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#### Abstract

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# ISOTROPIC RANDOM WALKS ON AFFINE BUILDINGS 

by James PARKINSON


#### Abstract

In this paper we apply techniques of spherical harmonic analysis to prove a local limit theorem, a rate of escape theorem, and a central limit theorem for isotropic random walks on arbitrary thick regular affine buildings of irreducible type. This generalises results of Cartwright and Woess where $\tilde{A}_{n}$ buildings are studied, Lindlbauer and Voit where $\tilde{A}_{2}$ buildings are studied, and Sawyer where homogeneous trees are studied (these are $\tilde{A}_{1}$ buildings).

RÉSumé. - Dans cet article, nous utilisons les techniques de l'analyse harmonique sphérique pour démontrer un théorème local limite, un théorème sur la vitesse de fuite et un théorème central limite pour les marches aléatoires isotropes sur des immeubles affines épais arbitraires de type irréductible. Cela généralise des résultats de Cartwright et Woess sur les immeubles de type $\tilde{A}_{n}$, de Lindlbauer et Voit sur les immeubles de type $\tilde{A}_{2}$ et de Sawyer sur les arbres homogènes (qui sont des immeubles de type $\tilde{A}_{1}$ ).


## Introduction

Let $\mathscr{X}$ be a thick locally finite regular affine building of irreducible type. By regular we mean that the number of chambers containing a panel depends only on the cotype of the panel, and by thick we mean that this number is always at least 3 . The simplest example of such a building is a homogeneous tree with degree $q+1 \geqslant 3$, where the chambers are the edges of the graph. In this case, Sawyer [18] studied isotropic random walks $\left(Z_{k}\right)_{k \geqslant 0}$ on the vertices of $\mathscr{X}$, meaning that the transition probabilities $p(x, y)=\mathbb{P}\left(Z_{k+1}=y \mid Z_{k}=x\right)$ depend only on the graph distance $d(x, y)$ between $x$ and $y$. To motivate our results, let us briefly describe these random walks on trees.

Let $V$ be the vertex set of the tree, and for each $x \in V$ and $k \in \mathbb{N}=$ $\{0,1, \ldots\}$, write $V_{k}(x)$ for the set of all $y \in V$ such that $d(x, y)=k$. It is easily seen that the cardinalities $\left|V_{k}(x)\right|$ are independent of the particular $x \in V$, and we write $N_{k}$ for this value. For each $k \in \mathbb{N}$ there is a natural operator $A_{k}$ acting on the space of all functions $f: V \rightarrow \mathbb{C}$, where for each $x \in V,\left(A_{k} f\right)(x)$ is the average value of $f$ over $V_{k}(x)$. The operator $A_{k}$ may be regarded as the transition operator of the isotropic random walk with matrix $\left(p_{k}(x, y)\right)_{x, y \in V}$, where $p_{k}(x, y)=\frac{1}{N_{k}}$ if $y \in V_{k}(x)$ and $p_{k}(x, y)=0$ otherwise. Indeed it is easily seen that a random walk on $V$ is isotropic if and only if it has a transition operator of the form

$$
A=\sum_{k \in \mathbb{N}} a_{k} A_{k}
$$

where $a_{k} \geqslant 0$ for all $k \in \mathbb{N}$ and $\sum_{k \in \mathbb{N}} a_{k}=1$.
The linear span over $\mathbb{C}$ of the operators $\left\{A_{k}\right\}_{k \in \mathbb{N}}$ is a commutative algebra $\mathscr{A}$ with a rich theory of harmonic analysis (see [7]). In particular, the algebra homomorphisms $h: \mathscr{A} \rightarrow \mathbb{C}$ may be explicitly described, and local limit theorems, central limit theorems, and rate of escape theorems can be proved as applications.

Now consider a regular affine building $\mathscr{X}$ of irreducible type. Thus $\mathscr{X}$ may be regarded as a simplicial complex made by 'gluing together' many copies of a given Coxeter complex, each Coxeter complex called an apartment of $\mathscr{X}$ (these are regular tessellations of Euclidean space by simplices). There is an irreducible (but not necessarily reduced) root system $R$ associated to $\mathscr{X}$, as described in Section 1.3, and the coweight lattice $P$ of $R$ is a subset of the vertex set of the standard Coxeter complex $\Sigma$. We consider random walks on a related subset $V_{P}$ of the vertices of $\mathscr{X}$, which in most cases is the set of all special vertices of $\mathscr{X}$.

Let $P^{+}$be the set of dominant coweights of $R$ (relative to some fixed base). For each $x \in V_{P}$ there is a natural partition of $V_{P}$ into subsets $V_{\lambda}(x), \lambda \in P^{+}$, as described in Definition 1.4. Roughly speaking, $y \in V_{\lambda}(x)$ means that there exists an apartment $\mathcal{A}$ containing $x$ and $y$ and a 'suitable' isomorphism $\psi: \mathcal{A} \rightarrow \Sigma$ such that $\psi(x)=0$ and $\psi(y)=\lambda \in P^{+}$(in other words, $y$ is in position $\lambda$ from $x$ ). It is shown in [15, Theorem 5.15] that for all $\lambda \in P^{+}$the cardinality of the set $V_{\lambda}(x)$ is independent of the particular $x \in V_{P}$, and we write $N_{\lambda}$ for this value. Following the tree case, for each $\lambda \in P^{+}$let $A_{\lambda}$ be the operator acting on the space of functions $f: V_{P} \rightarrow \mathbb{C}$ with $\left(A_{\lambda} f\right)(x)$ being the average value of $f$ over $V_{\lambda}(x)$. The linear span of these operators over $\mathbb{C}$ is a commutative algebra $\mathscr{A}$, which has been studied extensively in [15]. As shown in [16], the algebra homomorphisms
$h: \mathscr{A} \rightarrow \mathbb{C}$ may be explicitly described both in terms of the Macdonald spherical functions and in terms of an integral over the boundary of $\mathscr{X}$.
We call a random walk $\left(Z_{k}\right)_{k \in \mathbb{N}}$ on $V_{P}$ with transition matrix $(p(x, y))_{x, y \in V_{P}}$ isotropic if $p(x, y)=p\left(x^{\prime}, y^{\prime}\right)$ whenever $y \in V_{\lambda}(x)$ and $y^{\prime} \in V_{\lambda}\left(x^{\prime}\right)$ for some $\lambda \in P^{+}$. As in the tree case, the operators $A_{\lambda}$ may be regarded as the transition operators of isotropic random walks with matrices $\left(p_{\lambda}(x, y)\right)_{x, y \in V_{P}}$, where $p_{\lambda}(x, y)=\frac{1}{N_{\lambda}}$ if $y \in V_{\lambda}(x)$ and $p_{\lambda}(x, y)=0$ otherwise. It is easily seen that a random walk on $V_{P}$ is isotropic if and only if it has a transition operator of the form

$$
\begin{equation*}
A=\sum_{\lambda \in P^{+}} a_{\lambda} A_{\lambda} \tag{0.1}
\end{equation*}
$$

where $a_{\lambda} \geqslant 0$ for all $\lambda \in P^{+}$and $\sum_{\lambda \in P^{+}} a_{\lambda}=1$.
In this paper we apply the spherical harmonic analysis associated to the algebra $\mathscr{A}$ to prove a local limit theorem, a central limit theorem, and a rate of escape theorem for isotropic random walks on $V_{P}$. These results generalise the results in [5] where $\tilde{A}_{n}$ buildings are studied (which in turn generalise the corresponding results for homogeneous trees). Our results may also be viewed as 'building analogues' of well known results concerning random walks on semisimple Lie groups (see [1] for example).

Let us briefly outline the structure of this paper. In Section 1 we give a summary of some background material, mostly from [15] and [16]. This section includes a discussion of root systems, Coxeter complexes and buildings, the algebra $\mathscr{A}$, and the spherical harmonic analysis associated to this algebra. The main sections of this paper are Sections 2, 3 and 4.

In Section 2 we give our local limit theorem for isotropic random walks on $V_{P}$, describing the asymptotic behaviour of the $k$-step transition probabilities $p^{(k)}(x, y)=\mathbb{P}\left(Z_{k}=y \mid Z_{0}=x\right)$. We also give necessary and sufficient conditions for irreducibility and aperiodicity of the random walk, and in Remark 2.19 we outline some applications of our local limit theorem to random walks on groups acting on buildings.

In Section 3 we prove our rate of escape theorem. For each $k \in \mathbb{N}$, let $\nu_{k} \in P^{+}$be such that $Z_{k} \in V_{\nu_{k}}\left(Z_{0}\right)$. We show that, with probability 1 , the vector $\frac{1}{k} \nu_{k}$ converges to a vector $\gamma$ in the underlying vector space of the root system $R$. We apply our local limit theorem to show that each component of $\gamma$ (relative to a set of fundamental coweights of $P$ ) is strictly positive.

In Section 4 we prove our central limit theorem, showing that there is a positive definite matrix $\Gamma$ such that, with $\gamma$ as above, the vector $\frac{1}{\sqrt{k}}\left(\nu_{k}-k \gamma\right)$ tends in distribution to the normal distribution $N(0, \Gamma)$.

In Appendix A we determine the algebra homomorphisms $h: \mathscr{A} \rightarrow \mathbb{C}$ which are bounded (generalising [10, Theorem 4.7.1]).

## 1. Affine Buildings and Spherical Harmonic Analysis

### 1.1. Root Systems and Weyl Groups

Root systems play a significant role in this work. We fix the following notations and conventions, generally following [2, Chapter VI].

Let $R$ be an irreducible, but not necessarily reduced, root system in a real vector space $E$ with inner product $\langle\cdot, \cdot\rangle$. The rank of $R$ is $n$, the dimension of $E$. Let $I_{0}=\{1,2, \ldots, n\}$ and $I=\{0,1,2, \ldots, n\}$, and let $B=\left\{\alpha_{i} \mid i \in I_{0}\right\}$ be a fixed base of $R$. Write $R^{+}$for the set of positive roots (relative to $B$ ), and let $R^{\vee}=\left\{\alpha^{\vee} \mid \alpha \in R\right\}$ be the dual root system of $R$, where for each $\alpha \in R$ we write $\alpha^{\vee}=\frac{2 \alpha}{\langle\alpha, \alpha\rangle}$. Since $R$ is irreducible, by [2, VI, § $1, \mathrm{~N}^{\circ} 8$, Proposition 25] there is a unique highest root

$$
\begin{equation*}
\tilde{\alpha}=\sum_{i \in I_{0}} m_{i} \alpha_{i} \tag{1.1}
\end{equation*}
$$

with the property that if $\beta=\sum_{i \in I_{0}} k_{i} \alpha_{i} \in R$ then $m_{i} \geqslant k_{i}$ for each $i \in I_{0}$.
For each $i \in I_{0}$ define $\lambda_{i} \in E$ by $\left\langle\lambda_{i}, \alpha_{j}\right\rangle=\delta_{i, j}$ for all $j \in I_{0}$. The elements $\left\{\lambda_{i}\right\}_{i \in I_{0}}$ are called the fundamental coweights of $R$. The coweight lattice of $R$ is $P=\sum_{i \in I_{0}} \mathbb{Z} \lambda_{i}$, and elements $\lambda \in P$ are called coweights of $R$. A coweight $\lambda \in P$ is said to be dominant if $\left\langle\lambda, \alpha_{i}\right\rangle \geqslant 0$ for all $i \in I_{0}$, and we write $P^{+}$for the set of all dominant coweights. Let $Q=\sum_{\alpha \in R} \mathbb{Z} \alpha^{\vee}$ be the coroot lattice of $R$, and let $Q^{+}=\sum_{\alpha \in R^{+}} \mathbb{N} \alpha^{\vee}$. Note that $Q \subseteq P$, and by [2, VI, § 1, N ${ }^{\circ} 8$, Proposition 25] we have $\tilde{\alpha}^{\vee} \in P^{+}$.

For each $\alpha \in R$ and $k \in \mathbb{Z}$ let $H_{\alpha ; k}=\{x \in E \mid\langle x, \alpha\rangle=k\}$. We call these sets affine hyperplanes, or simply hyperplanes. For each $\alpha \in R$ and $k \in \mathbb{Z}$ let $s_{\alpha ; k}$ denote the orthogonal reflection in $H_{\alpha ; k}$. Thus $s_{\alpha ; k}(x)=$ $x-(\langle x, \alpha\rangle-k) \alpha^{\vee}$ for all $x \in E$. Write $s_{\alpha}$ in place of $s_{\alpha ; 0}, s_{i}$ in place of $s_{\alpha_{i}}\left(\right.$ for $\left.i \in I_{0}\right)$, and let $s_{0}=s_{\tilde{\alpha} ; 1}$. Let $W_{0}=W_{0}(R)$ be the Weyl group of $R$, and let $W=W(R)$ be the affine Weyl group of $R$. Thus $W_{0}$ is the subgroup of GL $(E)$ generated by $S_{0}=\left\{s_{i}\right\}_{i \in I_{0}}$, and $W$ is the subgroup of $\operatorname{Aff}(E)$ generated by $S=\left\{s_{i}\right\}_{i \in I}$. Both $\left(W_{0}, S_{0}\right)$ and $(W, S)$ are Coxeter systems, and clearly $W_{0} \leqslant W$. Given $w \in W$, we define the length $\ell(w)$ of $w$ to be smallest $k \in \mathbb{N}$ such that $w=s_{i_{1}} \ldots s_{i_{k}}$, with $i_{1}, \ldots, i_{k} \in I$.

The extended affine Weyl group of $R$ is $\widetilde{W}(R)=\widetilde{W}=W_{0} \ltimes P$. Since $W \cong W_{0} \ltimes Q\left(\right.$ see $\left[2, \mathrm{VI}, \S 2, \mathrm{~N}^{\circ} 1\right.$, Proposition 1]) and $Q \subseteq P$, we may
regard $W$ as a subgroup of $\widetilde{W}$. Note that $\widetilde{W}$ contains all translations by elements of $P$, while $W$ only contains those translations by elements of $Q$.

Remark 1.1. - We make the following comments for those readers not so familiar with the non-reduced root systems. For each $n \geqslant 1$ there is exactly one irreducible non-reduced root system (up to isomorphism) of rank $n$, denoted by $B C_{n}$. To describe this root system we may take $E=\mathbb{R}^{n}$ with the usual inner product, and let $R$ consist of the vectors $\pm e_{i}, \pm 2 e_{i}$ and $\pm e_{j} \pm e_{k}$ for $1 \leqslant i \leqslant n$ and $1 \leqslant j<k \leqslant n$. Let $\alpha_{i}=e_{i}-e_{i+1}$ for $1 \leqslant i<n$ and $\alpha_{n}=e_{n}$. Then $B=\left\{\alpha_{i}\right\}_{i \in I_{0}}$ is a base of $R$, and $R^{+}$consists of the vectors $e_{i}, 2 e_{i}$ and $e_{j} \pm e_{k}$ for $1 \leqslant i \leqslant n$ and $1 \leqslant j<k \leqslant n$. The fundamental coweights are $\lambda_{i}=e_{1}+\cdots+e_{i}$ for each $1 \leqslant i \leqslant n$. Note that $R^{\vee}=R$ and $Q=P$. The subsystem $R_{1}=\{\alpha \in R \mid 2 \alpha \notin R\}$ is a root system of type $C_{n}$, and the subsystem $R_{2}=\left\{\alpha \in R \left\lvert\, \frac{1}{2} \alpha \notin R\right.\right\}$ is a root system of type $B_{n}$ (with the convention that $C_{1}=B_{1}=A_{1}$ ). We have $Q(R)=P(R)=Q\left(R_{1}\right) \subset P\left(R_{1}\right)$ (with strict inclusion), and so $W(R)=W\left(R_{1}\right)$ but $\widetilde{W}(R) \neq \widetilde{W}\left(R_{1}\right)$.

### 1.2. The Coxeter Complex

There is a natural geometric realisation $\Sigma=\Sigma(R)$ of the Coxeter complex of $W=W(R)$. Let $\mathcal{H}$ denote the family of the hyperplanes $H_{\alpha ; k}, \alpha \in R$, $k \in \mathbb{Z}$, and define chambers of $\Sigma$ to be open connected components of $E \backslash \bigcup_{H \in \mathcal{H}} H$. Since $R$ is irreducible each chamber is an open (geometric) simplex $\left[2, \mathrm{~V}, \S 3, \mathrm{~N}^{\circ} 9\right.$, Proposition 8]. We call the extreme points of the closure of chambers vertices of $\Sigma$, and we write $V(\Sigma)$ for the set of all vertices of $\Sigma$.

The set $P$ of coweights of $R$ is a subset of $V(\Sigma)$, and we call elements of $P$ the good vertices of $\Sigma$. When $R$ is reduced, $P$ is the set of more familiar special vertices of $\Sigma\left[2, \mathrm{VI}, \S 2, \mathrm{~N}^{\circ} 2\right.$, Proposition 3].

The choice of the base $B=\left\{\alpha_{i}\right\}_{i=1}^{n}$ gives a natural choice of a fundamental chamber

$$
\begin{equation*}
C_{0}=\left\{x \in E \mid\left\langle x, \alpha_{i}\right\rangle>0 \text { for all } i \in I_{0} \text { and }\langle x, \tilde{\alpha}\rangle<1\right\} . \tag{1.2}
\end{equation*}
$$

In the notation of (1.1), the vertices of $C_{0}$ are the points $\{0\} \cup\left\{\lambda_{i} / m_{i}\right\}_{i \in I_{0}}$ $\left[2, \mathrm{VI}, \S 2, \mathrm{~N}^{\circ} 2\right]$. There is a natural simplicial complex structure on $\Sigma$ with maximal simplices being the vertex sets of chambers of $\Sigma$, and simplices being subsets of the maximal simplices. We define $\tau: V(\Sigma) \rightarrow I$ to be the unique labelling of $\Sigma$ (as a simplicial complex) such that $\tau(0)=0$ and $\tau\left(\lambda_{i} / m_{i}\right)=i$ for each $i \in I_{0}$.

We write $I_{P}=\{\tau(\lambda) \mid \lambda \in P\} \subseteq I$. Let $\left\{m_{i}\right\}_{i \in I_{0}}$ be as in (1.1), and define $m_{0}=1$. We have $I_{P}=\left\{i \in I \mid m_{i}=1\right\}$, which shows that $0 \in I_{P}$ for all root systems, and that $I_{P}=\{0\}$ if $R$ is non-reduced [15, Lemma 4.3]. This also shows that in the non-reduced case, and only in the non-reduced case, there are special vertices which are not good vertices.

We define the fundamental sector of $\Sigma$ to be the open simplicial cone

$$
\begin{equation*}
\mathcal{S}_{0}=\left\{x \in E \mid\left\langle x, \alpha_{i}\right\rangle>0 \text { for all } i \in I_{0}\right\} \tag{1.3}
\end{equation*}
$$

The sectors of $\Sigma$ are then the sets $\lambda+w \mathcal{S}_{0}$, where $w \in W_{0}$ and $\lambda \in P$ (equivalently, the sectors are the sets $\tilde{w} \mathcal{S}_{0}, \tilde{w} \in \widetilde{W}$ ).

An automorphism of $\Sigma$ is a bijection $\psi$ of $E$ which maps chambers, and only chambers, to chambers, with the property that chambers $C$ and $D$ are adjacent if and only if $\psi(C)$ is adjacent to $\psi(D)$. We write $\operatorname{Aut}(\Sigma)$ for the automorphism group of $\Sigma$. An automorphism $\psi$ of $\Sigma$ is called type preserving if $\tau(v)=\tau(\psi(v))$ for all $v \in V(\Sigma)$. By [17, Lemma 2.2] $\psi \in$ $\operatorname{Aut}(\Sigma)$ is type preserving if and only if $\psi \in W$. Generally we have $W_{0}<$ $W \leqslant \widetilde{W} \leqslant \operatorname{Aut}(\Sigma)$ (with the possibility that $W<\widetilde{W}$ and $\widetilde{W}<\operatorname{Aut}(\Sigma)$ ).

