Meanwhile,

$$
\begin{equation*}
\left(s_{2}^{-}\right)^{\prime} U_{\gamma_{5}}\left(k_{1}+y\right)=U_{\gamma_{5}}\left(k_{y}\right)\left[\prod_{i=1}^{11} U_{\gamma_{5}-\beta_{i}}\left(k_{y} c_{i}^{\prime}\right)\right]\left[\prod_{i=1}^{4} U_{\gamma_{5}-\gamma_{i}}\left(k_{y} d_{i}^{\prime}\right)\right] h_{y}^{\prime}\left(s_{2}^{-}\right)^{\prime \prime} \tag{3.19}
\end{equation*}
$$

with a suitable $\left(s_{2}^{-}\right)^{\prime \prime} \in N^{-}$, where

$$
k_{y}=\frac{k_{1}+y}{1+d_{5}^{\prime}\left(k_{1}+y\right)} \quad \text { and } \quad h_{y}^{\prime}=\Phi_{\gamma_{5}}\left(\begin{array}{cc}
\frac{1}{1+d_{5}^{\prime}\left(k_{1}+y\right)} & 0 \\
0 & 1+d_{5}^{\prime}\left(k_{1}+y\right)
\end{array}\right) \in T
$$

Thus, from equations (3.16) and (3.18), the $M$-part of $s_{1} s_{1}^{-} u$ is:

$$
M_{1}(x, y)=\left[\prod_{i=2}^{11} U_{-\tau_{i}}\left(-a_{i} x^{\prime}\right)\right] h_{x}\left[\prod_{i=1}^{11} U_{\gamma_{5}-\beta_{i}}\left(a_{i}^{\prime} y^{\prime}\right)\right]\left[\prod_{i=1}^{4} U_{\gamma_{5}-\gamma_{i}}\left(b_{i}^{\prime} y^{\prime}\right)\right] h_{y}
$$

While from equation (3.17) and (3.19), the $M$-part of $s_{2}^{-} s_{2} u$ is:

$$
M_{2}(x, y)=\left[\prod_{i=2}^{11} U_{-\tau_{i}}\left(-k_{x} c_{i}\right)\right] h_{x}^{\prime}\left[\prod_{i=1}^{11} U_{\gamma_{5}-\beta_{i}}\left(k_{y} c_{i}^{\prime}\right)\right]\left[\prod_{i=1}^{4} U_{\gamma_{5}-\gamma_{i}}\left(k_{y} d_{i}^{\prime}\right)\right] h_{y}^{\prime}
$$

Notice that all $-\tau_{i}$ are distinct negative roots while all $\gamma_{5}-\beta_{i}$ and $\gamma_{5}-\gamma_{i}$ are distinct positive roots (if they are roots). For $M_{1}(x, y)=M_{2}(x, y)$, the unipotent groups of the corresponding root vector must be equal, and their $T$ parts must be equal as well, as in the previous cases.

So we have $a_{i} x^{\prime}=k_{x} c_{i}$, for all $i, 2 \leq i \leq 11$, and almost all $x \in F$. Moreover, $h_{x}=h_{x}^{\prime}$ since $h_{x}\left(h_{x}^{\prime}\right), h_{y}\left(h_{y}^{\prime}\right)$ are dual to $\beta_{1}, \gamma_{5}$, respectively. As an analog of the proof in the $B_{l}\left(D_{l}\right)$ case, we get $a_{i}=c_{i} \equiv 0, \forall 1 \leq i \leq 11$. Thus from equation (3.16) and (3.17),

$$
s_{1}^{-}=s_{1}^{-\prime}=\prod_{i=1}^{5} U_{-\gamma_{i}}\left(b_{i}\right), \quad s_{2}^{-}=s_{2}^{-\prime}=\prod_{i=1}^{5} U_{-\gamma_{i}}\left(d_{i}\right)
$$

and from equations (3.18), (3.19),

$$
M_{1}(x, y)=\left[\prod_{i=1}^{4} U_{\gamma_{5}-\gamma_{i}}\left(b_{i} y^{\prime}\right)\right] h_{y}, \quad M_{2}(x, y)=\left[\prod_{i=1}^{4} U_{\gamma_{5}-\gamma_{i}}\left(d_{i} k_{y}\right)\right] h_{y}^{\prime}
$$

Since $s_{2}^{-} \neq 1$, there is one $i, 1 \leq i \leq 5$, such that $d_{i} \neq 0$. Together with the fact that $M_{1}(x, y)=M_{2}(x, y)$ for almost all $x, y \in F$, following the previous proofs, we can get $k_{1}=0$. So

$$
s_{2}^{-} s_{2}=\left[\prod_{i=1}^{5} U_{-\gamma_{i}}\left(d_{i}\right)\right] U_{\beta_{1}}\left(k_{0}\right)=U_{\beta_{1}}\left(k_{0}\right)\left[\prod_{i=1}^{5} U_{-\gamma_{i}}\left(d_{i}\right)\right]=s_{2} s_{2}^{-}=s_{1} s_{1}^{-}
$$

By the uniqueness of Bruhat decomposition, it must have $s_{1}^{-}=s_{2}^{-}$.
If at the beginning of this proof, we assume $k_{1} \neq 0$ instead of assuming $k_{0} \neq 0$, the proof will be similar.