### 1.3. Buildings and Regularity

Recall ([3]) that a building of type $W$ is a nonempty simplicial complex $\mathscr{X}$ which contains a family of subcomplexes called apartments such that:
(i) each apartment is isomorphic to the (simplicial) Coxeter complex of $W$,
(ii) given any two chambers of $\mathscr{X}$ there is an apartment containing both, and
(iii) given any two apartments $\mathcal{A}$ and $\mathcal{A}^{\prime}$ that contain a common chamber, there exists an isomorphism $\psi: \mathcal{A} \rightarrow \mathcal{A}^{\prime}$ fixing $\mathcal{A} \cap \mathcal{A}^{\prime}$ pointwise.
Since $W$ is an affine Weyl group, $\mathscr{X}$ is called an affine building.
It is a consequence of this definition that $\mathscr{X}$ is a labellable simplicial complex, and all the isomorphisms in the above definition may be taken to be type preserving (this ensures that the labellings of $\mathscr{X}$ and $\Sigma$ are compatible).

Let $V$ and $\mathcal{C}$ be the vertex and chamber sets of $\mathscr{X}$, respectively (with chambers being maximal simplices of $\mathscr{X}$ ). Chambers $c$ and $d$ are declared to be $i$-adjacent (written $c \sim_{i} d$ ) if and only if either $c=d$, or if all the vertices of $c$ and $d$ are the same except for those of type $i$.

Throughout this paper we assume that our buildings are:
(i) locally finite, meaning that $|I|<\infty$ and $\left|\left\{d \in \mathcal{C} \mid d \sim_{i} c\right\}\right|<\infty$ for each $c \in \mathcal{C}$ and $i \in I$,
(ii) regular, meaning that $\left|\left\{d \in \mathcal{C} \mid d \sim_{i} c\right\}\right|$ is independent of $c \in \mathcal{C}$ for each $i \in I$, and
(iii) thick, meaning that $\left|\left\{d \in \mathcal{C} \mid d \sim_{i} c\right\}\right| \geqslant 3$ for each $c \in \mathcal{C}$ and each $i \in I$.

By [15, Theorem 2.4] we see that thickness and regularity are intimately connected. Indeed, the only thick affine buildings of irreducible type which may fail to be regular are those of dimension 1 (thus regularity is a very weak hypothesis).

Since $\mathscr{X}$ is assumed to be regular, we may define numbers $q_{i}, i \in I$, called the parameters of the building, by $q_{i}+1=\left|\left\{d \in \mathcal{C} \mid d \sim_{i} c\right\}\right|$. These numbers satisfy $q_{j}=q_{i}$ if $s_{j}=w s_{i} w^{-1}$ for some $w \in W$ (see [15, Corollary 2.2]), and by thickness $q_{i}>1$ for all $i \in I$. If $w=s_{i_{1}} \cdots s_{i_{k}} \in$ $W$ is a reduced expression (that is, $\ell(w)=k$ ) we define $q_{w}=q_{i_{1}} \cdots q_{i_{k}}$, which is independent of the particular reduced expression for $w$ (see [15, Proposition 2.1(i)]).

To each locally finite regular affine building of irreducible type we associate an irreducible root system $R$ (depending on the parameter system of the building) as follows (see [15, Appendix]):
(i) If $\mathscr{X}$ is a regular $\tilde{A}_{1}$ building with $q_{0}=q_{1}$, then we take $R=A_{1}$ (these buildings are homogeneous trees).
(ii) If $\mathscr{X}$ is a regular $\tilde{A}_{1}$ building with $q_{0} \neq q_{1}$, then we take $R=B C_{1}$ (these buildings are semi-homogeneous trees).
(iii) If $\mathscr{X}$ is a regular $\tilde{C}_{n}$ building with $n \geqslant 2$ and $q_{0}=q_{n}$, then we take $R=C_{n}$.
(iv) If $\mathscr{X}$ is a regular $\tilde{C}_{n}$ building with $n \geqslant 2$ and $q_{0} \neq q_{n}$, then we take $R=B C_{n}$.
(v) If $\mathscr{X}$ is a regular building of type $\tilde{X}_{n}$, where $X=A$ and $n \geqslant 2$, or $X=B$ and $n \geqslant 3$, or $X=D$ and $n \geqslant 4$, or $X=E$ and $n=6,7$ or 8 , or $X=F$ and $n=4$, or $X=G$ and $n=2$, then we take $R=X_{n}$.

The choices above are made to ensure that the coweight lattice $P$ of $R$ preserves the parameter system of $\mathscr{X}$ in the sense that if $v \in V(\Sigma)$ then $q_{\tau(v)}=q_{\tau(v+\lambda)}$ for all $\lambda \in P$. Thus, for example, (iv) above is motivated by the general parameter system of a $\tilde{C}_{n}$ building (embodied in the Coxeter graph):


Figure 1.1
(see [15, Appendix]). If we take $R=C_{n}$, then by the definition of the type map (see Section 1.2) and the fact that $m_{n}=1$ (see (1.1) and [2, Plate III]) we have $\tau\left(\lambda_{n}\right)=n$. Thus in general we have $q_{\tau\left(\lambda_{n}\right)}=q_{n} \neq q_{0}=q_{\tau(0)}$. If we instead choose $R=B C_{n}$, then $P=Q$, and so $\tau(v+\lambda)=\tau(v)$ for all $v \in V(\Sigma)$ and $\lambda \in P\left(\right.$ and hence $\left.q_{\tau(v)}=q_{\tau(v+\lambda)}\right)$.

Definition 1.2. - Let $\mathscr{X}$ be a regular affine building with associated root system $R$ and vertex set $V$. A vertex $x \in V$ is said to be good if and only if $\tau(x) \in I_{P}$ (recall that $I_{P}=\{\tau(\lambda) \mid \lambda \in P\}$ ). Write $V_{P}$ for the set of all good vertices of $\mathscr{X}$.

It is clear that $V_{P}$ is a subset of the more familiar special vertices of $\mathscr{X}$. In fact if $R$ is reduced then $V_{P}$ equals the set of special vertices. If $R$ is non-reduced (so $R$ is of type $B C_{n}$ for some $n \geqslant 1$ ), then $V_{P}$ is the set of all type 0 vertices of $\mathscr{X}$ (whereas the special vertices are those with types 0 or $n$ ).

### 1.4. The Algebra $\mathscr{A}$

In this section we describe a commutative algebra $\mathscr{A}$ of vertex set averaging operators. This algebra has been studied in detail in [15], where it is shown that $\mathscr{A}$ is isomorphic to the center of an appropriate affine Hecke algebra.

Definition 1.3. - Let $\mathcal{A}$ be an apartment of $\mathscr{X}$. An isomorphism $\psi: \mathcal{A} \rightarrow \Sigma$ is called type rotating if and only if it is of the form $\psi=w \circ \psi_{0}$, where $\psi_{0}: \mathcal{A} \rightarrow \Sigma$ is a type preserving isomorphism, and $w \in \widetilde{W}$.

Definition 1.4. - Given $x \in V_{P}$ and $\lambda \in P^{+}$, let $V_{\lambda}(x)$ be the set of all $y \in V_{P}$ such that there exists an apartment $\mathcal{A}$ containing $x$ and $y$, and a type rotating isomorphism $\psi: \mathcal{A} \rightarrow \Sigma$ such that $\psi(x)=0$ and $\psi(y)=\lambda$. Equivalently, $y \in V_{\lambda}(x)$ if and only if there exists an apartment $\mathcal{A}$ containing $x$ and $y$ and a type rotating isomorphism $\psi: \mathcal{A} \rightarrow \Sigma$ such that $\psi(x)=0$ and $\psi(y) \in W_{0} \lambda$.

The requirement that $\psi$ is type rotating in Definition 1.4 ensures that if $y \in V_{\lambda}(x) \cap V_{\mu}(x)$ then $\lambda=\mu$. Indeed, in [15, Proposition 5.6] we showed that for each $x \in V_{P},\left\{V_{\lambda}(x)\right\}_{\lambda \in P^{+}}$forms a partition of $V_{P}$.

Remark 1.5. - To get a feel for Definition 1.4 in a special case, suppose that $\mathscr{X}$ is a homogeneous tree with degree $q+1$. Thus $R$ has type $A_{1}$, and we may take $R=\{\alpha,-\alpha\}$ where $\alpha=e_{1}-e_{2}$ (the underlying vector space here is $\left.E=\left\{\xi \in \mathbb{R}^{2} \mid \xi_{1}+\xi_{2}=0\right\}\right)$. Taking $B=\{\alpha\}$ we have $\lambda_{1}=\frac{\alpha}{2}$ and $P^{+}=\left\{k \lambda_{1} \mid k \in \mathbb{N}\right\}$. We have $V_{P}=V$ (all vertices are 'good'), and writing $V_{k}(x)$ in place of $V_{k \lambda_{1}}(x)$, we have

$$
V_{k}(x)=\{y \in V \mid d(x, y)=k\}
$$

where $d: V \times V \rightarrow \mathbb{N}$ is the usual graph metric.
Note that in this example all isomorphisms $\psi: \mathcal{A} \rightarrow \Sigma$ where $\mathcal{A}$ is an apartment of $\mathscr{X}$ are type rotating. To understand why the type rotating hypothesis becomes important, suppose that $\mathscr{X}$ is a regular $\tilde{A}_{2}$ building, and take vertices $x, y \in V_{P}$ with $y \in V_{\lambda_{1}}(x)$. Thus there exists an apartment $\mathcal{A}$ containing $x$ and $y$ and a type rotating isomorphism $\psi: \mathcal{A} \rightarrow \Sigma$ with $\psi(x)=0$ and $\psi(y)=\lambda_{1}$. The map $\varphi: \Sigma \rightarrow \Sigma$ given by $a_{1} \lambda_{1}+a_{2} \lambda_{2} \mapsto a_{2} \lambda_{1}+a_{1} \lambda_{2}$ is an automorphism of $\Sigma$, and so $\varphi \circ \psi: \mathcal{A} \rightarrow \Sigma$ is an isomorphism (however it is not type rotating). Notice that $(\varphi \circ \psi)(x)=0$ and $(\varphi \circ \psi)(y)=\lambda_{2}$, and so if we drop the type rotating hypothesis in Definition 1.4 we would conclude that $y \in V_{\lambda_{1}}(x) \cap V_{\lambda_{2}}(x)$.

For $\lambda \in P$ let $\lambda^{*}=-w_{0} \lambda$, where $w_{0}$ is the unique longest element of $W_{0}$. In [15, Proposition 5.8] we showed that if $\lambda \in P^{+}$then $\lambda^{*} \in P^{+}$, and that $y \in V_{\lambda}(x)$ if and only if $x \in V_{\lambda^{*}}(y)$. Note that $*$ is trivial unless $w_{0} \neq-1$, that is, unless $R=A_{n}, D_{2 n+1}$ or $E_{6}$ for some $n \geqslant 2$ (see [2, Plates I-IX]). For example, the map $\varphi$ from Remark 1.5 is $\lambda \mapsto \lambda^{*}$.

In [15, Theorem 5.15] we showed that $\left|V_{\lambda}(x)\right|=\left|V_{\lambda}(y)\right|$ for all $x, y \in V_{P}$ and $\lambda \in P^{+}$, and we denote this common value by $N_{\lambda}$ (see (1.5) for a formula for $N_{\lambda}$ ). For each $\lambda \in P^{+}$define an operator $A_{\lambda}$, acting on the space of functions $f: V_{P} \rightarrow \mathbb{C}$, by

$$
\left(A_{\lambda} f\right)(x)=\frac{1}{N_{\lambda}} \sum_{y \in V_{\lambda}(x)} f(y) \quad \text { for all } x \in V_{P}
$$

(thus $\left(A_{\lambda} f\right)(x)$ is the average value of $f$ over the set $\left.V_{\lambda}(x)\right)$. The linear span $\mathscr{A}$ of $\left\{A_{\lambda}\right\}_{\lambda \in P^{+}}$over $\mathbb{C}$ is a commutative algebra [15, Theorem 5.24].

Remark 1.6. - (i) In the situation of the first example of Remark 1.5, writing $N_{k}$ in place of $N_{k \lambda_{1}}$ we have $N_{0}=1$ and $N_{k}=(q+1) q^{k-1}$ for $k \geqslant 1$. In this case the operators $A_{k}=A_{k \lambda_{1}}$ have been studied by many authors (see [7, p.57], [18] or [20, § III.19.C]). They satisfy the simple recurrence

$$
A_{k} A_{1}=\frac{q}{q+1} A_{k+1}+\frac{1}{q+1} A_{k-1} \quad \text { for } k \geqslant 1
$$

although for general affine buildings such a formula is not readily available.
(ii) Let $\mathscr{A}_{Q}$ denote the linear span (over $\mathbb{C}$ ) of $\left\{A_{\lambda} \mid \lambda \in Q \cap P^{+}\right\}$. It is easily seen that $\mathscr{A}_{Q}$ is a subalgebra of $\mathscr{A}$. In the case when $\mathscr{X}$ is the Bruhat-Tits building of a group $G$ of $p$-adic type with maximal compact subgroup $K$ (as in [10, § 2.4-2.7]), $\mathscr{A}_{Q}$ is isomorphic to $\mathscr{L}(G, K)$, the space of continuous compactly supported bi- $K$-invariant functions on $G$.

### 1.5. Isotropic Random Walks

As mentioned in the introduction, we call a random walk on $V_{P}$ with transition probability matrix $A=(p(x, y))_{x, y \in V_{P}}$ isotropic if $p(x, y)=p\left(x^{\prime}, y^{\prime}\right)$ whenever $y \in V_{\lambda}(x)$ and $y^{\prime} \in V_{\lambda}\left(x^{\prime}\right)$ for some $\lambda \in P^{+}$. In particular, each operator $A_{\lambda}, \lambda \in P^{+}$, represents an isotropic random walk with transition matrix (also called $A_{\lambda}$ ) given by $A_{\lambda}=\left(p_{\lambda}(x, y)\right)_{x, y \in V_{P}}$, where $p_{\lambda}(x, y)=N_{\lambda}^{-1}$ if $y \in V_{\lambda}(x)$ and $p_{\lambda}(x, y)=0$ otherwise.

It is easily seen that a random walk is isotropic if and only if its transition matrix (operator) $A$ is as in (0.1). To avoid triviality we always assume that $a_{\lambda}>0$ for at least one $\lambda \neq 0$ (so that $A$ is not the identity). In this paper we will prove a local limit theorem, a rate of escape theorem, and a central limit theorem for such random walks, generalising the work of [18] (where homogeneous trees are studied) and [5] (where $\tilde{A}_{n}$ buildings are studied). The main techniques we use are those of spherical harmonic analysis, as recalled in the following sections. We note that isotropic random walks on $\tilde{A}_{2}$ buildings have also been studied by Lindlbauer and Voit [9] where more hypergroup oriented techniques are used (see [15, §7] for a discussion of the hypergroups that arise in the setting of general affine buildings).

In the case of Remark 1.6(ii), the theorems we prove in this paper can be translated into theorems concerning probability measures on groups of p-adic type. We briefly discuss this in Remark 2.19.

### 1.6. The Algebra Homomorphisms $h: \mathscr{A} \rightarrow \mathbb{C}$

Our proofs of the local limit theorem, rate of escape theorem and central limit theorem rely heavily on two formulae for the algebra homomorphisms $h: \mathscr{A} \rightarrow \mathbb{C}$. In this section we recall these formulae from [16]. The first formula is in terms of the Macdonald spherical functions, and the second is in terms of an integral over the boundary of $\mathscr{X}$.

To simultaneously deal with the reduced and non-reduced cases we introduce the following notation. Let $R_{1}=\{\alpha \in R \mid 2 \alpha \notin R\}, R_{2}=\{\alpha \in$ $\left.R \left\lvert\, \frac{1}{2} \alpha \notin R\right.\right\}$ and $R_{3}=R_{1} \cap R_{2}$. Notice that $R_{1}=R_{2}=R_{3}=R$ if $R$ is reduced. For $\alpha \in R_{2}$, write $q_{\alpha}=q_{i}$ if $|\alpha|=\left|\alpha_{i}\right|$ (if $|\alpha|=\left|\alpha_{i}\right|$ then necessarily $\alpha \in R_{2}$ ). Since $q_{j}=q_{i}$ whenever $s_{j}=w s_{i} w^{-1}$ for some $w \in W$ (see [15, Corollary 2.2]), it follows that $q_{i}=q_{j}$ whenever $\left|\alpha_{i}\right|=\left|\alpha_{j}\right|$, and so the definition of $q_{\alpha}$ is unambiguous. Note that $R=R_{3} \cup\left(R_{1} \backslash R_{3}\right) \cup\left(R_{2} \backslash R_{3}\right)$ where the union is disjoint. Define a set of numbers $\left\{\tau_{\alpha}\right\}_{\alpha \in R}$ related to the numbers $\left\{q_{\alpha}\right\}_{\alpha \in R_{2}}$ by

$$
\tau_{\alpha}= \begin{cases}q_{\alpha} & \text { if } \alpha \in R_{3} \\ q_{0} & \text { if } \alpha \in R_{1} \backslash R_{3} \\ q_{\alpha} q_{0}^{-1} & \text { if } \alpha \in R_{2} \backslash R_{3}\end{cases}
$$

It is convenient to define $\tau_{\alpha}=1$ if $\alpha \notin R$. Note that $\tau_{\alpha}=q_{\alpha}$ if $R$ is reduced (and many subsequent formulae will simplify in this case).

If $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$we write $u^{\lambda}$ in place of $u(\lambda)$. The homomorphism $r \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$defined by

$$
\begin{equation*}
r^{\lambda}=\prod_{\alpha \in R^{+}} \tau_{\alpha}^{\frac{1}{2}\langle\lambda, \alpha\rangle} \quad \text { for all } \lambda \in P \tag{1.4}
\end{equation*}
$$

plays an important role. By [16, Proposition 1.5] and [16, Proposition A.1] we have

$$
\begin{equation*}
N_{\lambda}=N_{\lambda^{*}}=\frac{W_{0}\left(q^{-1}\right)}{W_{0 \lambda}\left(q^{-1}\right)} r^{2 \lambda} \tag{1.5}
\end{equation*}
$$

where $W_{0 \lambda}=\left\{w \in W_{0} \mid w \lambda=\lambda\right\}$ and where $X\left(q^{-1}\right)=\sum_{w \in X} q_{w}^{-1}$ for subsets $X \subseteq W_{0}$.

For $w \in W_{0}$ and $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$we write $w u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$for the homomorphism with $(w u)^{\lambda}=u^{w \lambda}$ for all $\lambda \in P$. Following [10, Chapter IV], for $\lambda \in P^{+}$and $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$we define the Macdonald spherical function $P_{\lambda}(u)$ by

$$
\begin{align*}
P_{\lambda}(u) & =\frac{r^{-\lambda}}{W_{0}\left(q^{-1}\right)} \sum_{w \in W_{0}} c(w u) u^{w \lambda} \quad \text { where } \\
c(u) & =\prod_{\alpha \in R^{+}} \frac{1-\tau_{\alpha}^{-1} \tau_{\alpha / 2}^{-1 / 2} u^{-\alpha^{\vee}}}{1-\tau_{\alpha / 2}^{-1 / 2} u^{-\alpha^{\vee}}} \tag{1.6}
\end{align*}
$$

provided that the denominators of the $c(w u)$ functions do not vanish. Since $P_{\lambda}(u)$ is a Laurent polynomial (see $[16,(1.8)]$ ), these singular cases can
be obtained from the general formula by taking an appropriate limit (see Lemma 2.9 for one example).

For $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$, let $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ be the linear map with $h_{u}\left(A_{\lambda}\right)=$ $P_{\lambda}(u)$ for each $\lambda \in P^{+}$. By [16, Proposition 2.1] every algebra homomorphism $h: \mathscr{A} \rightarrow \mathbb{C}$ is of the form $h=h_{u}$ for some $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$, and $h_{u^{\prime}}=h_{u}$ if and only if $u^{\prime}=w u$ for some $w \in W_{0}$. We call the formula $h_{u}\left(A_{\lambda}\right)=P_{\lambda}(u)$ the Macdonald formula for the algebra homomorphisms $h: \mathscr{A} \rightarrow \mathbb{C}$.

Remark 1.7. - (i) In the situation of homogeneous trees from Remarks 1.5 and $1.6(\mathrm{i})$, if $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$, then writing $z=u^{\lambda_{1}} \in \mathbb{C}^{\times}$ we have

$$
h_{u}\left(A_{k}\right)=\frac{q^{-k / 2}}{1+q^{-1}}\left(\frac{1-q^{-1} z^{-2}}{1-z^{-2}} z^{k}+\frac{1-q^{-1} z^{2}}{1-z^{2}} z^{-k}\right)
$$

provided that $z \neq \pm 1$ (with the values at $z= \pm 1$ found by taking appropriate limits). More generally, in the $\tilde{A}_{n}$ case the functions $P_{\lambda}(u)$ are essentially the Hall-Littlewood polynomials of [12] (see [4]).
(ii) At times the $B C_{n}$ case (see Remark 1.1) requires separate treatment. Recall the description of the parameter system from Figure 1.1. For $u \in$ $\operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$, by writing $t_{i}=u^{e_{i}}$ for $1 \leqslant i \leqslant n$ (noting that in this case $e_{i} \in P$ for each $1 \leqslant i \leqslant n$ ), we have

$$
\begin{aligned}
c(u)= & \left\{\prod_{i=1}^{n} \frac{\left(1-a^{-1} t_{i}^{-1}\right)\left(1+b^{-1} t_{i}^{-1}\right)}{1-t_{i}^{-2}}\right\} \\
& \times\left\{\prod_{1 \leqslant j<k \leqslant n} \frac{\left(1-q_{1}^{-1} t_{j}^{-1} t_{k}\right)\left(1-q_{1}^{-1} t_{j}^{-1} t_{k}^{-1}\right)}{\left(1-t_{j}^{-1} t_{k}\right)\left(1-t_{j}^{-1} t_{k}^{-1}\right)}\right\}
\end{aligned}
$$

where $a=\sqrt{q_{n} q_{0}}$ and $b=\sqrt{q_{n} / q_{0}}$ (see [16, Section 5.2]).
We now recall the second formula for the homomorphisms $h: \mathscr{A} \rightarrow \mathbb{C}$. A sector of $\mathscr{X}$ is a subcomplex $\mathcal{S} \subset \mathscr{X}$ such that there exists an apartment $\mathcal{A}$ with $\mathcal{S} \subset \mathcal{A}$, and an isomorphism $\psi: \mathcal{A} \rightarrow \Sigma$ such that $\psi(\mathcal{S})$ is a sector of $\Sigma$. The base vertex of $\mathcal{S}$ is $\psi^{-1}(\lambda)$, where $\lambda \in P$ is the base vertex of $\psi(\mathcal{S})$. If $\mathcal{S}$ and $\mathcal{S}^{\prime}$ are sectors of $\mathscr{X}$ with $\mathcal{S}^{\prime} \subseteq \mathcal{S}$, then we say that $\mathcal{S}^{\prime}$ is a subsector of $\mathcal{S}$. The boundary $\Omega$ of $\mathscr{X}$ is the set of equivalence classes of sectors, where we declare two sectors to be equivalent if and only if they contain a common subsector. Given $x \in V_{P}$ and $\omega \in \Omega$ there exists a unique sector, denoted $\mathcal{S}^{x}(\omega)$, in the class $\omega$ with base vertex $x$ [17, Lemma 9.7]. For each $x \in V_{P}, \omega \in \Omega$ and $\lambda \in P^{+}$, the intersection $V_{\lambda}(x) \cap \mathcal{S}^{x}(\omega)$ contains exactly one vertex, denoted $v_{\lambda}^{x}(\omega)$. By [16, Theorem 3.4], for each
$\omega \in \Omega$ and $x, y \in V_{P}$ there exists a coweight $h(x, y ; \omega) \in P$ such that

$$
\begin{equation*}
v_{\mu}^{x}(\omega)=v_{\mu-h(x, y ; \omega)}^{y}(\omega) \tag{1.7}
\end{equation*}
$$

for $\mu \in P^{+}$with each $\left\langle\mu, \alpha_{i}\right\rangle, i \in I_{0}$, sufficiently large. Indeed, if $y \in$ $V_{\lambda}(x)$ then (1.7) holds, for all $\omega \in \Omega$, whenever $\mu-\Pi_{\lambda} \subset P^{+}$(see [16, Theorem 3.6]). Here $\Pi_{\lambda} \subset P$ is the saturated set with highest coweight $\lambda$ relative to the partial order on $P$ given by $\mu \preceq \lambda$ if and only if $\lambda-\mu \in Q^{+}$ (recall that $Q^{+}$is the $\mathbb{N}$-span of $\left\{\alpha^{\vee} \mid \alpha \in R^{+}\right\}$). We have

$$
\Pi_{\lambda}=\left\{w \nu \mid \nu \in P^{+}, \nu \preceq \lambda, w \in W_{0}\right\}
$$

(see [8, Lemma 13.4B] for example). The vectors $h(x, y ; \omega)$ are generalisations of the so called horocycle numbers for homogeneous trees.

By [16, Proposition 3.5], for all $\omega \in \Omega$ and all $x, y, z \in V_{P}$ we have the cocycle relation

$$
\begin{equation*}
h(x, y ; \omega)=h(x, z ; \omega)+h(z, y ; \omega) \tag{1.8}
\end{equation*}
$$

Thus $h(x, x ; \omega)=0$ and $h(x, y ; \omega)=-h(y, x ; \omega)$ for all $\omega \in \Omega$ and all $x, y \in V_{P}$.

There is a natural topology on $\Omega$ (discussed in [16]) in which for each $x \in V_{P}$ the sets $\Omega_{x}(y)=\left\{\omega \in \Omega \mid y \in \mathcal{S}^{y}(\omega)\right\}, y \in V_{P}$, form a basis of open and closed sets (this topology is independent of the particular $x \in V_{P}$ chosen). For each $x \in V_{P}$ there is a unique regular Borel probability measure $\nu_{x}$ on $\Omega$ such that $\nu_{x}\left(\Omega_{x}(y)\right)=N_{\lambda}^{-1}$ if $y \in V_{\lambda}(x)$. For $x, x^{\prime} \in V_{P}$ the measures $\nu_{x}$ and $\nu_{x^{\prime}}$ are mutually absolutely continuous with RadonNikodym derivative $\left(d \nu_{x^{\prime}} / d \nu_{x}\right)(\omega)=r^{2 h\left(x, x^{\prime} ; \omega\right)}$ (see [16, Theorem 3.17]).

The integral formula for the algebra homomorphisms $h: \mathscr{A} \rightarrow \mathbb{C}$ is

$$
\begin{equation*}
P_{\lambda}(u)=h_{u}\left(A_{\lambda}\right)=\int_{\Omega}(u r)^{h(x, y ; \omega)} d \nu_{x}(\omega) \tag{1.9}
\end{equation*}
$$

for any $x, y \in V_{P}$ with $y \in V_{\lambda}(x)$. Equality of the Macdonald and integral formulae is proved in [16, Corollary 3.23 and Theorem 6.2].

### 1.7. The Plancherel measure

The Plancherel measure of $\mathscr{A}$ is instrumental in our proof of the local limit theorem. In this section we recall some details about the Plancherel measure and the $\ell^{2}$-spectrum of $\mathscr{A}$ from [10] (see also [16]).

It is easy to see that each $A \in \mathscr{A}$ maps $\ell^{2}\left(V_{P}\right)$ into itself, and for $\lambda \in P^{+}$ and $f \in \ell^{2}\left(V_{P}\right)$ we have $\left\|A_{\lambda} f\right\|_{2} \leqslant\|f\|_{2}$ (see [4, Lemma 4.1] for a proof in a similar context). So we may regard $\mathscr{A}$ as a subalgebra of the $C^{*}$-algebra
$\mathscr{L}\left(\ell^{2}\left(V_{P}\right)\right)$ of bounded linear operators on $\ell^{2}\left(V_{P}\right)$. The facts that $y \in V_{\lambda}(x)$ if and only if $x \in V_{\lambda^{*}}(y)$, and $N_{\lambda^{*}}=N_{\lambda}$, imply that $A_{\lambda}^{*}=A_{\lambda^{*}}$, and so the adjoint $A^{*}$ of any $A \in \mathscr{A}$ is also in $\mathscr{A}$.

Let $\mathscr{A}_{2}$ denote the completion of $\mathscr{A}$ with respect to $\|\cdot\|$, the $\ell^{2}$-operator norm. So $\mathscr{A}_{2}$ is a commutative $C^{*}$-algebra. The algebra homomorphisms $h: \mathscr{A}_{2} \rightarrow \mathbb{C}$ are precisely the extensions $h=\tilde{h}_{u}$ of those algebra homomorphisms $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ which are continuous with respect to the $\ell^{2}$-operator norm. Let us describe the latter homomorphisms.

The analysis here splits into two cases. Following [10, Chapter V] we call the situation where $\tau_{\alpha} \geqslant 1$ for all $\alpha \in R$ the standard case, and the situation where $\tau_{\alpha}<1$ for some $\alpha \in R$ the exceptional case (the use of the word "exceptional" here is unrelated to the so called exceptional root systems). It is immediate from the definition of the numbers $\tau_{\alpha}$ that the exceptional case occurs exactly when $R=B C_{n}$ for some $n \geqslant 1$ and $q_{n}<q_{0}$ (see [16, Lemma 5.1]). In particular, if $R$ is reduced then we are in the standard case.

Let us consider the standard case first. Let

$$
\mathbb{U}=\left\{u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right):\left|u^{\lambda}\right|=1 \text { for all } \lambda \in P\right\} .
$$

In the standard case the algebra homomorphism $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ is continuous with respect to the $\ell^{2}$-operator norm if and only if $u \in \mathbb{U}$ (see [16, Corollary 5.4]). If $h=\tilde{h}_{u}, u \in \mathbb{U}$, we write $\widehat{A}(u)=h(A)$ for $A \in \mathscr{A}_{2}$. In particular, $\widehat{A}_{\lambda}(u)=P_{\lambda}(u)$ for $u \in \mathbb{U}$.

In the standard case, let $\pi$ be the measure on $\mathbb{U}$ given by $d \pi(u)=$ $\frac{W_{0}\left(q^{-1}\right)}{\left|W_{0}\right|}|c(u)|^{-2} d u$, where $d u$ is normalised Haar measure on $\mathbb{U}$ (note that in [16] we write $\pi_{0}$ instead of $\pi$ ). Then for $A \in \mathscr{A}_{2}$ we have

$$
\left(A \delta_{y}\right)(x)=\int_{\mathbb{U}} \widehat{A}(u) \overline{P_{\lambda}(u)} d \pi(u) \quad \text { whenever } y \in V_{\lambda}(x)
$$

where $\delta_{y}(x)=1$ if $x=y$ and $\delta_{y}(x)=0$ otherwise (see [16, Theorem 5.2 and Corollary 5.5]). The measure $\pi$ is essentially the Plancherel measure of $\mathscr{A}$ (more precisely, the Plancherel measure is the image of the measure $\pi$ under the homeomorphism $\left.\varpi: \mathbb{U} / W_{0} \rightarrow \operatorname{Hom}\left(\mathscr{A}_{2}, \mathbb{C}\right), u \mapsto \tilde{h}_{u}\right)$.

Let us consider the exceptional case, and so $R=B C_{n}$ for some $n \geqslant 1$ and $q_{n}<q_{0}$. For $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$, recall the definition of the numbers $t_{i}=$ $t_{i}(u), 1 \leqslant i \leqslant n$, from Remark 1.7(ii). We use the isomorphism $\mathbb{U} \rightarrow \mathbb{T}^{n}$, $u \mapsto\left(t_{1}, \ldots, t_{n}\right)$ to identify $\mathbb{U}$ with $\mathbb{T}^{n}$ (here $\mathbb{T}=\{t \in \mathbb{C}:|t|=1\}$ ). Define $\mathbb{U}^{\prime}=\{-b\} \times \mathbb{T}^{n-1}$, and write $U=\mathbb{U} \cup \mathbb{U}^{\prime}$ (recall from Remark 1.7(ii) that $\left.b=\sqrt{q_{n} / q_{0}}\right)$. Let $d t=d t_{1} \cdots d t_{n}$, where $d t_{i}$ is normalised Haar measure
on $\mathbb{T}$. Let $\phi_{0}(u)=c(u) c\left(u^{-1}\right)$, and let

$$
\phi_{1}(u)=\lim _{t_{1} \rightarrow-b} \frac{\phi_{0}(u)}{1+b^{-1} t_{1}} \quad \text { and } \quad d t^{\prime}=d \delta_{-b}\left(t_{1}\right) d t_{2} \cdots d t_{n}
$$

Note that this limit exists since there is a factor $1+b^{-1} t_{1}$ in $c\left(u^{-1}\right)$ (see Remark 1.7(ii)).

In the exceptional case, let $\pi$ be the measure on $U=\mathbb{U} \cup \mathbb{U}^{\prime}$ given by $d \pi(u)=\frac{W_{0}\left(q^{-1}\right)}{\left|W_{0}\right|} \frac{d t}{\phi_{0}(u)}$ on $\mathbb{U}$ and $d \pi(u)=\frac{W_{0}\left(q^{-1}\right)}{\left|W_{0}^{\prime}\right|} \frac{d t^{\prime}}{\phi_{1}(u)}$ on $\mathbb{U}^{\prime}$, where $W_{0}^{\prime}$ is the Coxeter group $C_{n-1}$ (with $C_{1}=A_{1}$ and $C_{0}=\{1\}$ ). Then for all $A \in \mathscr{A}_{2}$,

$$
\left(A \delta_{y}\right)(x)=\int_{U} \widehat{A}(u) \overline{P_{\lambda}(u)} d \pi(u) \quad \text { whenever } y \in V_{\lambda}(x)
$$

(see [16, Theorem 5.7 and Corollary 5.8]).
To conveniently state formulae in both the standard and exceptional cases simultaneously, we write $U=\mathbb{U}$ in the standard case and (as above) $U=\mathbb{U} \cup \mathbb{U}^{\prime}$ in the exceptional case. Thus (in all cases), for $A \in \mathscr{A}_{2}$,

$$
\begin{equation*}
\left(A \delta_{y}\right)(x)=\int_{U} \widehat{A}(u) \overline{P_{\lambda}(u)} d \pi(u) \quad \text { whenever } y \in V_{\lambda}(x) \tag{1.10}
\end{equation*}
$$

Remark 1.8. - The form of the Plancherel measure in the exceptional case requires that $q_{1} b \geqslant 1$, which follows from a theorem of D. Higman since the numbers $q_{i}, i \in I$, are the parameters of a building (see [16, Lemma 5.6]). We note that for the hypergroups associated to the $B C_{n}$ case the Plancherel measure is supported on $U=\mathbb{U} \cup \mathbb{U}^{\prime} \cup \mathbb{U}^{\prime \prime} \cup \cdots$, where there are $k$ components, with $k$ defined by $q_{1}^{k-1} b \geqslant 1>q_{1}^{k-2} b$. See [10, Theorem 5.2.10].

## 2. The Local Limit Theorem

The basic approach for the local limit theorem is as follows. Let $A$ be the transition operator for an isotropic random walk with matrix $(p(x, y))_{x, y \in V_{P}}$, as in (0.1). Then

$$
\begin{equation*}
p^{(k)}(x, y)=\left(A^{k} \delta_{y}\right)(x) \quad \text { for all } x, y \in V_{P} \text { and } k \in \mathbb{N} . \tag{2.1}
\end{equation*}
$$

Since $\|A\| \leqslant 1$, we may regard $A$ as in $\mathscr{A}_{2}$ and so $h_{u}(A), u \in U$, is defined. Writing $\widehat{A}(u)=h_{u}(A)$ for $u \in U$, we have $\widehat{A}_{\lambda}(u)=P_{\lambda}(u)$ and so

$$
\begin{equation*}
\widehat{A}(u)=\sum_{\lambda \in P^{+}} a_{\lambda} P_{\lambda}(u) . \tag{2.2}
\end{equation*}
$$

By (2.1) and (1.10) we have

$$
\begin{equation*}
p^{(k)}(x, y)=\int_{U}(\widehat{A}(u))^{k} \overline{P_{\lambda}(u)} d \pi(u) \quad \text { whenever } y \in V_{\lambda}(x) \tag{2.3}
\end{equation*}
$$

and we will prove the local limit theorem by determining the asymptotic behaviour of the integral in (2.3) as $k \rightarrow \infty$.

Lemma 2.1. - Let $\lambda \in P^{+}, \lambda \neq 0, x \in V_{P}$, and $y \in V_{\lambda}(x)$. Then
(i) there exists $z \in V_{\lambda}(x) \cap V_{\tilde{\alpha} \vee}(y)$, and
(ii) with $z$ as in (i), there exists $\omega \in \Omega$ such that $h(y, z ; \omega)=\tilde{\alpha}^{\vee}$.

Proof. - Note first that if $c$ and $d$ are distinct $i$-adjacent chambers, $i \in I_{P}$, with type $i$ vertices $u$ and $v$ respectively, then $v \in V_{\tilde{\alpha}^{\vee}}(u)$ (and $\left.u \in V_{\tilde{\alpha}}{ }^{v}(v)\right)$. To see this, let $\mathcal{A}$ be any apartment containing $c$ and $d$, and let $\psi: \mathcal{A} \rightarrow \Sigma$ be a type rotating isomorphism such that $\psi(u)=0$ and $\psi(c)=C_{0}$. Since $\psi(d)$ is 0 -adjacent to $\psi(c)$ we have $\psi(d)=s_{\tilde{\alpha} ; 1}\left(C_{0}\right)$, and so $\psi(v)=s_{\tilde{\alpha} ; 1}(0)=\tilde{\alpha}^{\vee}$. Thus $v \in V_{\tilde{\alpha}^{\vee}}(u)$.

Part (i) now follows exactly as in [5, Lemma 5.1]; we include the proof for completeness. Let $\mathcal{A}$ be an apartment containing $x$ and $y$, and let $c_{0}, c_{1}, \ldots, c_{m}$ be a gallery (that is, a sequence of adjacent chambers with $c_{i-1} \neq c_{i}$ for $1 \leqslant i \leqslant m$ ) with $x \in c_{0}$ and $y \in c_{m}$ and $m$ minimal. Let $\pi$ be the panel $c_{m} \backslash\{y\}$, and let $c^{\prime}$ be the chamber of $\mathcal{A}$ with $c^{\prime} \neq c_{m}$ and $\pi \subset c^{\prime}$ (so $c^{\prime}=c_{m-1}$ if $m \geqslant 1$ ). Let $H$ be the wall of $\mathcal{A}$ determined by $\pi$, and let $\mathcal{A}^{+}$be the half apartment of $\mathcal{A}$ determined by $H$ containing $c^{\prime}$. By thickness there exists a chamber $d \neq c^{\prime}, c_{m}$ with $\pi \subset d$, and writing $z$ for the vertex in $d \backslash \pi$ we have $z \in V_{\tilde{\alpha} \vee} \vee(y)$ by the above discussion. We now show that $z \in V_{\lambda}(x)$. By the proof of [17, Lemma 9.4] there exists an apartment $\mathcal{B}$ containing $\mathcal{A}^{+} \cup d$. Let $\rho_{\mathcal{A}, c^{\prime}}$ be the retraction of $\mathscr{X}$ onto $\mathcal{A}$ with center $c^{\prime}$ (see [3, §IV.3]), and so the map $\varphi=\left.\rho_{\mathcal{A}, c^{\prime}}\right|_{\mathcal{B}}: \mathcal{B} \rightarrow \mathcal{A}$ is a type preserving isomorphism with $\varphi(x)=x$ and $\varphi(z)=y\left(\right.$ since $\left.\varphi(d)=c_{m}\right)$. Since $y \in V_{\lambda}(x)$ there exists a type rotating isomorphism $\psi: \mathcal{A} \rightarrow \Sigma$ with $\psi(x)=0$ and $\psi(y)=\lambda$ (see [15, Proposition 5.6(iii)]), and so the map $\phi=\psi \circ \varphi: \mathcal{B} \rightarrow \Sigma$ is a type rotating isomorphism with $\phi(x)=0$ and $\phi(z)=\lambda$. Thus $z \in V_{\lambda}(x)$.