### 3.7 Type $E_{7}$



The longest root in this case is $2 \alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+2 \alpha_{4}+3 \alpha_{5}+2 \alpha_{6}+\alpha_{7}$, by Lemma 3.1; $N$ is abelian only when $\alpha=\alpha_{7}$.

Let

$$
\begin{aligned}
& \theta_{1}=\left\{\alpha ; \alpha+\alpha_{6} ; \quad \alpha+\alpha_{5}+\alpha_{6} ; \quad \alpha+\alpha_{3}+\alpha_{5}+\alpha_{6} ; \quad \alpha+\alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ;\right. \\
& \quad \alpha+\alpha_{2}+\alpha_{3}+\alpha_{5}+\alpha_{6} ; \quad \alpha+\alpha_{2}+\alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ; \\
& \quad \alpha+\alpha_{1}+\alpha_{2}+\alpha_{3}+\alpha_{5}+\alpha_{6} ; \quad \alpha+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ; \\
& \quad \alpha+\alpha_{1}+\alpha_{2}+\alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ; \quad \alpha+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+\alpha_{6} ; \\
& \quad \alpha+\alpha_{1}+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ; \quad \alpha+\alpha_{1}+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+\alpha_{6} ; \\
& \quad \alpha+\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} ; \quad \alpha+\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+\alpha_{6} ; \\
& \left.\quad \alpha+\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+\alpha_{6} ; \quad \alpha+\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+2 \alpha_{4}+2 \alpha_{5}+\alpha_{6} ;\right\} \\
& \theta_{2}=\left\{\alpha+2 \alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+2 \alpha_{4}+3 \alpha_{5}+2 \alpha_{6} ;\right. \\
& \\
& \alpha+\alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+2 \alpha_{4}+3 \alpha_{5}+2 \alpha_{6} ; \quad \alpha+\alpha_{1}+2 \alpha_{2}+4 \alpha_{3}+2 \alpha_{4}+3 \alpha_{5}+2 \alpha_{6} ; \\
& \alpha+\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+2 \alpha_{4}+3 \alpha_{5}+2 \alpha_{6} ; \quad \alpha+\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+\alpha_{4}+3 \alpha_{5}+2 \alpha_{6} ; \\
& \alpha+\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+2 \alpha_{4}+2 \alpha_{5}+2 \alpha_{6} ; \quad \alpha+\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+2 \alpha_{6} ; \\
& \alpha+\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+2 \alpha_{6} ; \quad \alpha+\alpha_{1}+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+2 \alpha_{6} ; \\
& \left.\alpha+\alpha_{2}+2 \alpha_{3}+\alpha_{4}+2 \alpha_{5}+2 \alpha_{6}\right\} .
\end{aligned}
$$

Then the positive roots in $N$ are $R^{+} \backslash \theta^{+}=\theta_{1} \cup \theta_{2}$.
Let $\beta_{1}, \beta_{2}, \ldots, \beta_{17}$ denote the roots in $\theta_{1}$ as the order listed, $\gamma_{1}, \gamma_{2}, \ldots, \gamma_{10}$ denote the roots in $\theta_{2}$ similarly. For any $\beta \in \theta_{1}, \beta-\alpha$ is a root (as is $E_{6}$ ); for $i=2, \ldots, 9, \gamma_{1}-$ $\gamma_{i}$ is a root while $\gamma_{1}-\gamma_{10}$ is not; for each $i, 1 \leq i \leq 10, \gamma_{i}-\beta_{1}$ is not a root. Notice for each root in $\theta_{1}$, the coefficient of $\alpha_{6}$ is 1 , and for each root in $\theta_{2}$, the coefficient of $\alpha_{6}$ is 2 .

We will define an order on $R$ : suppose $\beta, \gamma \in R$ and

$$
\beta-\gamma=\sum_{i=1}^{7} c_{i} \alpha_{i}
$$

If

$$
\sum_{i=1}^{7} c_{i}>0
$$

then $\beta \succ \gamma$; if

$$
\sum_{i=1}^{7} c_{i}=0
$$

and if the first nonzero coefficient is $>0$, then $\beta \succ \gamma$, otherwise $\beta \prec \gamma$. In particular, if $\beta \in R$ is a positive root, then $\beta \succ 0$. It is easily verified that this order is well defined and we have $\beta_{i} \prec \beta_{j}$ if $1 \leq i<j \leq 17$ and $\gamma_{i} \succ \gamma_{j}$ if $1 \leq i<j \leq 10$.

Suppose the root vectors are so chosen that the structure constants are normalized as in Lemma 2.1. Let

$$
N_{1}=\left\{\prod_{i=1}^{17} U_{\beta_{i}}\right\} \subset N, \quad N_{2}=\left\{\prod_{i=1}^{10} U_{\gamma_{i}}\right\} \subset N
$$

Every element of $u \in N$ can be written as $u=u_{1} u_{2}$ with $u_{i} \in N_{i}, i=1,2$. And we can similarly define $N_{1}^{-}, N_{2}^{-}$as subgroups of $N^{-}$. The roots of $N$ are divided into these two sets because, as we will prove, $U_{\beta_{1}}$ generates $N_{1}$ and $U_{\gamma_{1}}$ generates $N_{2}$ under the adjoint action of $M$. Each element in $U_{\gamma_{i}}$, with $1 \leq i \leq 10$, cannot be eliminated directly by an element in $U_{\beta_{1}}$ since $\gamma_{i}-\beta_{1}$ is not a root.

Lemma 3.9 For each $u \in N$, if $u=u_{1} u_{2}, u_{i} \in N_{i}, i=1$, 2, with $u_{1} \neq 1$. Then there exists $m \in M$, such that

$$
\operatorname{Int}(m) \circ u=\left\{\prod_{i=1}^{17} U_{\beta_{i}}\left(x_{i}^{\prime}\right)\right\}\left\{\prod_{i=1}^{10} U_{\gamma_{i}}\left(y_{i}^{\prime}\right)\right\} \quad \text { with } x_{1}^{\prime} \neq 0
$$

Proof This is analogous to Lemma 3.6, since for each $i, 2 \leq i \leq 17, \beta_{i}-\beta_{1}$ is a root, the proof is almost the same as of the proof for Lemma 3.6. The indices are the only changes.

Lemma 3.10 If

$$
u_{2}=\left\{\prod_{i=1}^{10} U_{\gamma_{i}}\left(x_{i}\right)\right\} \subset N_{2}
$$

and $u_{2} \neq 1$, we can find an $m \in M$ such that $\operatorname{Int}(m) \circ u_{2}=U_{\gamma_{1}}\left(a_{2}\right)$ with $a_{2} \neq 0$ and $\operatorname{Int}(m)$ fixes every element in $U_{\beta_{1}}$.

Proof First we prove the following claim:
Claim There is $m_{1} \in M$, such that

$$
\operatorname{Int}\left(m_{1}\right) \circ u_{2}=\prod_{i=1}^{10} U_{\gamma_{i}}\left(x_{i}^{\prime}\right) \quad \text { with } x_{1}^{\prime} \neq 0, x_{2}^{\prime} \neq 0
$$

(This claim is needed because $\gamma_{1}-\gamma_{10}$ is not a root, and $U_{\gamma_{10}}$ cannot be eliminated directly through $U_{\gamma_{1}}$. So we use $U_{\gamma_{2}}$ to eliminate it.)

Let $k$ be the smallest positive integer such that $x_{k} \neq 0$. If $k=1$, i.e., $x_{1} \neq 0$. And if $x_{2} \neq 0$, then the claim is trivial.

Case $k=1, x_{2}=0$ : Let $m_{1}=U_{\gamma_{2}-\gamma_{1}}(1)$. For any $i, 3 \leq i \leq 10$, by Lemma 2.1,

$$
\operatorname{Int}\left(U_{\gamma_{2}-\gamma_{1}}(1)\right) \circ U_{\gamma_{i}}\left(x_{i}\right)=\left(\prod_{\substack{k, n>0 \\ k\left(\gamma_{2}-\gamma_{1}\right)+n \gamma_{i} \in R}} U_{k\left(\gamma_{2}-\gamma_{1}\right)+n \gamma_{i}}\left(C_{\gamma_{2}-\gamma_{1}, \gamma_{i}, k, n} x_{i}^{n}\right)\right) \cdot U_{\gamma_{i}}\left(x_{i}\right),
$$

where the $C_{\gamma_{2}-\gamma_{1}, \gamma_{i}, k, n}$ 's are structure constants.
Since $\gamma_{2}-\gamma_{1} \in \operatorname{span}\left\{\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4}, \alpha_{5}\right\}$, for any pair of positive integers $\{k, n\}$, the coefficient of $\alpha_{6}$ in $k\left(\gamma_{2}-\gamma_{1}\right)+n \gamma_{i}$ is $2 n$. For it to be a root, $n$ must be 1 . Moreover, if this is the case, then $k\left(\gamma_{2}-\gamma_{1}\right)+\gamma_{i} \in \theta_{2}$.