Part (ii) is a consequence of the following fact. Let $u, v \in V_{P}$ with $v \in$ $V_{\lambda}(u)$. Then there exists $\omega \in \Omega$ such that $h(u, v ; \omega)=\lambda$. To see this, let $\mathcal{A}$ be any apartment containing $u$ and $v$, and let $\psi: \mathcal{A} \rightarrow \Sigma$ be a type rotating isomorphism such that $\psi(u)=0$ and $\psi(v)=\lambda$. Let $\omega$ be the class of $\psi^{-1}\left(\mathcal{S}_{0}\right)$. Since $\psi^{-1}\left(\mathcal{S}_{0}\right)=\mathcal{S}^{u}(\omega)$ and $\psi^{-1}\left(\lambda+\mathcal{S}_{0}\right)=\mathcal{S}^{v}(\omega)$, we have $\psi^{-1}(\mu)=v_{\mu}^{u}(\omega)=v_{\mu-\lambda}^{v}(\omega)$ for sufficiently large $\mu \in P^{+}$, and so $h(u, v ; \omega)=\lambda$.

Recall that $\mathbb{U}=\left\{u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right):\left|u^{\lambda}\right|=1\right.$ for all $\left.\lambda \in P\right\}$. Let

$$
\mathbb{U}_{Q}=\left\{u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right) \mid u^{\gamma}=1 \text { for all } \gamma \in Q\right\} .
$$

Thus $\mathbb{U}_{Q}$ is isomorphic to the dual of the finite abelian group $P / Q$, and so $\mathbb{U}_{Q} \cong P / Q$. Hence $\mathbb{U}_{Q}$ is finite, and $\mathbb{U}_{Q} \subset \mathbb{U}$.

Proposition 2.2. - The set $W_{0} \tilde{\alpha}^{\vee}$ spans $Q$ over $\mathbb{Z}$. Thus if $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$and $u^{w \tilde{\alpha}^{\vee}}=1$ for all $w \in W_{0}$, then $u \in \mathbb{U}_{Q}$.

Proof. - Let $Q^{\prime}$ denote the $\mathbb{Z}$-span of $W_{0} \tilde{\alpha}^{\vee}$. We show that $R^{\vee} \subset Q^{\prime}$, from which it follows that $Q=Q^{\prime}$, hence the result. Suppose first that $R$ is reduced, and let $\beta \in R$. By [2, VI, $\S 1, \mathrm{~N}^{\circ} 3$, Proposition 11] all roots of a given length are conjugate under $W_{0}$, and so if $|\beta|=|\tilde{\alpha}|$ then $\beta^{\vee} \in Q^{\prime}$. Suppose that $|\beta| \neq|\tilde{\alpha}|$. Since at most 2 root lengths occur in $R$ (see [8, Lemma 10.4.C]), and since $R$ is irreducible, there exists $v, v^{\prime} \in W_{0}$ such that $\left\langle v \tilde{\alpha}, v^{\prime} \beta\right\rangle \neq 0$ (for otherwise $W_{0} \tilde{\alpha} \cup W_{0} \beta$ is a partition of $R$ into nonempty pairwise orthogonal sets). Thus $\langle w \tilde{\alpha}, \beta\rangle \neq 0$, where $w=v^{\prime-1} v$, and so by $\left[2, \mathrm{VI}, \S 1, \mathrm{~N}^{\circ} 8\right.$, Proposition $\left.25(\mathrm{iv})\right]$ we have $\left\langle w \tilde{\alpha}^{\vee}, \beta\right\rangle=1$ or $\left\langle w \tilde{\alpha}^{\vee}, \beta\right\rangle=-1$, depending on if $w^{-1} \beta \in R^{+}$or $w^{-1} \beta \in R^{-}$. Since $s_{\beta}\left(w \tilde{\alpha}^{\vee}\right)=w \tilde{\alpha}^{\vee}-\left\langle w \tilde{\alpha}^{\vee}, \beta\right\rangle \beta^{\vee}$ we have $\beta^{\vee} \in Q^{\prime}$.

Finally, if $R$ is non-reduced, then $Q(R)=Q\left(R_{1}\right)$ and $W_{0}(R)=W_{0}\left(R_{1}\right)$. Since $\tilde{\alpha}$ is also the highest root of the reduced root system $R_{1}$ (with respect to the natural base), we have $Q^{\prime}(R)=Q^{\prime}\left(R_{1}\right)$, and so $Q^{\prime}=Q$ in all cases.

Remark 2.3. - Since $\tilde{\alpha}$ is a long root [8, Lemma 10.4.D] (with the convention that all roots are called long if there is only one root length), Proposition 2.2 is true whenever $\tilde{\alpha}$ is replaced with an arbitrary long root $\alpha$ (for $W_{0} \alpha=W_{0} \tilde{\alpha}$ ). However for general $\alpha \in R$ the proposition fails, despite the fact that $W_{0} \alpha^{\vee}$ spans $E$ (by [8, Lemma 10.4.B]). For example let $R$ be the standard $B_{2}$ root system, and take $\alpha=e_{2}$. Then the $\mathbb{Z}$-span of $W_{0} \alpha^{\vee}$ is $2 \mathbb{Z}^{2}$, whereas $Q=\mathbb{Z}^{2}$.

As usual, if $u, v \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$, define $u v \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$by $(u v)^{\lambda}=$ $u^{\lambda} v^{\lambda}$ for all $\lambda \in P$.

Lemma 2.4. - Let $u \in \mathbb{U}$ and $\lambda \in P^{+}$. Then $\left|P_{\lambda}(u)\right| \leqslant P_{\lambda}(1)$, and equality holds for $\lambda \neq 0$ if and only if $u \in \mathbb{U}_{Q}$. Moreover, if $u_{0} \in \mathbb{U}_{Q}$ then $P_{\lambda}\left(u_{0} u\right)=u_{0}^{\lambda} P_{\lambda}(u)$ for all $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$.

Proof. - (cf. [5, Lemma 5.3]) Let $x, y \in V_{P}$ be any vertices with $y \in$ $V_{\lambda}(x)$. The inequality is clear from the integral formula (1.9). Suppose equality holds for some $\lambda \neq 0$. Write $f(\omega)$ for the integrand in (1.9).

Then $f$ is a continuous function on $\Omega$ and $f(\omega) \neq 0$ for all $\omega \in \Omega$. So $\left|\int_{\Omega} f(\omega) d \nu_{x}(\omega)\right|=\int_{\Omega}|f(\omega)| d \nu_{x}(\omega)$ implies that $f(\omega) /|f(\omega)|$ is constant, since $\nu_{x}(O)>0$ for all non-empty open sets $O \subset \Omega$. Thus $u^{h(x, y ; \omega)}$ takes the constant value $P_{\lambda}(u) / P_{\lambda}(1)$ for all $\omega \in \Omega$. Let $z$ be as in Lemma 2.1(i). Since the value of the integral in (1.9) is unchanged if $y$ is replaced by $z$, it follows that $u^{h(x, y ; \omega)}=u^{h(x, z ; \omega)}$ for all $\omega \in \Omega$. Choosing $\omega \in \Omega$ as in Lemma 2.1(ii) and using the cocycle relations we have $u^{\tilde{\alpha}^{\vee}}=u^{h(y, z ; \omega)}=1$. Furthermore, since the value of the integral in (1.9) is unchanged if $u$ is replaced by $w u$ for any $w \in W_{0}$, then $u^{w \tilde{\alpha}^{\vee}}=1$ for all $w \in W_{0}$. It follows from Proposition 2.2 that $u \in \mathbb{U}_{Q}$.

Conversely, if $u_{0} \in \mathbb{U}_{Q}$ and $y \in V_{\lambda}(x)$, then $u_{0}^{h(x, y ; \omega)}=u_{0}^{\lambda}$ for all $\omega \in \Omega$, because $\lambda-h(x, y ; \omega) \in Q$ (see [16, Theorem 3.4(ii)]). Thus it follows from (1.9) that $P_{\lambda}\left(u_{0} u\right)=u_{0}^{\lambda} P_{\lambda}(u)$ for all $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$. In particular, $\left|P_{\lambda}\left(u_{0}\right)\right|=P_{\lambda}(1)$.

In the following series of estimates we will write $C$ for a positive constant, whose value may vary from line to line.

For each $\omega \in \Omega, x, y \in V_{P}$ and $1 \leqslant j \leqslant n$, define

$$
h_{j}(x, y ; \omega)=\left\langle h(x, y ; \omega), \alpha_{j}\right\rangle
$$

Lemma 2.5. - Let $x \in V_{P}$ and $\lambda \in P^{+}$. Then $|h(x, y ; \omega)| \leqslant|\lambda|$ and $\left|h_{j}(x, y ; \omega)\right| \leqslant C|\lambda|$ for all $\omega \in \Omega$, all $y \in V_{\lambda}(x)$, and all $j=1, \ldots, n$.

Proof. - Recall from [16, Theorem 3.4(ii)] that $h(x, y ; \omega) \in \Pi_{\lambda}$ for all $\omega \in \Omega$ and $y \in V_{\lambda}(x)$. By [13, (2.6.2)] we have that $\Pi_{\lambda} \subset \operatorname{conv}\left(W_{0} \lambda\right)$ (the usual convex hull in $E$ here), and since $|w \lambda|=|\lambda|$ for all $w \in W_{0}$, this implies that $|h(x, y ; \omega)| \leqslant|\lambda|$ for all $\omega \in \Omega$ and for all $y \in V_{\lambda}(x)$. We have $\left|\left\langle h(x, y ; \omega), \alpha_{j}\right\rangle\right| \leqslant|h(x, y ; \omega)|\left|\alpha_{j}\right|$, proving the final claim.

Remark 2.6. - There is a natural graph with vertex set $V_{P}$ and vertices $x, y \in V_{P}$ joined by an edge if and only if $y \in V_{\lambda_{i}}(x)$ for some $i \in I_{0}$. In this graph we have $d(x, y)=\sum_{i=1}^{n}\left\langle\lambda, \alpha_{i}\right\rangle$ if $y \in V_{\lambda}(x)$. Lemma 2.5 shows that $|h(x, y ; \omega)|$ and $\left|h_{j}(x, y ; \omega)\right|$ are bounded by $C d(x, y)$.

Notation. - Let $\theta_{1}, \ldots, \theta_{n} \in \mathbb{R}$ and write $\theta=\theta_{1} \alpha_{1}+\cdots+\theta_{n} \alpha_{n}$ (so $\theta \in E)$. Write $e^{i \theta}$ for the element of $\operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$with $\left(e^{i \theta}\right)^{\lambda}=e^{i\langle\lambda, \theta\rangle}$ for all $\lambda \in P^{+}$. With this notation (1.9) gives

$$
\begin{equation*}
P_{\lambda}\left(e^{i \theta}\right)=\int_{\Omega} r^{h(x, y ; \omega)} e^{i\langle h(x, y ; \omega), \theta\rangle} d \nu_{x}(\omega) \quad \text { for all } y \in V_{\lambda}(x) \tag{2.4}
\end{equation*}
$$

and since $P_{\lambda}\left(w^{-1} e^{i \theta}\right)=P_{\lambda}\left(e^{i \theta}\right)$ for all $w \in W_{0}$, it follows that

$$
\begin{equation*}
P_{\lambda}\left(e^{i \theta}\right)=\int_{\Omega} r^{h(x, y ; \omega)} e^{i\langle h(x, y ; \omega), w \theta\rangle} d \nu_{x}(\omega) \quad \text { for all } w \in W_{0}, y \in V_{\lambda}(x) \tag{2.5}
\end{equation*}
$$

Corollary 2.7. - For all $\lambda \in P^{+}, P_{\lambda}\left(e^{i \theta}\right)=P_{\lambda}(1)\left(1+E_{\lambda}(\theta)\right)$, where $\left|E_{\lambda}(\theta)\right| \leqslant|\lambda||\theta|$.

Proof. - We have

$$
\left|P_{\lambda}\left(e^{i \theta}\right)-P_{\lambda}(1)\right| \leqslant \int_{\Omega} r^{h(x, y ; \omega)}\left|e^{i\langle h(x, y ; \omega), \theta\rangle}-1\right| d \nu_{x}(\omega)
$$

and the result follows from Lemma 2.5 since $\left|e^{i z}-1\right| \leqslant|z|$ for all $z \in \mathbb{R}$.
Let $\lambda \in P^{+}$and $y \in V_{\lambda}(x)$. For each $1 \leqslant j, k \leqslant n$ define

$$
\begin{equation*}
b_{j, k}^{\lambda}=\frac{1}{2} \int_{\Omega} h_{j}(x, y ; \omega) h_{k}(x, y ; \omega) r^{h(x, y ; \omega)} d \nu_{x}(\omega) . \tag{2.6}
\end{equation*}
$$

This is independent of the particular pair $x, y \in V_{P}$ with $y \in V_{\lambda}(x)$, for by (2.4) we have

$$
\left.\frac{\partial^{2}}{\partial \theta_{j} \partial \theta_{k}} P_{\lambda}\left(e^{i \theta}\right)\right|_{\theta=0}=-\int_{\Omega} h_{j}(x, y ; \omega) h_{k}(x, y ; \omega) r^{h(x, y ; \omega)} d \nu_{x}(\omega)
$$

(Indeed any expression $\int_{\Omega} p\left(h_{1}(x, y ; \omega), \ldots, h_{n}(x, y ; \omega)\right) r^{h(x, y ; \omega)} d \nu_{x}(\omega)$, where $p$ is a polynomial, is independent of the particular pair $x, y \in V_{P}$ with $\left.y \in V_{\lambda}(x)\right)$.

Lemma 2.8. - Let $\lambda \in P^{+}$, and $\theta_{1}, \ldots, \theta_{n} \in \mathbb{R}$, and as usual write $\theta=\theta_{1} \alpha_{1}+\cdots+\theta_{n} \alpha_{n}$. Then

$$
\begin{equation*}
P_{\lambda}\left(e^{i \theta}\right)=P_{\lambda}(1)-\sum_{j, k=1}^{n} b_{j, k}^{\lambda} \theta_{j} \theta_{k}+R_{\lambda}(\theta) \tag{2.7}
\end{equation*}
$$

where $\left|R_{\lambda}(\theta)\right| \leqslant C|\lambda|^{3}|\theta|^{3} P_{\lambda}(1)$. Furthermore, $\sum_{j, k=1}^{n} b_{j, k}^{\lambda} \theta_{j} \theta_{k} \geqslant 0$, and when $\lambda \neq 0$, equality holds if and only if $\theta=0$.

Proof. - For $\varphi \in \mathbb{R}$ we have $e^{i \varphi}=1+i \varphi-\frac{1}{2} \varphi^{2}+R(\varphi)$ where $|R(\varphi)| \leqslant$ $\frac{1}{6}|\varphi|^{3}$. Applying this to $\varphi=\langle h(x, y ; \omega), \theta\rangle$ and using (2.4) we have

$$
\begin{aligned}
P_{\lambda}\left(e^{i \theta}\right)= & P_{\lambda}(1)+i \int_{\Omega}\langle h(x, y ; \omega), \theta\rangle r^{h(x, y ; \omega)} d \nu_{x}(\omega) \\
& -\frac{1}{2} \int_{\Omega}\langle h(x, y ; \omega), \theta\rangle^{2} r^{h(x, y ; \omega)} d \nu_{x}(\omega)+R_{\lambda}(\theta),
\end{aligned}
$$

where $\left|R_{\lambda}(\theta)\right| \leqslant \frac{1}{6}|\langle h(x, y ; \omega), \theta\rangle|^{3} P_{\lambda}(1) \leqslant \frac{1}{6}|h(x, y ; \omega)|^{3}|\theta|^{3} P_{\lambda}(1)$. The bound for $\left|R_{\lambda}(\theta)\right|$ follows from Lemma 2.5.

We claim that for all $j=1, \ldots, n$ and for all $y \in V_{P}$,

$$
\int_{\Omega} h_{j}(x, y ; \omega) r^{h(x, y ; \omega)} d \nu_{x}(\omega)=0
$$

To see this, let $j \in\{1, \ldots, n\}$ and set $\theta=\theta_{j} \alpha_{j}$ (that is, $\theta_{k}=0$ for all $k \neq j$ ). By differentiating (2.5) with respect to $\theta_{j}$, and then evaluating at $\theta_{j}=0$, firstly with $w=1$ and secondly with $w=s_{j}$, we see that

$$
\int_{\Omega} h_{j}(x, y ; \omega) r^{h(x, y ; \omega)} d \nu_{x}(\omega)=-\int_{\Omega} h_{j}(x, y ; \omega) r^{h(x, y ; \omega)} d \nu_{x}(\omega)
$$

proving the claim. Thus (2.7) holds, and $\sum_{j, k=1}^{n} b_{j, k}^{\lambda} \theta_{j} \theta_{k} \geqslant 0$. If equality holds, then

$$
\int_{\Omega}\langle h(x, y ; \omega), \theta\rangle^{2} r^{h(x, y ; \omega)} d \nu_{x}(\omega)=0 .
$$

Thus $\langle h(x, y ; \omega), \theta\rangle=0$ for almost all $\omega \in \Omega$, and thus for all $\omega \in \Omega$. Thus, since $\langle h(x, y ; \omega), t \theta\rangle=0$ for all $t \in \mathbb{R}$ and $\omega \in \Omega$, we have $P_{\lambda}\left(e^{i(t \theta)}\right)=P_{\lambda}(1)$ for all $t \in \mathbb{R}$ by (2.4), and so $e^{i(t \theta)} \in \mathbb{U}_{Q}$ for all $t \in \mathbb{R}$ by Lemma 2.4. Thus $\theta=0$ since $\left|\mathbb{U}_{Q}\right|<\infty$.