Since $\gamma_{2}-\gamma_{1} \prec 0, k\left(\gamma_{2}-\gamma_{1}\right)+\gamma_{i} \prec \gamma_{i}$. So

$$
\operatorname{Int}\left(m_{1}\right) \circ U_{\gamma_{i}}\left(x_{i}\right) \subset \prod_{j \geq i} U_{\gamma_{j}}
$$

consequently,

$$
\operatorname{Int}\left(m_{1}\right) \circ\left(\prod_{i=3}^{10} U_{\gamma_{i}}\left(x_{i}\right)\right) \subset \prod_{i=3}^{10} U_{\gamma_{i}}
$$

And by Lemma 2.1, $\operatorname{Int}\left(m_{1}\right) \circ U_{\gamma_{1}}\left(x_{1}\right)=U_{\gamma_{1}}\left(x_{1}\right) U_{\gamma_{2}}\left(x_{1}\right)$. Therefore, $\operatorname{Int}\left(m_{1}\right) \circ u=$ $U_{\gamma_{1}}\left(x_{1}\right) U_{\gamma_{2}}\left(x_{1}\right) \cdot u^{\prime}$ with

$$
u^{\prime} \in \prod_{i=3}^{10} U_{\gamma_{i}}
$$

Set $x_{1}^{\prime}=x_{2}^{\prime}=x_{1}$, and the claim is proved.
Case $k=2$ : Let $m_{1}=U_{\gamma_{1}-\gamma_{2}}(1)=U_{\alpha_{1}}(1)$. For each $i, 3 \leq i \leq 10$, and each pair $\{k, n\}$ of positive integers, the coefficient of $\alpha_{1}$ in $k \alpha_{1}+n \gamma_{i}$ is $k+n$. So for $k \alpha_{1}+n \gamma_{i}$ to be a root, we must have $k=n=1$. But it is easily checked that $\alpha_{1}+\gamma_{i}$ is not a root when $i \geq 3$. So by Lemma 2.1, $\operatorname{Int}\left(U_{\alpha_{1}}(1)\right) \circ U_{\gamma_{i}}\left(x_{i}\right)=U_{\gamma_{i}}\left(x_{i}\right)$. Also for any pair $\{k, n\}$ of positive integers, $k\left(\gamma_{1}-\gamma_{2}\right)+n \gamma_{2}$ can be a root only when $k=n=1$. So by applying Lemma 2.1, $\operatorname{Int}\left(U_{\alpha_{1}}(1)\right) \circ U_{\gamma_{2}}\left(x_{2}\right)=U_{\gamma_{1}}\left(x_{2}\right) U_{\gamma_{2}}\left(x_{2}\right)$, with $x_{2} \neq 0$. Then

$$
\operatorname{Int}\left(U_{\alpha_{1}}(1)\right) \circ u=U_{\gamma_{1}}\left(x_{2}\right) U_{\gamma_{2}}\left(x_{2}\right)\left[\prod_{i=3}^{10} U_{\gamma_{i}}\left(x_{i}\right)\right]
$$

Setting $x_{1}^{\prime}=x_{2}$ will prove our claim.
Case $3 \leq k<10$ : Let $m_{1}=U_{\gamma_{1}-\gamma_{k}}(1) U_{\gamma_{2}-\gamma_{k}}(1)$, with a similar discussion as the second case, but this time take the coefficients of $\alpha_{1}$ and $\alpha_{2}$ into account. We can figure out that the $U_{\gamma_{1}} U_{\gamma_{2}}$ part of $\operatorname{Int}\left(m_{1}\right) \circ u$ is $U_{\gamma_{1}}\left(x_{k}\right) U_{\gamma_{2}}\left(x_{k}\right)$.

Case $k=10$ : This case is handled separately because $\gamma_{1}-\gamma_{10}$ is not a root. Let $m_{1}=U_{\gamma_{2}-\gamma_{10}}(1)$, then $\operatorname{Int}\left(m_{1}\right) \circ u=\operatorname{Int}\left(m_{1}\right) \circ U_{\gamma_{10}}\left(x_{10}\right)=U_{\gamma_{2}}\left(x_{10}\right) U_{\gamma_{10}}\left(x_{10}\right)$ by Lemma 2.1, since for any positive integers $k$ and $n, k\left(\gamma_{2}-\gamma_{10}\right)+n \gamma_{10}$ is a root only when $k=n=1$. Now it will fall into the second case which has already been proved.

Now

$$
\operatorname{Int}\left(m_{1}\right) \circ u=\prod_{i=1}^{10} U_{\gamma_{i}}\left(x_{i}^{\prime}\right) \quad \text { with } x_{1}^{\prime} \neq 0, x_{2}^{\prime} \neq 0
$$

Let

$$
m_{2}=U_{\gamma_{10}-\gamma_{2}}\left(-\frac{x_{10}^{\prime}}{x_{2}^{\prime}}\right)
$$

It can be checked for any $i \geq 3$, and any pair of positive integers $\{k, n\}$, that $k\left(\gamma_{10}-\gamma_{2}\right)+n \gamma_{i}$ is not a root. So $\operatorname{Int}\left(m_{2}\right)$ fixes all $U_{\gamma_{i}}$.

For any pair of positive integers $\{k, n\}, k \gamma_{1}+n\left(\gamma_{10}-\gamma_{2}\right)$ or $k \gamma_{2}+n\left(\gamma_{10}-\gamma_{2}\right)$ can be a root only when $k=n=1$. And $\gamma_{1}+\left(\gamma_{10}-\gamma_{2}\right)=\gamma_{9} ; \gamma_{2}+\left(\gamma_{10}-\gamma_{2}\right)=\gamma_{10}$.

By Lemma 2.1,

$$
\begin{aligned}
& \operatorname{Int}\left(m_{2}\right) \circ U_{\gamma_{2}}\left(x_{2}^{\prime}\right)=U_{\gamma_{10}}\left(-x_{10}^{\prime}\right) U_{\gamma_{2}}\left(x_{2}^{\prime}\right), \\
& \operatorname{Int}\left(m_{2}\right) \circ U_{\gamma_{1}}\left(x_{1}^{\prime}\right)=U_{\gamma_{9}}\left(\frac{x_{1}^{\prime} x_{10}^{\prime}}{x_{2}^{\prime}}\right) U_{\gamma_{1}}\left(x_{1}^{\prime}\right)
\end{aligned}
$$

Consequently,

$$
\operatorname{Int}\left(m_{2}\right) \circ\left(\prod_{i=1}^{10} U_{\gamma_{i}}\left(x_{i}^{\prime}\right)\right)=\left[\prod_{i=1}^{8} U_{\gamma_{i}}\left(x_{i}^{\prime}\right)\right] U_{\gamma_{9}}\left(x_{9}^{\prime}-\frac{x_{1}^{\prime} x_{10}^{\prime}}{x_{2}^{\prime}}\right)
$$

For convenience of notation, let the right side of the above equation be

$$
\prod_{i=1}^{9} U_{\gamma_{i}}\left(y_{i}\right)
$$

Let

$$
m_{3}=\prod_{i=2}^{9} U_{\gamma_{i}-\gamma_{1}}\left(-\frac{y_{i}}{y_{1}}\right) .
$$

By Lemma 2.1, we have:

$$
\begin{equation*}
\operatorname{Int}\left(U_{\gamma_{9}-\gamma_{1}}\left(-\frac{y_{9}}{y_{1}}\right)\right) \circ U_{\gamma_{1}}\left(y_{1}\right)=U_{\gamma_{1}}\left(y_{1}\right) U_{\gamma_{9}}\left(-y_{9}\right) \tag{3.20}
\end{equation*}
$$

Remark For all $i$ with $i \neq 1$, and any pair $\{k, n\}$ of positive integers, the coefficient of $\alpha$ in $k\left(\gamma_{9}-\gamma_{1}\right)+n \gamma_{i}$ is $n$. So for it to be a root, $n$ must be 1 . Then the coefficient of $\alpha_{1}$ in $k\left(\gamma_{9}-\gamma_{1}\right)+n \gamma_{i}$ is $n-k=1-k$. For $k\left(\gamma_{9}-\gamma_{1}\right)+n \gamma_{i}$ to be a root, $1-k=1$ or 2 which is impossible. So $\operatorname{Int}\left(U_{\gamma_{9}-\gamma_{1}}\right)$ fixes all $U_{\gamma_{i}}$ with $i \neq 1$.

So by equation (3.20) and the above remark,

$$
\operatorname{Int}\left(U_{\gamma_{9}-\gamma_{1}}\left(-\frac{y_{9}}{y_{1}}\right)\right) \circ\left(\prod_{i=1}^{9} U_{\gamma_{i}}\left(y_{i}\right)\right)=\prod_{i=1}^{8} U_{\gamma_{i}}\left(y_{i}\right)
$$

By induction, and with the same discussion on the cases of roots as in the remark, we can prove:

$$
\operatorname{Int}\left(\prod_{i=j}^{9} U_{\gamma_{i}-\gamma_{1}}\left(-\frac{y_{i}}{y_{1}}\right)\right) \circ\left(\prod_{i=1}^{9} U_{\gamma_{i}}\left(y_{i}\right)\right)=\prod_{i=1}^{j-1} U_{\gamma_{i}}\left(y_{i}\right) .
$$

And in particular when $j=1$, then

$$
\operatorname{Int}\left(m_{3}\right) \circ\left(\prod_{i=1}^{9} U_{\gamma_{i}}\left(y_{i}\right)\right)=U_{\gamma_{1}}\left(y_{1}\right)
$$

Set $m=m_{3} m_{2} m_{1}$. We can then see from the above process that $\operatorname{Int}(m) \circ u=$ $U_{\gamma_{1}}\left(y_{1}\right)$.