Lemma 2.9. - There exists a polynomial $p\left(x_{1}, \ldots, x_{n}\right)$ of degree at most $M$ such that

$$
\begin{equation*}
P_{\lambda}(1)=r^{-\lambda} p\left(\left\langle\lambda, \alpha_{1}\right\rangle, \ldots,\left\langle\lambda, \alpha_{n}\right\rangle\right) \tag{2.8}
\end{equation*}
$$

for all $\lambda \in P^{+}$, where $M>0$ is some integer depending only on the underlying root system. Furthermore, (by thickness) there exists some $q>1$ such that

$$
\begin{equation*}
P_{\lambda}(1) \leqslant C(|\lambda|+1)^{M} q^{-|\lambda|} \tag{2.9}
\end{equation*}
$$

Proof. - Assuming that $u^{-\alpha^{\vee}} \neq 1$ for any $\alpha \in R_{2}^{+}$, by (1.6) and the definition of the numbers $\tau_{\alpha}$ we have

$$
\begin{equation*}
c(u)=\prod_{\alpha \in R_{2}^{+}} \frac{\left(1-\tau_{2 \alpha}^{-1} \tau_{\alpha}^{-1 / 2} u^{-\alpha^{\vee} / 2}\right)\left(1+\tau_{\alpha}^{-1 / 2} u^{-\alpha^{\vee} / 2}\right)}{1-u^{-\alpha^{\vee}}} \tag{2.10}
\end{equation*}
$$

Write $\sigma=\lambda_{1}+\cdots+\lambda_{n}$. It follows from [2, VI, § 3, $\mathrm{N}^{\circ} 3$, Proposition 2] that

$$
\prod_{\alpha \in R_{2}^{+}}\left(1-u^{-w \alpha^{\vee}}\right)=(-1)^{\ell(w)} u^{\sigma-w \sigma} \prod_{\alpha \in R_{2}^{+}}\left(1-u^{-\alpha^{\vee}}\right)
$$

for all $w \in W_{0}$, and so by (1.6) and (2.10) we have

$$
\begin{equation*}
P_{\lambda}(u)=r^{-\lambda} \frac{F(\lambda)}{\prod_{\alpha \in R_{2}^{+}}\left(1-u^{-\alpha^{v}}\right)} \tag{2.11}
\end{equation*}
$$

where $F(\lambda)$ equals $\frac{1}{W_{0}\left(q^{-1}\right)}$ times

$$
\sum_{w \in W_{0}}\left\{(-1)^{\ell(w)} u^{w \lambda+w \sigma-\sigma} \prod_{\alpha \in R_{2}^{+}}\left(1-\tau_{2 \alpha}^{-1} \tau_{\alpha}^{-1 / 2} u^{-w \alpha^{\vee} / 2}\right)\left(1+\tau_{\alpha}^{-1 / 2} u^{-w \alpha^{\vee} / 2}\right)\right\}
$$

We know that $P_{\lambda}(u)$ is a Laurent polynomial in $u_{1}, \ldots, u_{n}$, and so (2.8) follows from (2.11) by repeated applications of L'Hôpital's rule. The inequality (2.9) follows from Proposition 3.3(ii) and the proof of Proposition 3.3(iv) in Section 3.

Let $A$ be as in (0.1) and $\widehat{A}(u)=h_{u}(A)$ be as in (2.2). It follows from Lemma 2.5 that $\left|b_{j, k}^{\lambda}\right| \leqslant C|\lambda|^{2} P_{\lambda}(1)$, and thus the inequality (2.9) implies that $\sum_{\lambda \in P^{+}} a_{\lambda} b_{j, k}^{\lambda}$ is absolutely convergent for each $1 \leqslant j, k \leqslant n$. We define

$$
\begin{equation*}
b_{j, k}=\frac{1}{\widehat{A}(1)} \sum_{\lambda \in P^{+}} a_{\lambda} b_{j, k}^{\lambda} \tag{2.12}
\end{equation*}
$$

Corollary 2.10. - Let $A$ be as in (0.1), and let $\theta_{1}, \ldots, \theta_{n} \in \mathbb{R}$. Then

$$
\widehat{A}\left(e^{i \theta}\right)=\widehat{A}(1)\left(1-\sum_{j, k=1}^{n} b_{j, k} \theta_{j} \theta_{k}+R(\theta)\right)
$$

where $\sum_{j, k=1}^{n} b_{j, k} \theta_{j} \theta_{k}>0$ unless $\theta=0$, and where $|R(\theta)| \leqslant C|\theta|^{3}$.
Proof. - This follows from Lemma 2.8, using (2.9) to bound $R(\theta)$.
Lemma 2.11. - Let $\theta_{1}, \ldots, \theta_{n} \in \mathbb{R}$. Then

$$
\frac{1}{\left|c\left(e^{i \theta}\right)\right|^{2}}=\prod_{\alpha \in R_{2}^{+}} \frac{\left\langle\alpha^{\vee}, \theta\right\rangle^{2}}{\left(1-\tau_{2 \alpha}^{-1} \tau_{\alpha}^{-1 / 2}\right)^{2}\left(1+\tau_{\alpha}^{-1 / 2}\right)^{2}}\left(1+E_{\alpha}(\theta)\right)
$$

where $\left|E_{\alpha}(\theta)\right| \leqslant C\left\langle\alpha^{\vee}, \theta\right\rangle^{2}$ for each $\alpha \in R_{2}^{+}$.
Proof. - Observe that for $x \in \mathbb{R}$ and $p>1$

$$
\begin{equation*}
\left|\frac{1-e^{-i x}}{1-p^{-1} e^{-i x}}\right|^{2}=\frac{x^{2}}{\left(1-p^{-1}\right)^{2}}\left(1+E_{1}(x)\right), \tag{2.13}
\end{equation*}
$$

where $\left|E_{1}(x)\right| \leqslant C x^{2}$, and for $p>0$

$$
\begin{equation*}
\left|\frac{1+e^{-i x}}{1+p^{-1} e^{-i x}}\right|^{2}=\frac{4}{\left(1+p^{-1}\right)^{2}}\left(1+E_{2}(x)\right) \tag{2.14}
\end{equation*}
$$

where $\left|E_{2}(x)\right| \leqslant C x^{2}$. The result follows by using (2.10), (2.13) and (2.14).

Let $\mathbb{U}_{A}=\{u \in \mathbb{U}:|\widehat{A}(u)|=\widehat{A}(1)\}$. This set will play a role in the local limit theorem, and in the conditions for irreducibility and aperiodicity of the random walk. The following lemma gives a description of $\mathbb{U}_{A}$ in terms of the coefficients $a_{\lambda}$ appearing in (0.1).

Lemma 2.12. - We have

$$
\mathbb{U}_{A}=\left\{u \in \mathbb{U}_{Q} \mid u^{\mu}=u^{\nu} \text { for all } \mu, \nu \in P^{+} \text {with } a_{\mu}, a_{\nu}>0\right\} .
$$

If $u_{0} \in \mathbb{U}_{A}$ then $\widehat{A}\left(u_{0} u\right)=u_{0}^{\mu} \widehat{A}(u)$ for all $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$and all $\mu \in P^{+}$ such that $a_{\mu}>0$.

Proof. - For $u \in \mathbb{U}$ we have

$$
\begin{equation*}
|\widehat{A}(u)| \leqslant \sum_{\lambda \in P^{+}} a_{\lambda}\left|P_{\lambda}(u)\right| \leqslant \sum_{\lambda \in P^{+}} a_{\lambda} P_{\lambda}(1)=\widehat{A}(1) . \tag{2.15}
\end{equation*}
$$

Suppose that $|\widehat{A}(u)|=\widehat{A}(1)$. If $u=u_{0} \in \mathbb{U}_{A}$, then since equality must hold in the second inequality in (2.15) we have $\left|P_{\lambda}\left(u_{0}\right)\right|=P_{\lambda}(1)$ whenever $a_{\lambda}>0$. Since we assume that $a_{\lambda}>0$ for at least one nonzero $\lambda \in P^{+}$we have $u_{0} \in \mathbb{U}_{Q}$ by Lemma 2.4. Thus by Lemma 2.4 we have $P_{\lambda}\left(u_{0}\right)=u_{0}^{\lambda} P_{\lambda}(1)$ for all $\lambda \in P^{+}$, and so since equality must hold in the first inequality in (2.15) we have $u_{0}^{\mu}=u_{0}^{\nu}$ whenever $a_{\mu}, a_{\nu}>0$, proving that

$$
\mathbb{U}_{A} \subseteq\left\{u \in \mathbb{U}_{Q} \mid u^{\mu}=u^{\nu} \text { for all } \mu, \nu \in P^{+} \text {with } a_{\mu}, a_{\nu}>0\right\} .
$$

Conversely, if $u_{0} \in \mathbb{U}_{Q}$ and $u_{0}^{\mu}=u_{0}^{\nu}$ for all $\mu, \nu \in P^{+}$with $a_{\mu}, a_{\nu}>0$, then by Lemma 2.4 we see that $\widehat{A}\left(u_{0} u\right)=u_{0}^{\mu} \widehat{A}(u)$ for all $u \in \mathbb{U}$ and any $\mu \in P^{+}$such that $a_{\mu}>0$, and so taking $u=1$ we have $\left|\widehat{A}\left(u_{0}\right)\right|=\widehat{A}(1)$, so $u_{0} \in \mathbb{U}_{A}$.

For $k \in \mathbb{N}$ and $\lambda \in P^{+}$let

$$
I_{k, \lambda}=\int_{\mathbb{U}}(\widehat{A}(u))^{k} \overline{P_{\lambda}(u)} d \pi(u)
$$

If $y \in V_{\lambda}(x)$, then by (2.3)

$$
p^{(k)}(x, y)= \begin{cases}I_{k, \lambda} & \text { in the standard case, and }  \tag{2.16}\\ I_{k, \lambda}+I_{k, \lambda}^{\prime} & \text { in the exceptional case }\end{cases}
$$

where

$$
\begin{equation*}
I_{k, \lambda}^{\prime}=\int_{\mathbb{U}^{\prime}}(\widehat{A}(u))^{k} \overline{P_{\lambda}(u)} d \pi(u) \tag{2.17}
\end{equation*}
$$

Thus to give an asymptotic formula for $p^{(k)}(x, y)$ we need to give estimates for $I_{k, \lambda}$ and $I_{k, \lambda}^{\prime}$.

Given $\epsilon>0$ and $u_{0} \in \mathbb{U}$, let $N_{\epsilon}\left(u_{0}\right)=\left\{u \in \mathbb{U}:\left|u^{\lambda_{i}}-u_{0}^{\lambda_{i}}\right|<\epsilon\right.$ for all $\left.i \in I_{0}\right\}$. Since $\left|\mathbb{U}_{A}\right|<\infty$ we may choose $\epsilon>0$ sufficiently small so that

$$
\begin{equation*}
N_{\epsilon}\left(u_{0}\right) \cap N_{\epsilon}\left(u_{0}^{\prime}\right)=\emptyset \quad \text { whenever } u_{0}, u_{0}^{\prime} \in \mathbb{U}_{A} \text { are distinct. } \tag{2.18}
\end{equation*}
$$

Write $N_{\epsilon}=N_{\epsilon}(1)$ and $N_{\epsilon}\left(\mathbb{U}_{A}\right)=\bigcup_{u_{0} \in \mathbb{U}_{A}} N_{\epsilon}\left(u_{0}\right)$.
Define $\rho_{1}=\rho_{1}(\epsilon)=\sup \left\{|\widehat{A}(u)| / \widehat{A}(1): u \in \mathbb{U} \backslash N_{\epsilon}\left(\mathbb{U}_{A}\right)\right\}$, and so $0<\rho_{1}<1$. Let

$$
I_{k, \lambda}^{\epsilon}=\int_{N_{\epsilon}}(\widehat{A}(u))^{k} \overline{P_{\lambda}(u)} d \pi(u)
$$

Lemma 2.13. - Fix $\mu \in P^{+}$such that $a_{\mu}>0$, and let $\epsilon>0$ satisfy (2.18). If $u_{0}^{k \mu}=u_{0}^{\lambda}$ for all $u_{0} \in \mathbb{U}_{A}$, then

$$
I_{k, \lambda}=\left|\mathbb{U}_{A}\right| I_{k, \lambda}^{\epsilon}+\mathcal{O}\left(\rho_{1}^{k} \widehat{A}(1)^{k}\right)
$$

Otherwise, $I_{k, \lambda}=0$.
Proof. - It is clear from the formula for $c(u)$ that $c\left(u_{0} u\right)=c(u)$ for all $u_{0} \in \mathbb{U}_{Q}$ and $u \in \mathbb{U}$. Thus by Lemmas 2.4 and 2.12, if $u_{0} \in \mathbb{U}_{A}$ we have

$$
\begin{equation*}
I_{k, \lambda}=u_{0}^{k \mu-\lambda} \int_{\mathbb{U}}\left(\widehat{A}\left(u_{0}^{-1} u\right)\right)^{k} \overline{P_{\lambda}\left(u_{0}^{-1} u\right)} d \pi\left(u_{0}^{-1} u\right)=u_{0}^{k \mu-\lambda} I_{k, \lambda} . \tag{2.19}
\end{equation*}
$$

This shows that $I_{k, \lambda}=0$ if there exists $u_{0} \in \mathbb{U}_{A}$ such that $u_{0}^{k \mu-\lambda} \neq 1$.
Suppose now that $u_{0}^{k \mu-\lambda}=1$ for all $u_{0} \in \mathbb{U}_{A}$. It is clear that

$$
\begin{equation*}
I_{k, \lambda}=\int_{N_{\epsilon}\left(\mathbb{U}_{A}\right)}(\widehat{A}(u))^{k} \overline{P_{\lambda}(u)} d \pi(u)+\mathcal{O}\left(\rho_{1}^{k} \widehat{A}(1)^{k}\right) \tag{2.20}
\end{equation*}
$$

and since $N_{\epsilon}\left(u_{0}\right)=u_{0} N_{\epsilon}$, the calculation in (2.19) shows that for each $u_{0} \in \mathbb{U}_{A}$,

$$
\int_{N_{\epsilon}\left(u_{0}\right)}(\widehat{A}(u))^{k} \overline{P_{\lambda}(u)} d \pi(u)=u_{0}^{k \mu-\lambda} \int_{N_{\epsilon}}(\widehat{A}(u))^{k} \overline{P_{\lambda}(u)} d \pi(u)=I_{k, \lambda}^{\epsilon}
$$

since $u_{0}^{k \mu-\lambda}=1$. The result follows from (2.20) by the choice of $\epsilon$.
It is clear from Corollary 2.10 that if each $\left|\theta_{j}\right|, j=1, \ldots, n$, is sufficiently small, then

$$
\begin{array}{rlr}
\widehat{A}\left(e^{i \theta}\right) & =\widehat{A}(1) e^{-\sum_{i, j=1}^{n} b_{i, j} \theta_{i} \theta_{j}+G(\theta)} \quad \text { where } \\
G(\theta) & =o\left(\sum_{i, j=1}^{n} b_{i, j} \theta_{i} \theta_{j}\right) \tag{2.21}
\end{array}
$$

Writing $\delta=2 \sin ^{-1}(\epsilon / 2)$ we have $N_{\epsilon}=\left\{e^{i \theta}:\left|\theta_{j}\right|<\delta\right.$ for $\left.j=1, \ldots, n\right\}$, and so we may choose $\epsilon>0$ sufficiently small so that

$$
\begin{equation*}
|G(\theta)| \leqslant \frac{1}{2} \sum_{i, j=1}^{n} b_{i, j} \theta_{i} \theta_{j} \tag{2.22}
\end{equation*}
$$

whenever $e^{i \theta} \in N_{\epsilon}$ and $\left|\theta_{j}\right| \leqslant \pi$ for $j=1, \ldots, n$.
Define constants $K_{1}, K_{2}$ and $K_{3}$ by $K_{1}=W_{0}\left(q^{-1}\right)\left|W_{0}\right|^{-1}(2 \pi)^{-n}$,

$$
\begin{align*}
K_{2} & =\prod_{\alpha \in R_{2}^{+}}\left(1-\tau_{2 \alpha}^{-1} \tau_{\alpha}^{-1 / 2}\right)^{-2}\left(1+\tau_{\alpha}^{-1 / 2}\right)^{-2} \\
K_{3} & =\int_{\mathbb{R}^{n}} e^{-\sum_{i, j=1}^{n} b_{i, j} \varphi_{i} \varphi_{j}} \prod_{\alpha \in R_{2}^{+}}\left\langle\alpha^{\vee}, \varphi\right\rangle^{2} d \varphi_{1} \cdots d \varphi_{n} \tag{2.23}
\end{align*}
$$

where $\varphi=\varphi_{1} \alpha_{1}+\cdots+\varphi_{n} \alpha_{n}$.
Lemma 2.14. - Let $\epsilon>0$ be such that (2.18) and (2.22) hold. Then

$$
I_{k, \lambda}^{\epsilon}=K P_{\lambda}(1) \widehat{A}(1)^{k} k^{-\left|R_{2}^{+}\right|-n / 2}\left(1+\mathcal{O}\left(k^{-1 / 2}\right)\right)
$$

where $K=K_{1} K_{2} K_{3}$.
Proof. - Let $\delta=2 \sin ^{-1}(\epsilon / 2)$ as above. We have

$$
I_{k, \lambda}^{\epsilon}=K_{1} \int_{-\delta}^{\delta} \ldots \int_{-\delta}^{\delta}\left(\widehat{A}\left(e^{i \theta}\right)\right)^{k} \frac{P_{\lambda}\left(e^{-i \theta}\right)}{\left|c\left(e^{i \theta}\right)\right|^{2}} d \theta_{1} \cdots d \theta_{n}
$$

and so by making the change of variable $\varphi_{j}=\sqrt{k} \theta_{j}$ for each $j=1, \ldots, n$ we see that

$$
\begin{equation*}
I_{k, \lambda}^{\epsilon}=K_{1} k^{-n / 2} \int_{-\sqrt{k} \delta}^{\sqrt{k} \delta} \cdots \int_{-\sqrt{k} \delta}^{\sqrt{k} \delta}\left(\widehat{A}\left(e^{i \varphi / \sqrt{k}}\right)\right)^{k} \frac{P_{\lambda}\left(e^{-i \varphi / \sqrt{k}}\right)}{\left|c\left(e^{i \varphi / \sqrt{k}}\right)\right|^{2}} d \varphi_{1} \cdots d \varphi_{n} \tag{2.24}
\end{equation*}
$$

where $\varphi=\varphi_{1} \alpha_{1}+\cdots+\varphi_{n} \alpha_{n}$.
By Corollary 2.7 we have

$$
P_{\lambda}\left(e^{-i \varphi / \sqrt{k}}\right)=P_{\lambda}(1)\left(1+E_{1}(\varphi)\right) \quad \text { where } \quad\left|E_{1}(\varphi)\right| \leqslant \frac{|\lambda||\varphi|}{\sqrt{k}}
$$

and it follows from Lemma 2.11 that

$$
\frac{1}{\left|c\left(e^{i \varphi / \sqrt{k}}\right)\right|^{2}}=K_{2} k^{-\left|R_{2}^{+}\right|}\left(1+E_{2}(\varphi)\right) \prod_{\alpha \in R_{2}^{+}}\left\langle\alpha^{\vee}, \varphi\right\rangle^{2}
$$

where $\left|E_{2}(\varphi)\right| \leqslant k^{-1} p\left(\varphi_{1}, \ldots, \varphi_{n}\right)$ for some polynomial $p\left(x_{1}, \ldots, x_{n}\right)$. Using these estimates (along with (2.21)) in (2.24), we see that $I_{k, \lambda}^{\epsilon}$ equals

$$
\begin{aligned}
& K_{1} K_{2} P_{\lambda}(1) \widehat{A}(1)^{k} k^{-\left|R_{2}^{+}\right|-n / 2} \text { times } \\
& \qquad \begin{array}{l}
\int_{X} e^{-\sum_{i, j=1}^{n} b_{i, j} \varphi_{i} \varphi_{j}+k G(\varphi / \sqrt{k})}\left(\prod_{\alpha \in R_{2}^{+}}\left\langle\alpha^{\vee}, \varphi\right\rangle^{2}\right) \\
\\
\times\left(1+E_{1}(\varphi)\right)\left(1+E_{2}(\varphi)\right) d \varphi_{1} \cdots d \varphi_{n}
\end{array}
\end{aligned}
$$

where $X=[-\sqrt{k} \delta, \sqrt{k} \delta]^{n}$. By (2.22), the above integrand is bounded by

$$
e^{-\frac{1}{2} \sum_{i, j=1}^{n} b_{i, j} \varphi_{i} \varphi_{j}}\left(\prod_{\alpha \in R_{2}^{+}}\left\langle\alpha^{\vee}, \varphi\right\rangle^{2}\right)\left(1+\frac{|\lambda||\varphi|}{\sqrt{k}}\right)\left(1+\frac{p\left(\varphi_{1}, \ldots, \varphi_{n}\right)}{k}\right)
$$

and the lemma follows by the Dominated Convergence Theorem.
Lemma 2.15. - Let $\lambda \in P^{+}$and $k \in \mathbb{N}$. In the exceptional case, there exists $0<\rho_{2}<1$ such that

$$
\int_{\mathbb{U}^{\prime}}(\widehat{A}(u))^{k} \overline{P_{\lambda}(u)} d \pi(u)=\mathcal{O}\left(\rho_{2}^{k} \widehat{A}(1)^{k}\right)
$$

Proof. - Let us sketch the proof of this result. The details are given in [14, Appendix B.3]. Since we are in the exceptional case, we have $R=$ $B C_{n}$ for some $n \geqslant 1$ and $q_{n}<q_{0}$. Use the isomorphism $\operatorname{Hom}\left(P, \mathbb{C}^{\times}\right) \rightarrow$ $\left(\mathbb{C}^{\times}\right)^{n}, u \mapsto\left(t_{1}, \ldots, t_{n}\right)$, where $t_{i}=u^{e_{i}}$, to identify $\operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$with $\left(\mathbb{C}^{\times}\right)^{n}$ (and so $\mathbb{U}$ is identified with $\mathbb{T}^{n}$ ). Recall that $\mathbb{U}^{\prime}$ consists of those $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$such that $t_{1}=-\sqrt{q_{n} / q_{0}}$ and $t_{j} \in \mathbb{T}$ for $2 \leqslant j \leqslant n$. Write $\xi_{t}=\left(-\sqrt{q_{n} / q_{0}}, t_{2}, \ldots, t_{n}\right)$, and set $t_{j}=e^{i \theta_{j}}$ for $2 \leqslant j \leqslant n$.