For any $1 \leq i, j \leq 10, \gamma_{i}-\gamma_{j} \in \operatorname{span}\left\{\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4}, \alpha_{5}\right\}$. But for any $\gamma \in$ $\operatorname{span}\left\{\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4}, \alpha_{5}\right\}$ and any pair $\{k, n\}$ of positive integers, $k \beta_{1}+n \gamma$ cannot be a root by Lemma 3.2. So each $\operatorname{Int}\left(U_{\gamma_{i}-\gamma_{j}}\right)$ fixes $U_{\beta_{1}}$ and consequently, all $\operatorname{Int}\left(m_{1}\right)$, $\operatorname{Int}\left(m_{2}\right), \operatorname{Int}\left(m_{3}\right)$ fix $U_{\beta_{1}}$ and therefore $\operatorname{Int}(m)$ fixes $U_{\beta_{1}}$.

Theorem 3.11 (Gaussian Elimination) For any $u \in N$, there exists $m \in M$, such that $\operatorname{Int}(m) \circ u=U_{\beta_{1}}\left(a_{1}\right) U_{\gamma_{1}}\left(a_{2}\right)$, with $a_{1}, a_{2} \in F$.

Proof Write $u=u_{1} u_{2}$, where

$$
u_{1}=\prod_{i=1}^{17} U_{\beta_{i}}\left(x_{i}\right) \in N_{1}, \quad u_{2}=\prod_{i=1}^{10} U_{\gamma_{i}}\left(y_{i}\right) \in N_{2}
$$

If $u_{1}=1$, then it is the case of Lemma 3.10.
If $u_{1} \neq 1$, by applying a suitable $\operatorname{Int}(m)$ on $u$ from Lemma 3.9, we can assume $x_{1} \neq 1$. Let

$$
m_{1}=\prod_{i=2}^{17} U_{\beta_{i}-\beta_{1}}\left(\frac{x_{i}}{x_{1}}\right)
$$

then $\beta_{i}-\beta_{1}$ is a positive root and the coefficient of $\alpha_{6}$ in $\beta_{i}-\beta_{1}$ is 1 .
For any fixed $j$, with $2 \leq j \leq 17$, and for each pair of positive integers $\{k, n\}$, the coefficient of $\alpha_{6}$ in $k\left(\beta_{i}-\beta_{1}\right)+n \beta_{j}$ is $k+n \geq 2$, so $k\left(\beta_{i}-\beta_{1}\right)+n \beta_{j} \in \theta_{2}$ if it is a root. Moreover, for any $\gamma \in \theta_{2}$, the coefficient of $\alpha_{6}$ in $k\left(\beta_{i}-\beta_{1}\right)+n \gamma$ is $k+2 n \geq 3$, so $k\left(\beta_{i}-\beta_{1}\right)+n \gamma$ cannot be a root, hence $\operatorname{Int}\left(U_{\beta_{i}-\beta_{1}}\right)$ fixes every element in $N_{2}$.

So by Lemma 2.1, we have:

$$
\operatorname{Int}\left(U_{\beta_{i}-\beta_{1}}\left(\frac{x_{i}}{x_{1}}\right)\right) \circ U_{\beta_{j}}\left(x_{j}\right)=U_{\beta_{j}}\left(x_{j}\right) \cdot n_{i, j}, \quad \text { with } n_{i, j} \in N_{2}
$$

Consequently, $\operatorname{Int}\left(m_{1}\right) \circ U_{\beta_{j}}\left(x_{j}\right)=U_{\beta_{j}}\left(x_{j}\right) n_{j}$ with

$$
n_{j}=\prod_{i=2}^{17} n_{i, j} \in N_{2}
$$

and

$$
\operatorname{Int}\left(m_{1}\right) \circ U_{\beta_{1}}\left(x_{1}\right)=U_{\beta_{1}}\left(x_{1}\right) \cdot \prod_{i=2}^{17} U_{\beta_{i}}\left(-x_{i}\right) \cdot n_{1} \quad \text { with } n_{1} \in N_{2}
$$

So

$$
\operatorname{Int}\left(m_{1}\right) \circ u_{1}=\operatorname{Int}\left(m_{1}\right) \circ\left(\prod_{i=1}^{17} U_{\beta_{i}}\left(x_{i}\right)\right)=U_{\beta_{1}}\left(x_{1}\right) \cdot n \text { where } n=\prod_{i=1}^{17} n_{i} \in N_{2}
$$