We claim that $\left|P_{\lambda}\left(\xi_{t}\right)\right|<P_{\lambda}(1)$ for all $\lambda \neq 0$ and all $t_{2}, \ldots, t_{n} \in \mathbb{T}$, from which the result follows (since $\widehat{A}$ is continuous on $\mathbb{U}^{\prime}$, and $\mathbb{U}^{\prime}$ is compact). The first step is to explicitly compute $P_{\lambda_{1}}(u)$ for arbitrary $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$. By [14, Lemma B.3.2] we have
$P_{\lambda_{1}}(u)=N_{\lambda_{1}}^{-1}\left(\left(q_{0}-1\right)\left(1+q_{1}+\cdots+q_{1}^{n-1}\right)+\sqrt{q_{0} q_{n}} q_{1}^{n-1} \sum_{j=1}^{n}\left(t_{j}+t_{j}^{-1}\right)\right)$
for all $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$. From this formula we deduce that $\left|P_{\lambda_{1}}\left(\xi_{t}\right)\right|<$ $P_{\lambda_{1}}(1)$ for all $t_{2}, \ldots, t_{n} \in \mathbb{T}$ (see [14, Theorem B.3.3]). We now use this fact to show that $\left|P_{\lambda}\left(\xi_{t}\right)\right|<P_{\lambda}(1)$ for all $\lambda \neq 0$ and all $t_{2}, \ldots, t_{n} \in \mathbb{T}$.

Recall from [15, Corollary 5.22] that the operators $\left\{A_{\lambda}\right\}_{\lambda \in P^{+}}$satisfy

$$
A_{\lambda} A_{\mu}=\sum_{\nu \in P^{+}} a_{\lambda, \mu ; \nu} A_{\nu}
$$

where

$$
a_{\lambda, \mu ; \nu}=\frac{N_{\nu}}{N_{\lambda} N_{\mu}}\left|V_{\lambda}(x) \cap V_{\mu^{*}}(y)\right| \geqslant 0
$$

and where $x, y \in V_{P}$ is any pair with $y \in V_{\nu}(x)$. Since $\lambda_{1}=\tilde{\alpha}^{\vee}$ here, an analogous argument to that given in Lemma 2.1(i) shows that $a_{\lambda, \lambda ; \lambda_{1}}>0$ for all $\lambda \neq 0$ (see [14, Lemma B.3.4]).

Since the algebra homomorphisms $h_{\xi_{t}}: \mathscr{A} \rightarrow \mathbb{C}$ are continuous with respect to the $\ell^{2}$-operator norm, and since $\left\|A_{\mu}\right\|=P_{\mu}(1)$ for all $\mu \in P^{+}$ (see $\left[16\right.$, Theorem 6.3]), we have $\left|P_{\mu}\left(\xi_{t}\right)\right| \leqslant P_{\mu}(1)$ for all $\mu \in P^{+}$. Hence for $\lambda \neq 0$,

$$
\left|P_{\lambda}\left(\xi_{t}\right)\right|^{2}=\left|h_{\xi_{t}}\left(A_{\lambda}^{2}\right)\right| \leqslant \sum_{\mu \in P^{+}} a_{\lambda, \lambda ; \mu}\left|P_{\mu}\left(\xi_{t}\right)\right|<\sum_{\mu \in P^{+}} a_{\lambda, \lambda ; \mu} P_{\mu}(1)=P_{\lambda}(1)^{2}
$$

where we have used the facts that $\left|P_{\lambda_{1}}\left(\xi_{t}\right)\right|<P_{\lambda_{1}}(1)$ and $a_{\lambda, \lambda ; \lambda_{1}}>0$.
We now give our local limit theorem.
Theorem 2.16. - Let $y \in V_{\lambda}(x)$ and $k \in \mathbb{N}$, and suppose that $a_{\mu}>0$. If $u_{0}^{k \mu}=u_{0}^{\lambda}$ for all $u_{0} \in \mathbb{U}_{A}$, then

$$
p^{(k)}(x, y)=\left|\mathbb{U}_{A}\right| K P_{\lambda}(1) \widehat{A}(1)^{k} k^{-\left|R_{2}^{+}\right|-n / 2}\left(1+\mathcal{O}\left(k^{-1 / 2}\right)\right),
$$

where $K$ is as in Lemma 2.14. If $u_{0}^{k \mu} \neq u_{0}^{\lambda}$ for some $u_{0} \in \mathbb{U}_{A}$, then $p^{(k)}(x, y)=0$.

Proof. - In the standard case the result follows from (2.16) and Lemmas 2.13 and 2.14. In the exceptional case, $Q=P$, and so $\mathbb{U}_{Q}=\{1\}$, and so $\mathbb{U}_{A}=\{1\}$. The result now follows from (2.16) and Lemmas 2.13, 2.14, and 2.15 .

A random walk on a state-space $X$ is called irreducible if for each pair $x, y \in X$ there exists $k=k(x, y) \in \mathbb{N}$ such that $p^{(k)}(x, y)>0$. The period of an irreducible random walk is $\mathfrak{p}=\operatorname{gcd}\left\{k \geqslant 1 \mid p^{(k)}(x, x)>0\right\}$, which is independent of $x \in X$ by irreducibility (see [20]). An irreducible random walk is called aperiodic if $\mathfrak{p}=1$.

Corollary 2.17. - Let $A$ be as in (0.1), and suppose that $a_{\mu}>0$. Then
(i) $A$ is irreducible if and only if for each $\lambda \in P^{+}$there exists $k=$ $k(\lambda) \in \mathbb{N}$ such that $u_{0}^{k \mu}=u_{0}^{\lambda}$ for all $u_{0} \in \mathbb{U}_{A}$, and
(ii) $A$ is irreducible and aperiodic if and only if $\left|\mathbb{U}_{A}\right|=1$.

Proof. - First let us note that in the exeptional case it is easy to see that any walk with $a_{\mu}>0$ for some $\mu \neq 0$ is both aperiodic and irreducible, and since $Q=P$ we have $\mathbb{U}_{A}=\{1\}$. So consider the standard case, and suppose that $a_{\mu}>0$. Let $y \in V_{\lambda}(x)$. If $A$ is irreducible, then there exists $k \in \mathbb{N}$ such that $p^{(k)}(x, y)>0$, and so $u_{0}^{k \mu}=u_{0}^{\lambda}$ for all $u_{0} \in \mathbb{U}_{A}$, by (2.16) and Lemma 2.13. Conversely, if for each $\lambda \in P^{+}$there exists $k_{0} \in \mathbb{N}$ such that
$u_{0}^{k_{0} \mu}=u_{0}^{\lambda}$ for all $u_{0} \in \mathbb{U}_{A}$, then writing $r=\left|\mathbb{U}_{A}\right|$ we have $u_{0}^{\left(k_{0}+r l\right) \mu}=u_{0}^{\lambda}$ for all $u_{0} \in \mathbb{U}_{A}$ and all $l \geqslant 0$. As $k \rightarrow \infty$ through the values $k_{0}+r l$, Theorem 2.16 implies irreducibility.

If $\left|\mathbb{U}_{A}\right|=1$ then $A$ is irreducible and aperiodic, by Theorem 2.16. Conversely, if $A$ is irreducible and aperiodic, then

$$
1=\operatorname{gcd}\left\{k \geqslant 1 \mid p^{(k)}(x, x)>0\right\}=\operatorname{gcd}\left\{k \geqslant 1 \mid u_{0}^{k \mu}=1 \text { for all } u_{0} \in \mathbb{U}_{A}\right\}
$$

and so $\mathbb{U}_{A}=\{1\}$.
Remark 2.18. - It is possible to explicitly compute the constant $K_{3}$ from (2.23) (at least in most cases). We refer the reader to [14, Section 8.4] for details. A key step in the calculation is to observe that there is a number $b>0$ such that $b_{j, k}=\left\langle\alpha_{j}, \alpha_{k}\right\rangle b$ for all $1 \leqslant j, k \leqslant n$, and so $K_{3}=$ $b^{-\left|R_{2}^{+}\right|-n / 2} J$ where

$$
J=\int_{\mathbb{R}^{n}} e^{-\sum_{j, k=1}^{n}\left\langle\alpha_{j}, \alpha_{k}\right\rangle \theta_{j} \theta_{k}} \prod_{\alpha \in R_{2}^{+}}\left\langle\alpha^{\vee}, \theta\right\rangle^{2} d \theta_{1} \ldots d \theta_{n}
$$

and $\theta=\theta_{1} \alpha_{1}+\cdots+\theta_{n} \alpha_{n}$. The integral $J$ depends only on the underlying root system, and has been computed using Gram's identity in the cases when $R=B_{n}, C_{n}, D_{n}$ or $B C_{n}$ (there are other techniques using orthogonal polynomials). We have

$$
J= \begin{cases}\pi^{n / 2} 2^{-n(n-1)} \prod_{i=1}^{n}(2 i)! & \text { if } R=B_{n} \text { or } R=B C_{n} \\ \pi^{n / 2} 2^{-n^{2}-n-1} \prod_{i=1}^{n}(2 i)! & \text { if } R=C_{n} \\ \pi^{n / 2} 2^{-n^{2}+n-1} n!\prod_{i=1}^{n-1}(2 i)! & \text { if } R=D_{n} .\end{cases}
$$

When $R=A_{n}$ the integral $J$ may be written as

$$
\int_{\mathbb{R}^{n}} e^{-\left(x_{1}^{2}+\cdots+x_{n+1}^{2}\right)} \prod_{1 \leqslant i<j \leqslant n+1}\left(x_{i}-x_{j}\right)^{2} d x_{1} \cdots d x_{n}
$$

(up to some constant factors), where $x_{n+1}=-\left(x_{1}+\cdots+x_{n}\right)$. We have been unable to compute this integral. In principle the integrals for the $E, F$ and $G$ cases could be explicitly computed using a computer package.

Remark 2.19. - Let us briefly discuss some applications of our local limit theorem to probability measures on groups acting on $\mathscr{X}$. An automorphism $\psi$ of $\mathscr{X}$ is called type rotating if there exists a type rotating automorphism $\sigma$ of the Coxeter graph of $W$ (in the sense of [15, § 4.8]) such that $\tau(\psi(x))=\sigma(\tau(x))$ for all $x \in V$. Suppose that $G$ is a locally compact group acting on $V$ such that $G$ acts transitively on $V_{P}$ and such that for each $x \in V_{P}$ and $g \in G$ the automorphism $x \mapsto g x$ is type rotating. Assume that $K=\{g \in G \mid g o=o\}$ acts transitively on each set $V_{\lambda}(o)$,
$\lambda \in P^{+}$, where $o \in V_{P}$ is some fixed vertex. Let $\varphi$ be the density function of a bi- $K$-invariant probability measure on $G$. Then, exactly as in [5, Lemma 8.1], setting $p(g o, h o)=\varphi\left(g^{-1} h\right)$ for $g, h \in G$ defines an isotropic random walk on $V_{P}$. Since the $k$-th convolution power $\varphi^{(* k)}(g)$ is simply $p^{(k)}(o, g o)$, Theorem 2.16 may immediately be interpreted as a local limit theorem for bi- $K$-invariant probability measures on $G$ (the assumption that $a_{\lambda}>0$ for some $\lambda \neq 0$ simply means that $\varphi$ is not the indicator function on $K$ ).

As an important modification, suppose now that $G$ is a group of type preserving simplicial complex automorphisms acting strongly transitively on $\mathscr{X}$, meaning that $G$ acts transitively on pairs $(\mathcal{A}, c)$ of apartments $\mathcal{A}$ and chambers $c \subset \mathcal{A}$. Fix an apartment $\mathcal{A}_{0}$ and a chamber $c_{0} \subset \mathcal{A}_{0}$. The subgroups $B=\operatorname{stab}_{G}\left(c_{0}\right)$ and $N=\operatorname{stab}_{G}\left(\mathcal{A}_{0}\right)$ form a $B N$-pair in $G$ with associated Weyl group $N /(B \cap N)$ isomorphic to $W$ [17, Theorem 5.2]. Indeed the set of left cosets $\{g B \mid g \in G\}$ defines an affine building (as a chamber system) isomorphic to $\mathscr{X}$, where $g B \sim_{i} h B$ if and only if $g^{-1} h \in B\left\langle s_{i}\right\rangle B$ (where $w B$ means $n B$ for any $n \in N$ with image $w \in W)$. Let $o$ be the type 0 vertex of $c_{0}$. The subgroup $K=\operatorname{stab}_{G}(o)$ of $G$ equals $B W_{0} B=\bigcup_{w \in W_{0}} B w B$ (see [17, Theorem 5.4(iii)]), and since $G$ acts strongly transitively and $B \cap N$ is transitive on the chambers of $\mathcal{A}_{0}$, it follows that $K$ is transitive on each set $V_{\lambda}(o), \lambda \in Q \cap P^{+}$.

Let $\varphi$ be the density function of a bi- $K$-invariant probability measure on $G$. To study convolution powers $\varphi^{(* k)}(g), g \in G$, it is natural to study an associated random walk on $V_{Q}=\left\{x \in V_{P} \mid \tau(x)=0\right\} \subseteq V_{P}$, where we define $p(g o, h o)=\varphi\left(g^{-1} h\right)$ for $g, h \in G$. To apply our local limit theorem we consider these random walks as reducible isotropic random walks on $V_{P}$ by setting $p(x, y)=\varphi\left(g^{-1} h\right)$ if $y \in V_{\lambda}(x)$ and $g o \in V_{\lambda}(h o)$ for some $\lambda \in P^{+}$(necessarily $\lambda \in Q \cap P^{+}$), and $p(x, y)=0$ otherwise. These random walks have the property that $p(x, y)=0$ if $\tau(x) \neq \tau(y)$, and it is simple to see that they are indeed isotropic. Theorem 2.16 is now applicable, and in particular, by taking $\mathscr{X}$ to be the Bruhat-Tits building of a group of $p$-adic type (see Remark 1.6(ii)) we have a local limit theorem for these groups.

Finally we remark that the methods here can be extended to deal with groups acting (in a type rotating fashion) on subsets $V_{L}$ of $V_{P}$. Here $L$ is a lattice in $E$ with $Q \subseteq L \subseteq P$, and $V_{L}=\left\{x \in V_{P} \mid \tau(x) \in I_{L}\right\}$, where $I_{L}=\{\tau(\lambda) \mid \lambda \in L\} \subseteq I_{P}$. Thus $V_{Q} \subseteq V_{L} \subseteq V_{P}$, and our discussion above deals with the extreme cases of $L=P$ and $L=Q$.

## 3. The Rate of Escape Theorem

Let $X$ be any set, and let $P=(p(x, y))_{x, y \in X}$ be a transition probability matrix. Let $X=\bigcup_{i \in I} X_{i}$ be a partition of $X$. We call $P$ factorisable over $I$ if for each $i, j \in I$, the sum

$$
\sum_{y \in X_{j}} p(x, y)
$$

has the same value for all $x \in X_{i}$. In this case we write $\bar{p}(i, j)$ for this value, and let $\bar{P}=(\bar{p}(i, j))_{i, j \in I}$. Clearly $\bar{p}(i, j) \geqslant 0$ for all $i, j \in I$, and for each $i \in I$,

$$
\sum_{j \in I} \bar{p}(i, j)=\sum_{j \in I} \sum_{y \in X_{j}} p(x, y)=\sum_{y \in X} p(x, y)=1
$$

where $x \in X_{i}$. Thus $\bar{P}$ is a transition probability matrix (on $I$ ). Furthermore, if $\left(Z_{k}\right)_{k \in \mathbb{N}}$ is a Markov chain on $X$ with transition probability matrix $P$, then $\left(\bar{Z}_{k}\right)_{k \in \mathbb{N}}$, where $\bar{Z}_{k}=i$ if $Z_{k} \in X_{i}$, defines a Markov chain on $I$ with transition probability matrix $\bar{P}$.

In our setting, consider the partition (for fixed $o \in V_{P}$ and $\omega \in \Omega$ ) $V_{P}=\bigcup_{\lambda \in P} V_{\lambda}$, where

$$
V_{\lambda}=\left\{x \in V_{P} \mid h(o, x ; \omega)=\lambda\right\}
$$

Proposition 3.1. - The matrices (operators) $A_{\lambda}=\left(p_{\lambda}(x, y)\right)_{x, y \in V_{P}}$, $\lambda \in P^{+}$, are factorisable over $P$. Moreover, $\bar{p}_{\lambda}(\mu, \nu)$ does not depend on o or $\omega$, and $\bar{p}_{\lambda}(\mu, \nu)=\bar{p}_{\lambda}(0, \nu-\mu)$.

Proof. - Let $\mu, \nu \in P$ and $x \in V_{\mu}$. By the cocycle relations we have $h(o, y ; \omega)=h(x, y ; \omega)+\mu$ for all $y \in V_{P}$, and so

$$
\begin{align*}
\sum_{y \in V_{\nu}} p_{\lambda}(x, y) & =\frac{1}{N_{\lambda}}\left|\left\{y \in V_{\lambda}(x) \mid h(o, y ; \omega)=\nu\right\}\right|  \tag{3.1}\\
& =\frac{1}{N_{\lambda}}\left|\left\{y \in V_{\lambda}(x) \mid h(x, y ; \omega)=\nu-\mu\right\}\right| .
\end{align*}
$$

It follows from [16, Lemma 3.19] that $A_{\lambda}$ is factorisable, and that $\bar{p}_{\lambda}(\mu, \nu)$ does not depend on $\omega \in \Omega$ or $o \in V_{P}$. The transitional invariance is immediate from (3.1).

Corollary 3.2. - Let $A=(p(x, y))_{x, y \in V_{P}}$ be as in (0.1). Then $A$ is factorisable over $P$. Moreover, for each $\mu, \nu \in P$ we have $\bar{p}(\mu, \nu)=$ $\bar{p}(0, \nu-\mu)$, and this value does not depend on o or $\omega$. Finally, if $\left(Z_{k}\right)_{k \in \mathbb{N}}$ is a Markov chain with transition probability matrix $A$, then $\bar{Z}_{k}=h\left(o, Z_{k} ; \omega\right)$, so that $\bar{p}(\mu, \nu)=\mathbb{P}\left(h\left(o, Z_{k+1} ; \omega\right)=\nu \mid h\left(o, Z_{k} ; \omega\right)=\mu\right)$.

Proof. - The first statements follow from Proposition 3.1 and the elementary fact that a (finite or infinite) convex combination of factorisable transition matrices is again factorisable. The final claim is immediate from the definition of $\bar{Z}_{k}$.

Let $\left\{T_{j}\right\}_{j \in J}$ be a partition of $R_{2}$ according to root length (so $|J|=1$ or $|J|=2$ ). For $j \in J$, let $T_{j}^{+}=R_{2}^{+} \cap T_{j}$, and $B_{j}=B \cap T_{j}$ (so $B=\bigcup_{j \in J} B_{j}$, as $B \subset R_{2}$ ). For each $j \in J$, let

$$
\rho_{j}=\frac{1}{2} \sum_{\alpha \in T_{j}^{+}} \alpha .
$$

Finally, for each $j \in J$ fix some $\beta_{j} \in T_{j}^{+}$.
Proposition 3.3. - With the above definitions:
(i) For $i, j \in J$, if $\alpha \in B_{i}$, then $\left\langle\alpha^{\vee}, \rho_{j}\right\rangle=\delta_{i, j}$. Thus, for each $j \in J$, $\rho_{j} \in \overline{\mathcal{S}}_{0}$.
(ii) Let $\lambda \in P$. Then

$$
r^{\lambda}=\prod_{j \in J}\left(\tau_{\beta_{j}} \tau_{2 \beta_{j}}^{2}\right)^{\left\langle\lambda, \rho_{j}\right\rangle}
$$

(note that this product has at most two factors).
(iii) $r^{w_{0} \lambda}=r^{-\lambda}$ for all $\lambda \in P$ (that is, $r^{\lambda^{*}}=r^{\lambda}$ ).
(iv) If $\lambda \in P$ and $\mu \preceq \lambda$, then $r^{\mu} \leqslant r^{\lambda}$, with equality if and only if $\mu=\lambda$.
(v) For $w \in W_{0}$, we have $c(w r)=\delta_{w, 1} W_{0}\left(q^{-1}\right)$.

Proof. - (i) Since $\alpha \in B, s_{\alpha}$ permutes $R_{2}^{+} \backslash\{\alpha\}$, and since the sets $T_{j}$, $j \in J$, are $W_{0}$-invariant, we have $s_{\alpha}\left(T_{j}^{+}\right)=T_{j}^{+}$if $j \in J \backslash\{i\}$. Thus for any $j \in J$ we have

$$
s_{\alpha}\left(\rho_{j}\right)=\rho_{j}-\delta_{i, j} \alpha,
$$

and so $\left\langle\alpha^{\vee}, \rho_{j}\right\rangle=\delta_{i, j}$. Then $\left\langle\alpha^{\vee}, \rho_{j}\right\rangle \geqslant 0$ for all $\alpha \in B$, and so $\rho_{j} \in \overline{\mathcal{S}}_{0}$ for all $j \in J$.
(ii) By (1.4) and the fact that $R_{1}^{+} \backslash R_{3}^{+}=2\left(R_{2}^{+} \backslash R_{3}^{+}\right)$we calculate

$$
r^{\lambda}=\left(\prod_{\alpha \in R_{3}^{+}} \tau_{\alpha}^{\frac{1}{2}\langle\lambda, \alpha\rangle}\right)\left(\prod_{\beta \in R_{2}^{+} \backslash R_{3}^{+}}\left(\tau_{\beta} \tau_{2 \beta}^{2}\right)^{\frac{1}{2}\langle\lambda, \beta\rangle}\right)
$$

Since $\tau_{2 \alpha}=1$ if $\alpha \in R_{3}^{+}$, and since $\tau_{\beta}=\tau_{\beta_{j}}$ if $\beta \in T_{j}$, it follows that

$$
r^{\lambda}=\prod_{\beta \in R_{2}^{+}}\left(\tau_{\beta} \tau_{2 \beta}^{2}\right)^{\frac{1}{2}\langle\lambda, \beta\rangle}=\prod_{j \in J}\left(\tau_{\beta_{j}} \tau_{2 \beta_{j}}^{2}\right)^{\left\langle\lambda, \rho_{j}\right\rangle} .
$$

(iii) Since $w_{0} \rho_{j}=-\rho_{j}$ for $j \in J$, by (ii) we have $r^{w_{0} \lambda}=r^{-\lambda}$ for all $\lambda \in P$.
(iv) Observe that $\tau_{\alpha} \tau_{2 \alpha}^{2}=q_{\alpha}$ if $\alpha \in R_{3}$, and $\tau_{\alpha} \tau_{2 \alpha}^{2}=q_{0} q_{\alpha}$ if $\alpha \in R_{2} \backslash R_{3}$. Thus, by thickness, $\tau_{\alpha} \tau_{2 \alpha}^{2}>1$ for all $\alpha \in R_{2}$. Since $\rho_{j} \in \overline{\mathcal{S}}_{0}$ for $j \in J$, and since $\mu \preceq \lambda$ implies that $\lambda-\mu \in Q^{+}$, it follows from (ii) that $r^{\lambda-\mu} \geqslant 1$, with equality if and only if $\mu=\lambda$ (for if $\mu \neq \lambda$, then $\left\langle\lambda-\mu, \rho_{j}\right\rangle>0$ for at least one $j \in J$ ).
(v) Observe that if $w \neq 1$, then $w R_{2}^{+} \cap(-B) \neq \emptyset$. To see this, if $\alpha \in R_{2}^{+}$, and if $-\alpha \notin w R_{2}^{+}$, then $-w^{-1} \alpha \notin R_{2}^{+}$, and so $w^{-1} \alpha \in R_{2}^{+}$. It follows that if $w R_{2}^{+} \cap(-B)=\emptyset$, then $w^{-1} B \subset R_{2}^{+}$, and so $w^{-1} R_{2}^{+}=R_{2}^{+}$. Thus $w=1$ (for by [2, VI, $\S 1, \mathrm{~N}^{\circ} 6$, Corollary 2] we have $\ell(w)=\left|\left\{\alpha \in R_{2}^{+} \mid w \alpha \in R_{2}^{-}\right\}\right|$ for all $w \in W_{0}$ ).

Suppose that $w \neq 1$, and take $\alpha \in R_{2}^{+}$such that $w \alpha=-\beta \in-B$. Then by (i) and (ii),

$$
r^{-w \alpha^{\vee}}=r^{\beta^{\vee}}=\tau_{\beta} \tau_{2 \beta}^{2}=\tau_{\alpha} \tau_{2 \alpha}^{2}
$$

and so $1-\tau_{2 \alpha}^{-1} \tau_{\alpha}^{-1 / 2} r^{-w \alpha^{\vee} / 2}=0$. Thus by (2.10) we see that $c(w r)=0$ whenever $w \neq 1$. Since $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ is an algebra homomorphism we have

$$
1=h_{u}\left(A_{0}\right)=\frac{1}{W_{0}\left(q^{-1}\right)} \sum_{w \in W_{0}} c(w u)
$$

for all nonsingular $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$. Evaluating at $u=r$ shows that $c(r)=$ $W_{0}\left(q^{-1}\right)$.

Remark 3.4. - Most of Proposition 3.3 can be found on page 61 of [10]. Notice in particular that Proposition 3.3(v) gives a nice factorisation of the Poincaré polynomial of $W_{0}$, namely $W_{0}\left(q^{-1}\right)=c(r)$ (at least when the $q$ 's come from a building). See also [11].

For each $j \in I_{0}$, let $j^{*} \in I_{0}$ be defined by $-w_{0} \alpha_{j}=\alpha_{j^{*}}$. Note that $\left(j^{*}\right)^{*}=j$.

Corollary 3.5. - Let $x \in V_{P}, \lambda \in P^{+}$, and $y \in V_{\lambda}(x)$.
(i) We have

$$
P_{\lambda}\left(r e^{i \theta}\right)=\int_{\Omega} e^{i\left\langle h(x, y ; \omega), w_{0} \theta\right\rangle} d \nu_{x}(\omega)
$$

where $w_{0}$ is the longest element of $W_{0}$.
(ii) For each $j \in I_{0}$ the integral

$$
\gamma_{j}^{(\lambda)}=\int_{\Omega} h_{j^{*}}(y, x ; \omega) d \nu_{x}(\omega)
$$

is independent of the particular pair $x, y \in V_{P}$ with $y \in V_{\lambda}(x)$ (the $j^{*}$ here makes the statements of the main theorems simpler).

Proof. - By Proposition 3.3(iii) (and the fact that $w_{0}^{-1}=w_{0}$ ) we have

$$
P_{\lambda}\left(r e^{i \theta}\right)=P_{\lambda}\left(w_{0}\left(r e^{i \theta}\right)\right)=P_{\lambda}\left(r^{-1} w_{0}\left(e^{i \theta}\right)\right)=\int_{\Omega} e^{i\left\langle h(x, y ; \omega), w_{0} \theta\right\rangle} d \nu_{x}(\omega)
$$

proving (i).
Since $w_{0} \theta=-\sum_{j=1}^{n} \theta_{j} \alpha_{j^{*}}$, by (i) we have

$$
\begin{equation*}
\left.\frac{\partial}{\partial \theta_{j}} P_{\lambda}\left(r e^{i \theta}\right)\right|_{\theta=0}=i \gamma_{j}^{(\lambda)} \tag{3.2}
\end{equation*}
$$

proving (ii).
The following proposition gives a symmetry property of the numbers $\gamma_{j}^{(\lambda)}$ generalising [5, Proposition 3.5(iii)]. We will not use this result in this paper.

Proposition 3.6. - Let $j \in I_{0}$ and $\lambda \in P^{+}$. We have $\gamma_{j}^{\left(\lambda^{*}\right)}=\gamma_{j^{*}}^{(\lambda)}$.
Proof. - Observe that for $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$and $\lambda \in P^{+}, P_{\lambda^{*}}(u)=$ $P_{\lambda}\left(u^{-1}\right)$. It suffices to prove this for nonsingular $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$. Using Proposition 3.3(iii) we see that if $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$is nonsingular then

$$
c\left(w_{0} u\right)=\prod_{\alpha \in R^{+}} \frac{1-\tau_{\alpha}^{-1} \tau_{\alpha / 2}^{-1 / 2} u^{-w_{0} \alpha^{\vee}}}{1-\tau_{\alpha / 2}^{-1 / 2} u^{-w_{0} \alpha^{\vee}}}=\prod_{\alpha \in R^{+}} \frac{1-\tau_{\alpha}^{-1} \tau_{\alpha / 2}^{-1 / 2} u^{\alpha^{\vee}}}{1-\tau_{\alpha / 2}^{-1 / 2} u^{\alpha^{\vee}}}=c\left(u^{-1}\right)
$$

(we have used the facts that $w_{0} R^{+}=R^{-}$and $\tau_{\alpha}=\tau_{\beta}$ if $\beta \in W_{0} \alpha$ ). Thus

$$
\begin{aligned}
P_{\lambda^{*}}(u) & =\frac{r^{-\lambda^{*}}}{W_{0}\left(q^{-1}\right)} \sum_{w \in W_{0}} c(w u) u^{w \lambda^{*}} \\
& =\frac{r^{-\lambda}}{W_{0}\left(q^{-1}\right)} \sum_{w \in W_{0}} c\left(w w_{0} u\right) u^{w w_{0} \lambda^{*}} \\
& =P_{\lambda}\left(u^{-1}\right)
\end{aligned}
$$

and so by (3.2) we have

$$
\begin{aligned}
\gamma_{j}^{\left(\lambda^{*}\right)} & =-\left.i \frac{\partial}{\partial \theta_{j}} P_{\lambda}\left(r^{-1} e^{-i \theta}\right)\right|_{\theta=0} \\
& =-\left.i \frac{\partial}{\partial \theta_{j}} P_{\lambda}\left(w_{0}\left(r e^{-i w_{0} \theta}\right)\right)\right|_{\theta=0} \\
& =-\left.i \frac{\partial}{\partial \varphi_{j^{*}}} P_{\lambda}\left(r e^{i \varphi}\right)\right|_{\varphi=0}
\end{aligned}
$$

and so $\gamma_{j}^{\left(\lambda^{*}\right)}=\gamma_{j^{*}}^{(\lambda)}$.

Lemma 3.7. - Let $\lambda \in P^{+}$and $j \in I_{0}$. Then

$$
\gamma_{j}^{(\lambda)}=\left\langle\lambda, \alpha_{j}\right\rangle+\mathcal{O}(1) .
$$

Proof. - Let us temporarily write $\rho_{\theta}$ in place of $r e^{i \theta}$. Then for $w \in W_{0}$,

$$
\frac{\partial}{\partial \theta_{j}} c\left(w \rho_{\theta}\right) \rho_{\theta}^{w \lambda}=i\left\langle w \lambda, \alpha_{j}\right\rangle c\left(w \rho_{\theta}\right) \rho_{\theta}^{w \lambda}+\rho_{\theta}^{w \lambda} \frac{\partial}{\partial \theta_{j}} c\left(w \rho_{\theta}\right) .
$$

It follows from Proposition 3.3(v) that for $w \in W_{0}$,

$$
\left.\frac{r^{-\lambda}}{W_{0}\left(q^{-1}\right)} \frac{\partial}{\partial \theta_{j}} c\left(w \rho_{\theta}\right) \rho_{\theta}^{w \lambda}\right|_{\theta=0}= \begin{cases}i\left\langle\lambda, \alpha_{j}\right\rangle+\mathcal{O}(1) & \text { if } w=1 \\ \mathcal{O}(1) & \text { if } w \neq 1\end{cases}
$$

The result follows from (3.2).
Lemma 3.8. - If $\left(Z_{k}\right)_{k \in \mathbb{N}}$ is a Markov chain in $V_{P}$ with $Z_{0}=x$ and transition operator $A_{\lambda}$, then for any $\omega \in \Omega, \mathbb{E}\left(h_{j}\left(Z_{1}, x ; \omega\right)\right)=\gamma_{j^{*}}^{(\lambda)}$.

Proof. - Since $Z_{1} \in V_{\lambda}(x)$ with probability 1 , we have $\int_{\Omega} h_{j}\left(Z_{1}, x ; \omega\right)$ $d \nu_{x}(\omega)=\gamma_{j^{*}}^{(\lambda)}$. As in [5, Proposition 3.5(ii)] we see that we may take expectations under the integral sign, and so $\gamma_{j^{*}}^{(\lambda)}=\int_{\Omega} \mathbb{E}\left(h_{j}\left(Z_{1}, x ; \omega\right)\right) d \nu_{x}(\omega)$. By Corollary 3.2, the distribution of $h_{j}\left(Z_{1}, x ; \omega\right)$, and hence $\mathbb{E}\left(h_{j}\left(Z_{1}, x ; \omega\right)\right)$, is independent of $\omega \in \Omega$. The result follows.

We now prove our rate of escape theorem.
Theorem 3.9. - Let $A$ be as in (0.1), and suppose that

$$
\sum_{\lambda \in P^{+}}|\lambda| a_{\lambda}<\infty .
$$

Let $\left(Z_{k}\right)_{k \in \mathbb{N}}$ be the corresponding Markov chain, and for each $k \in \mathbb{N}$ let $\nu_{k} \in P^{+}$be such that $Z_{k} \in V_{\nu_{k}}(x)$, where $x=Z_{0}$. Then for each $j \in I_{0}$, with probability 1

$$
\frac{1}{k}\left\langle\nu_{k}, \alpha_{j}\right\rangle \rightarrow \gamma_{j} \quad \text { as } k \rightarrow \infty
$$

where $\gamma_{j}=\sum_{\lambda \in P^{+}} a_{\lambda} \gamma_{j}^{(\lambda)}$. That is, $\frac{1}{k} \nu_{k} \rightarrow \gamma_{1} \lambda_{1}+\cdots+\gamma_{n} \lambda_{n}$.
Proof. - Observe first that $\gamma_{j}<\infty$ by Lemma 2.5 and the finite first moment assumption. By Lemma 3.7 we have $\frac{1}{k}\left\langle\nu_{k}, \alpha_{j}\right\rangle=\frac{1}{k} \gamma_{j}^{\left(\nu_{k}\right)}+\mathcal{O}\left(k^{-1}\right)$, and so it suffices to prove that

$$
\int_{\Omega} \frac{h_{j^{*}}\left(Z_{k}, x ; \omega\right)}{k} d \nu_{x}(\omega) \rightarrow \gamma_{j}
$$

with probability 1.
By Corollary 3.2 we see that for each fixed $\omega \in \Omega, h_{j^{*}}\left(Z_{k}, x ; \omega\right)$ is a random variable distributed like a sum of $k$ independent real random variables,
each with the distribution of $h_{j^{*}}\left(Z_{1}, x ; \omega\right)$. Now $\mathbb{E}\left(h_{j^{*}}\left(Z_{1}, x ; \omega\right)\right)=\gamma_{j}$, and so by the classical law of large numbers we have

$$
\frac{h_{j^{*}}\left(Z_{k}, x ; \omega\right)}{k} \rightarrow \gamma_{j}
$$

with probability 1.
By Remark 2.6 and the second part of [20, Proposition 8.8(a)] we see that $h_{j^{*}}\left(Z_{k}, x ; \omega\right) / k$ is bounded with probability 1 . Thus by the Bounded Convergence Theorem we have

$$
\lim _{k \rightarrow \infty} \int_{\Omega} \frac{h_{j^{*}}\left(Z_{k}, x ; \omega\right)}{k} d \nu_{x}(\omega)=\int_{\Omega} \lim _{k \rightarrow \infty} \frac{h_{j^{*}}\left(Z_{k}, x ; \omega\right)}{k} d \nu_{x}(\omega)=\gamma_{j}
$$

with probability 1 , completing the proof.
Corollary 3.10. - The numbers $\gamma_{j}, j=1, \ldots, n$, from Theorem 3.9 are nonnegative.

Proof. - This follows from the rate of escape theorem, since $\nu_{k} \in P^{+}$ for each $k \in \mathbb{N}$.

We can strengthen Corollary 3.10 by applying the local limit theorem.
Theorem 3.11. - For $j=1, \ldots, n$ we have $\gamma_{j}>0$. In particular, by taking $A=A_{\lambda}$ we have $\gamma_{j}^{(\lambda)}>0$ for all $\lambda \neq 0$.

Proof. - We follow the outline given in [5, Remark 4.7]. By our local limit theorem we may choose a pair $(\nu, K) \in P^{+} \times \mathbb{N}$ with $K$ large and each $\left\langle\nu, \alpha_{j}\right\rangle, j=1, \ldots, n$, large, such that $p^{(K)}(x, y)>0$ whenever $y \in V_{\nu}(x)$. With $A$ as in (0.1), write $A^{K}=\sum_{\lambda \in P^{+}} a_{\lambda}^{(K)} A_{\lambda}$, and so $a_{\nu}^{(K)}>0$. For each $j=1, \ldots, n$ let

$$
\begin{align*}
\gamma_{j, K} & =\sum_{\lambda \in P^{+}} a_{\lambda}^{(K)} \gamma_{j}^{(\lambda)}  \tag{3.3}\\
& =-\left.i \frac{\partial}{\partial \theta_{j}} \widehat{A}^{K}\left(r e^{i \theta}\right)\right|_{\theta=0} \tag{3.4}
\end{align*}
$$

where we have used (3.2). Note that $\gamma_{j, K}$ is simply the $\gamma_{j}$ for the transition matrix $A^{K}$.