Now let $u_{2}^{\prime}=n \cdot u_{2}$ and apply Lemma 3.10 to $u_{2}^{\prime}$. There exists $m_{2} \in M$ such that $\operatorname{Int}\left(m_{2}\right) \circ u_{2}^{\prime}=U_{\gamma_{1}}\left(a_{2}\right)$ and $\operatorname{Int}\left(m_{2}\right) \circ U_{\beta_{1}}\left(x_{1}\right)=U_{\beta_{1}}\left(x_{1}\right)$. Let $m=m_{2} m_{1}$ and $a_{1}=x_{1}$. Then $\operatorname{Int}(m) \circ u=U_{\beta_{1}}\left(a_{1}\right) U_{\gamma_{1}}\left(a_{2}\right)$.

Now start from $s_{1} s_{1}^{-}=s_{2}^{-} s_{2}$ acting on $\operatorname{Int}(m)$ on both sides, we can assume $s_{2}=$ $U_{\beta_{1}}\left(a_{1}\right) U_{\gamma_{1}}\left(a_{2}\right)$. The proof of the main theorem is almost the same as that of $E_{6}$. We need only make a small justification of the fact that $\gamma_{1}-\gamma_{10}$ is not a root, but this does not make much difference. Each step in the proof of the $E_{6}$ case can be paralleled to finish the proof in the $E_{7}$ case.

## 4 Application to Intertwining Operators

Now by Theorem 2.2, $M_{m_{i}}^{t}=M_{n_{i}}$. This can be used to refine the main results in [8]. To be more precise, let $X(\mathbf{M})_{F}$ be the group of $F$-rational characters of $\mathbf{M}$. Denote by $\mathbf{A}$ the split component of the center of $\mathbf{M}$. Then $\mathbf{A} \subset \mathbf{A}_{0}$. Let

$$
\left.\mathfrak{a}=\operatorname{Hom}\left(X(\mathbf{M})_{F}\right), \mathbb{R}\right)=\operatorname{Hom}\left(X(\mathbf{A})_{F}, \mathbb{R}\right)
$$

be the real Lie algebra of $\mathbf{A}$. Set $\mathfrak{a}^{*}=X(\mathbf{M})_{F} \otimes_{\mathbb{Z}} \mathbb{R}$ and $\mathfrak{a}_{\mathbb{C}}^{*}=\mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$ to denote its real and complex duals.

For $\nu \in \mathfrak{a}_{\mathbb{C}}^{*}$ and $\sigma$ an irreducible admissible representation of $M$, let $I(\nu, \sigma)=$ $\operatorname{Ind}_{M N \uparrow G} \sigma \otimes q^{\left\langle\nu, H_{P}(\cdot)\right\rangle} \otimes 1$, where $H_{P}$ is the extension of the homomorphism $H_{M}: M \rightarrow \mathfrak{a}=\operatorname{Hom}\left(X(\mathbf{M})_{F}, \mathbb{R}\right)$ to $P$, extended trivially along $N$, defined by $q^{\left\langle\chi, H_{P}(m)\right\rangle}=|\chi(m)|_{F}$ for all $\chi \in X(\mathbf{M})_{F}$. Let $V(\nu, \sigma)$ be the space of $I(\nu, \sigma)$, for $h \in V(\nu, \sigma)$, and let

$$
A(\nu, \sigma, w) h(g)=\int_{N_{\bar{w}}} h\left(w^{-1} n g\right) d n
$$

where $N_{\tilde{w}}=U \cap w N^{-} w^{-1}$, be the standard intertwining operator from $I(\nu, \sigma)$ into $I(w(\nu), w(\sigma))$.

Let $I(\sigma)=I(0, \sigma)$ and $V(\sigma)=V(0, \sigma)$ be the induced representation and its space at $\nu=0$, respectively. Since $w_{0}(M)=M, I(\sigma)$ is irreducible if and only if $A\left(\nu, \sigma, w_{0}\right)$ has a pole at $\nu=0(c f .[6-8])$. By [7, Lemma 4.1], it is enough to determine the pole of $\int_{N} h\left(w_{0}^{-1} n\right) d n$ at $\nu=0$ for any $h$ in $V(\nu, \sigma)$ which is supported in $P N^{-}$.

For $n_{i} \in N$, suppose $n_{i}$ is inside an open orbit under $\operatorname{Int}(M)$, with $w_{0}^{-1} n_{i} \in P N^{-}$. Write $w_{0}^{-1} n_{i}=m_{i} n_{i}^{\prime} n_{i}^{-}$as before, define $d^{*} n_{i}=q^{\left\langle\rho, H_{M}\left(m_{i}\right)\right\rangle} d n$ where $\rho$ is half the summation of the positive roots in $N$. Then by [8, Lemma 2.3], the measure $d^{*} n_{i}$ is an invariant measure on $M / M_{n_{i}}$ and thus induces a measure on the quotient $M / M_{n_{i}}$.

For the purpose of computing the residue we may assume that there exists a Schwartz function $\phi$ on $\mathfrak{N}^{-}$, the Lie algebra of $N^{-}$, such that

$$
h\left(\exp \left(\mathfrak{n}^{-}\right)\right)=\phi\left(\mathfrak{n}^{-}\right) h(e)
$$

where $\mathfrak{n}^{-} \in \mathfrak{N}^{-}$. Let $n_{i}^{-}=\exp \left(\mathfrak{n}_{i}^{-}\right)$, with $\mathfrak{n}_{i}^{-} \in \mathfrak{N}^{-}$. Given a representation $\sigma$, let $\psi(m)$ be among the matrix coefficients of $\sigma$, i.e, choose an arbitrary element $\tilde{v}$ in the contragredient space of $\sigma$. Let $\psi(m)=\langle\sigma(m) h(e), \tilde{v}\rangle$. With these notations and applying Theorem 2.2, [8, Proposition 2.4] can be restated as:

Proposition 4.1 Let $\sigma$ be an irreducible admissible representation of $M$. Then the poles of $A\left(\nu, \sigma, w_{0}\right)$ are the same as those of

$$
\sum_{n_{i} \in O_{i}} \int_{M / M_{n_{i}}} q^{\left\langle\nu, H_{M}\left(w_{0}(m) m_{i} m^{-1}\right)\right\rangle} \phi\left(\operatorname{Ad}\left(m^{-1}\right) \mathfrak{n}_{i}^{-}\right) \psi\left(w_{0}(m) m_{i} m^{-1}\right) d \dot{m}
$$

where $O_{i}$ runs through a finite number of open orbits of $\mathfrak{M}$ under $\operatorname{Ad}(M) ; \mathfrak{n}_{i}$ is a representative of $O_{i}$, under the correspondence that $w_{0}^{-1} n_{i}=m_{i} n_{i}^{\prime} n_{i}^{-}$, with $n_{i}=\exp \left(\mathfrak{n}_{i}\right)$, $n_{i}^{-}=\exp \left(\mathfrak{n}_{i}^{-}\right)$and d $\dot{m}$ is the measure on $M / M_{n_{i}}$ induced from $d^{*} n_{i}$.

Let $\tilde{\mathbf{A}}$ be the center of $\mathbf{M}$. Then there exists a function $f \in C_{c}^{\infty}(M)$ such that $\psi(m)=\int_{\tilde{A}} f(a m) \omega^{-1}(a) d a$, where $\omega$ is the central character of $\sigma$.

Define

$$
\theta: M \rightarrow M, \quad \theta(m)=w_{0}^{-1} m w_{0}, \forall m \in M
$$

Given $f \in C_{c}^{\infty}(M)$ and $m_{0} \in M$, define the $\theta$-twisted orbit integral for $f$ at $m_{0}$ by:

$$
\phi_{\theta}\left(m_{0}, f\right)=\int_{M / M_{\theta, m_{0}}} f\left(\theta(m) m_{0} m^{-1}\right) d \dot{m}
$$

where

$$
M_{\theta, m_{0}}=M_{\theta, m_{0}}(F)=\left\{m \in M(F) \mid \theta(m) m_{0} m^{-1}=m_{0}\right\}
$$

is the $\theta$-twisted centralizer of $m_{0}$ in $M(F), d \dot{m}$ is the measure on $M / M_{\theta, m_{0}}$ induced from $d m$.

Applying our Theorem 2.2, the main theorem in [8] (Theorem 2.5) can be modified as:

Proposition 4.2 Assume $\sigma$ is supercuspidal and $w_{0}(\sigma) \cong \sigma$. The intertwining operator $A\left(\nu, \sigma, w_{0}\right)$ has a pole at $\nu=0$ if and only if

$$
\sum_{i} \int_{Z(G) / Z(G) \cap w_{0}(\tilde{A}) \tilde{A}^{-1}} \phi_{\theta}\left(z m_{i}, f\right) \omega^{-1}(z) d z \neq 0
$$

for $f$ as above. Here $Z(G)$ is the center of $G$ and

$$
\phi_{\theta}\left(z m_{i}, f\right)=\int_{M / M_{n_{i}}} f\left(z \theta(m) m_{i} m^{-1}\right) d \dot{m}
$$

is the $\theta$-twisted orbital integral for $f$ at $z m_{i}$, where $m_{i}$ corresponds to the representatives $\left\{n_{i}\right\}$ for the open orbits in $N$ under $\operatorname{Int}(M)$, with $w_{0}^{-1} n_{i}=m_{i} n_{i}^{\prime} n_{i}^{-}$, as $n_{i}$ runs through the finite number of open orbits in $N$.

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