It follows from (3.4) and (3.2) that $\gamma_{j, K}=K \gamma_{j}$, and from Corollary 3.10 and (3.3) we have $\gamma_{j, K} \geqslant a_{\nu}^{(K)} \gamma_{j}^{(\nu)}$. Now, by Lemma 3.7 we see that each $\gamma_{j}^{(\nu)}, j=1, \ldots, n$, is strictly positive (remember that each component of $\nu$ may be chosen to be large), and thus $\gamma_{j, K} \geqslant a_{\nu}^{(K)} \gamma_{j}^{(\nu)}>0$. Thus $\gamma_{j}=$ $\frac{1}{K} \gamma_{j, K}>0$ for each $j=1, \ldots, n$.

Remark 3.12. - Fix a vertex $o \in V_{P}$, and recall that $\mathcal{S}^{\circ}(\omega)$ denotes the unique sector in the class $\omega$ based at $o$, and that for each $\lambda \in P^{+}$we write $v_{\lambda}^{o}(\omega)$ for the unique vertex in $\mathcal{S}^{o}(\omega) \cap V_{\lambda}(o)$. Given vertices $x, y \in V_{P}$, there is a natural notion of the convex hull conv $\{x, y\}$, as studied in [16, Appendix B]. We say that a sequence $\left(x_{k}\right)_{k \in \mathbb{N}}$ of vertices in $V_{P}$ converges to $\omega \in \Omega$ if for each $\lambda \in P^{+}$there exists $k_{\lambda} \in \mathbb{N}$ such that $v_{\lambda}^{o}(\omega)$ is in conv $\left\{o, x_{k}\right\}$ whenever $k \geqslant k_{\lambda}$. This definition is independent of the $o \in V_{P}$ chosen. Theorem 3.11 shows that for an isotropic random walk $\left(Z_{k}\right)_{k \in \mathbb{N}}$ we have, with probability $1, Z_{k} \rightarrow \omega$ for some random element $\omega \in \Omega$. The key point to observe to show this is that if $Z_{k} \in V_{\nu_{k}}\left(Z_{0}\right)$, then by Theorems 3.9 and $3.11\left\langle\nu_{k}, \alpha_{j}\right\rangle, j=1, \ldots, n$, becomes large as $k \rightarrow \infty$.

Remark 3.13. - We note that the random walk $\left(\bar{Z}_{k}\right)_{k \in \mathbb{N}}$ on $P$ from Corollary 3.2 can be studied using classical methods, since $P \cong \mathbb{Z}^{n}$. In the notation of (A.1), by (3.1) and [16, Lemma 3.19 and Theorem 6.2] we have $\bar{p}_{\lambda}(0, \mu)=r^{-\mu} a_{\lambda, \mu}$. Assuming that $\sum_{\mu \in P}|\mu| \bar{p}(0, \mu)<\infty$, a calculation using (A.1) and (3.2) shows that the mean $\mathfrak{m}=\sum_{\mu \in P} \mu \bar{p}(0, \mu)$ of the random walk $\left(\bar{Z}_{k}\right)_{k \in \mathbb{N}}$ is $\mathfrak{m}=\sum_{j=1}^{n} \gamma_{j^{*}} \lambda_{j}$, where $\gamma_{j}$ is as in Theorem 3.9. A similar calculation shows that the characteristic function for the walk is

$$
\sum_{\mu \in P} \bar{p}(0, \mu) e^{i\langle\mu, \theta\rangle}=\widehat{A}\left(r^{-1} e^{i \theta}\right) .
$$

By Corollary 3.2 this walk is transitionally invariant, and so the usual Fourier inversion (as in [19, § II.6, Proposition 3]) gives

$$
\bar{p}^{(k)}(0, \mu)=\frac{1}{(2 \pi)^{n}} \int_{-\pi}^{\pi} \cdots \int_{-\pi}^{\pi}\left(\widehat{A}\left(r^{-1} e^{i \theta}\right)\right)^{k} e^{-i\langle\mu, \theta\rangle} d \theta_{1} \cdots d \theta_{n}
$$

The asymptotic behaviour may now be extracted using the methods in [20, § III.13] and the calculations in Lemma 4.4.

## 4. The Central Limit Theorem

Lemma 4.1. - Let $\lambda \in P^{+}$. There exists a constant $C$, independent of $\theta$ and $\lambda$, such that

$$
\left|h_{r e^{i \theta}}\left(A_{\lambda}\right)-e^{i\langle\lambda, \theta\rangle}\right| \leqslant C|\theta| .
$$

Proof. - Recall that $r^{w \lambda} \leqslant r^{\lambda}$ and $c(w r)=\delta_{w, 1} W_{0}\left(q^{-1}\right)$ for all $w \in W_{0}$ (see Proposition 3.3). Thus

$$
\left|h_{r e^{i \theta}}\left(A_{\lambda}\right)-e^{i\langle\lambda, \theta\rangle}\right| \leqslant \frac{1}{W_{0}\left(q^{-1}\right)} \sum_{w \in W_{0}}\left|c\left(w\left(r e^{i \theta}\right)\right)-c(w r)\right| .
$$

The result follows since each $c\left(w\left(r e^{i \theta}\right)\right)$ is a smooth function in $\theta_{1}, \ldots, \theta_{n}$.

Lemma 4.2. - (See Theorem A.3) The homomorphisms $h_{r e^{i \theta}}: \mathscr{A} \rightarrow \mathbb{C}$, $\theta \in E$, are bounded.

Proof. - For each $\lambda \in P^{+}$we have $\left|h_{r e^{i \theta}}\left(A_{\lambda}\right)\right| \leqslant 1$ by Corollary 3.5(i).

Let $x \in V_{P}$. The spherical function (with respect to $x$ ) associated to $h_{u}$ is the function $F_{u}^{x}: V_{P} \rightarrow \mathbb{C}$ which for each $\lambda \in P^{+}$takes the constant value $h_{u}\left(A_{\lambda}\right)$ on the set $V_{\lambda}(x)$.

Lemma 4.3. - Let $Z_{0}=x$, and suppose that $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$is such that $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ is bounded. Then $\mathbb{E}\left(F_{u}^{x}\left(Z_{k}\right)\right)=(\widehat{A}(u))^{k}$.

Proof. - We have $A^{k}=\sum_{\lambda \in P^{+}} a_{\lambda}^{(k)} A_{\lambda}$ where $a_{\lambda}^{(k)}=\mathbb{P}\left(Z_{k} \in V_{\lambda}(x)\right)$. Since $F_{u}^{x}\left(Z_{k}\right)=h_{u}\left(A_{\lambda}\right)$ if $Z_{k} \in V_{\lambda}(x)$, we have

$$
\mathbb{E}\left(F_{u}^{x}\left(Z_{k}\right)\right)=\sum_{\lambda \in P^{+}} a_{\lambda}^{(k)} h_{u}\left(A_{\lambda}\right)=h_{u}\left(A^{k}\right)=(\widehat{A}(u))^{k}
$$

where we have used the continuity of $h_{u}$ on the closure of $\mathscr{A}$ in the space of bounded linear operators on $\ell^{1}\left(V_{P}\right)$ to justify the last two equalities.

For $1 \leqslant j, k \leqslant n$ and $\lambda \in P^{+}$, let

$$
\begin{equation*}
\gamma_{j, k}^{(\lambda)}=\int_{\Omega} h_{j^{*}}(y, x ; \omega) h_{k^{*}}(y, x ; \omega) d \nu_{x}(\omega) \tag{4.1}
\end{equation*}
$$

where $x, y \in V_{P}$ are any vertices with $y \in V_{\lambda}(x)$ (as in Corollary 3.5(ii) this is independent of the particular pair $x, y \in V_{P}$ with $y \in V_{\lambda}(x)$ chosen). If we suppose a finite second moment assumption:

$$
\begin{equation*}
\sum_{\lambda \in P^{+}}|\lambda|^{2} a_{\lambda}<\infty \tag{4.2}
\end{equation*}
$$

then for all $1 \leqslant j, k \leqslant n$ we have $\sum_{\lambda \in P^{+}} a_{\lambda} \gamma_{j, k}^{(\lambda)}<\infty$, and we denote this value by $\gamma_{j, k}$.

Lemma 4.4. - Suppose that (4.2) holds. Then with $\gamma_{j}, 1 \leqslant j \leqslant n$ as defined in Theorem 3.9, and $\gamma_{j, k}, 1 \leqslant j, k \leqslant n$ as defined above,

$$
\widehat{A}\left(r e^{i \theta}\right)=1+i \sum_{j=1}^{n} \gamma_{j} \theta_{j}-\frac{1}{2} \sum_{j, k=1}^{n} \gamma_{j, k} \theta_{j} \theta_{k}+o\left(|\theta|^{2}\right)
$$

Furthermore, if $\theta \neq 0$ then $\left(\sum_{j=1}^{n} \gamma_{j} \theta_{j}\right)^{2}<\sum_{j, k=1}^{n} \gamma_{j, k} \theta_{j} \theta_{k}$.

Proof. - Consider the case $A=A_{\lambda}, \lambda \neq 0$. Using Corollary 3.5(i), the elementary result $e^{i \varphi}=1+i \varphi-\frac{1}{2} \varphi^{2}+o\left(\varphi^{2}\right)$ implies that

$$
\widehat{A}_{\lambda}\left(r e^{i \theta}\right)=1+i \sum_{j=1}^{n} \gamma_{j}^{(\lambda)} \theta_{j}-\frac{1}{2} \sum_{j, k=1}^{n} \gamma_{j, k}^{(\lambda)} \theta_{j} \theta_{k}+o(|\lambda \| \theta|),
$$

where we have used Lemma 2.5. The first claim follows.
To deduce the final claim, let $B_{\lambda}=\sum_{j=1}^{n} \gamma_{j}^{(\lambda)} \theta_{j}$ and $C_{\lambda}=\sum_{j, k=1}^{n} \gamma_{j, k}^{(\lambda)} \theta_{j} \theta_{k}$. Then
$B_{\lambda}^{2}=\left(\int_{\Omega} \sum_{j=1}^{n} h_{j^{*}}(y, x ; \omega) d \nu_{x}(\omega)\right)^{2} \leqslant \int_{\Omega}\left(\sum_{j=1}^{n} h_{j^{*}}(y, x ; \omega)\right)^{2} d \nu_{x}(\omega)=C_{\lambda}$,
and

$$
\left(\sum_{j=1}^{n} \gamma_{j} \theta_{j}\right)^{2}=\left(\sum_{\lambda \in P^{+}} a_{\lambda} B_{\lambda}\right)^{2} \leqslant \sum_{\lambda \in P^{+}} a_{\lambda} B_{\lambda}^{2} \leqslant \sum_{\lambda \in P^{+}} a_{\lambda} C_{\lambda}=\sum_{j, k=1}^{n} \gamma_{j, k} \theta_{j} \theta_{k}
$$

To see that the inequality is strict if $\theta \neq 0$, recall that by hypothesis there exists $\lambda \neq 0$ such that $a_{\lambda}>0$. If equality holds in the inequality $B_{\lambda}^{2} \leqslant C_{\lambda}$, then for $y \in V_{\lambda}(x)$,

$$
\left\langle h(x, y ; \omega), w_{0} \theta\right\rangle=\sum_{j=1}^{n} h_{j^{*}}(y, x ; \omega) \theta_{j}
$$

is independent of $\omega \in \Omega$, and thus by Corollary 3.5 (ii) this quantity is independent of the particular pair $x, y \in V_{P}$ with $y \in V_{\lambda}(x)$ too. Choosing $z \in V_{\lambda}(x) \cap V_{\tilde{\alpha}^{\vee}}(y)$ as in Lemma 2.1(i), we have

$$
\left\langle h(y, z ; \omega), w_{0} \theta\right\rangle=\left\langle h(x, z ; \omega), w_{0} \theta\right\rangle-\left\langle h(x, y ; \omega), w_{0} \theta\right\rangle=0 .
$$

By modifying the proof of Lemma 2.1(ii), it is easy to see that for each $w \in W_{0}$ there exists $\omega_{w} \in \Omega$ such that $h\left(y, z ; \omega_{w}\right)=w \tilde{\alpha}^{\vee}$, and thus by the above $\left\langle w \tilde{\alpha}^{\vee}, w_{0} \theta\right\rangle=0$ for all $w \in W_{0}$. Thus $\theta=0$, since $W_{0} \tilde{\alpha}^{\vee}$ spans $E[8$, Lemma 10.4.B].

Let $\Gamma_{1}(\theta)=\sum_{j=1}^{n} \gamma_{j} \theta_{j}$ and $\Gamma_{2}(\theta)=\sum_{j, k=1}^{n} \gamma_{j, k} \theta_{j} \theta_{k}$, and write

$$
\Gamma(\theta)=\Gamma_{2}(\theta)-\Gamma_{1}^{2}(\theta)=\sum_{j, k=1}^{n} g_{j, k} \theta_{j} \theta_{k}
$$

By Lemma 4.4, $\Gamma=\left(g_{j, k}\right)_{j, k=1}^{n}$ is a positive definite matrix.
Theorem 4.5. - Let $A$ be as in (0.1) and suppose that (4.2) holds. As in Theorem 3.9, for each $k \in \mathbb{N}$ let $\nu_{k} \in P^{+}$be such that $Z_{k} \in V_{\nu_{k}}(x)$,
where $x=Z_{0}$. Then

$$
\left(\frac{\left\langle\nu_{k}, \alpha_{1}\right\rangle-\gamma_{1} k}{\sqrt{k}}, \ldots, \frac{\left\langle\nu_{k}, \alpha_{n}\right\rangle-\gamma_{n} k}{\sqrt{k}}\right)
$$

converges in distribution to the normal distribution $N(0, \Gamma)$, with $\Gamma$ as above.

Proof. - Following the proof of the classical Central Limit Theorem (see [19, Proposition II.8] for example), it suffices to show that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \mathbb{E}\left(e^{i\left(\left\langle\nu_{k}, \theta\right\rangle-k \Gamma_{1}(\theta)\right) / \sqrt{k}}\right)=e^{-\frac{1}{2} \Gamma(\theta)} \tag{4.3}
\end{equation*}
$$

By Lemma 4.1 we have

$$
e^{i\left\langle\nu_{k}, \theta\right\rangle / \sqrt{k}}=P_{\nu_{k}}\left(r e^{i \theta / \sqrt{k}}\right)+o\left(k^{-1 / 2}\right)=F_{r e^{i \theta / \sqrt{k}}}^{x}\left(Z_{k}\right)+o\left(k^{-1 / 2}\right),
$$

and so by Lemmas 4.2 and 4.3 we have

$$
\mathbb{E}\left(e^{i\left\langle\nu_{k}, \theta\right\rangle / \sqrt{k}}\right)=\left(\widehat{A}\left(r e^{i \theta / \sqrt{k}}\right)\right)^{k}+o\left(k^{-1 / 2}\right)
$$

Thus using Lemma 4.4 we compute

$$
\begin{aligned}
\mathbb{E}\left(e^{i\left(\left\langle\nu_{k}, \theta\right\rangle-k \Gamma_{1}(\theta)\right) / \sqrt{k}}\right) & =\left(1-\frac{1}{2 k} \Gamma(\theta)+o\left(k^{-1}\right)\right)^{k}+o\left(k^{-1 / 2}\right) \\
& =e^{-\frac{1}{2} \Gamma(\theta)}+o\left(k^{-1 / 2}\right)
\end{aligned}
$$

Thus (4.3) holds, completing the proof.

## Appendix A. Bounded Spherical Functions

It is straightforward to see that each $A \in \mathscr{A}$ maps $\ell^{1}\left(V_{P}\right)$ into itself. Let $\mathscr{A}_{1}$ denote the closure of $\mathscr{A}$ in the space $\mathscr{L}\left(\ell^{1}\left(V_{P}\right)\right)$ of bounded linear operators on $\ell^{1}\left(V_{P}\right)$. Thus $\mathscr{A}_{1}$ is a commutative unital Banach $*$-algebra. The algebra homomorphisms $h: \mathscr{A}_{1} \rightarrow \mathbb{C}$ are precisely the extensions of those algebra homomorphisms $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ which are bounded. In this appendix we determine the $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$for which this holds.

In the notation of Remark 1.6(ii), it is shown in [10, Theorem 4.7.1] that $h_{u}: \mathscr{A}_{Q} \rightarrow \mathbb{C}$ is bounded if and only $\left|u^{w \lambda}\right| \leqslant r^{\lambda}$ for all $\lambda \in Q \cap P^{+}$and all $w \in W_{0}$. The proof given in [10] requires some knowledge of the singular cases (when the denominator of a $c(w u)$ function vanishes). While it should be possible to generalise the proof in [10] to cover the more general setting of homomorphisms $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$, we will provide a different proof which does not require any specific details of the singular cases (instead our proof uses the Plancherel measure).

We restrict our attention to the standard case (where $\tau_{\alpha} \geqslant 1$ for all $\alpha \in R$ ). In the exceptional case (where $\tau_{\alpha}<1$ for some $\alpha \in R$ ) we have $R=B C_{n}$ for some $n \geqslant 1$ and so $Q=P$ and $\mathscr{A}_{Q}=\mathscr{A}$. Thus Macdonald's analysis in [10] covers this specific case.

Remark A.1. - (i) In fact in [10, Theorem 4.7.1] Macdonald proves a geometric analog of the result stated above. For $u \in \operatorname{Hom}\left(Q, \mathbb{C}^{\times}\right)$, identify $\log |u| \in \operatorname{Hom}(Q, \mathbb{R})$ with the unique element $x_{u} \in E$ which satisfies $\left\langle\lambda, x_{u}\right\rangle=\log \left|u^{\lambda}\right|$ for all $\lambda \in Q$. Let $D=\left\{x_{w r} \mid w \in W_{0}\right\}$. Then [10, Theorem 4.7.1] says that $h_{u}: \mathscr{A}_{Q} \rightarrow \mathbb{C}$ is bounded if and only if $x_{u} \in \operatorname{conv}(D)$.
(ii) Note that we have already seen in Lemma 4.2 that the homomorphisms $h_{r e^{i \theta}}: \mathscr{A} \rightarrow \mathbb{C}$ are bounded, and that this was enough information to prove our central limit theorem. It is, of course, still desirable to have the much more accurate Theorem A. 3 below.

For $\lambda \in P^{+}$and $u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right)$, define the monomial symmetric function $m_{\lambda}(u)$ by

$$
m_{\lambda}(u)=\sum_{\mu \in W_{0} \lambda} u^{\mu}
$$

where $W_{0} \lambda=\left\{w \lambda \mid w \in W_{0}\right\}$. By $[16,(6.1)]$ there are numbers $a_{\lambda, \mu}$ such that

$$
\begin{equation*}
P_{\lambda}(u)=\sum_{\mu \in P} a_{\lambda, \mu} u^{\mu}=\sum_{\mu \preceq \lambda} a_{\lambda, \mu} m_{\mu}(u) \tag{A.1}
\end{equation*}
$$

(where the second sum is over those $\mu \in P^{+}$with $\lambda-\mu \in Q^{+}$), and it follows (using [15, Theorem 6.11] for example) that for $\lambda, \mu \in P^{+}$there are numbers $b_{\lambda, \mu}$ such that

$$
\begin{equation*}
m_{\lambda}(u)=\sum_{\mu \preceq \lambda} b_{\lambda, \mu} P_{\mu}(u) . \tag{A.2}
\end{equation*}
$$

Lemma A.2. - Let $\lambda, \mu \in P^{+}$and $\mu \preceq \lambda$. In the standard case there exists a constant $K>0$ independent of $\lambda$ and $\mu$ such that $\left|b_{\lambda, \mu}\right| \leqslant K r^{\mu}$. Thus there is a constant $C>0$ independent of $\lambda \in P^{+}$such that

$$
\sum_{\mu \preceq \lambda}\left|b_{\lambda, \mu}\right| \leqslant C\left|\Pi_{\lambda}\right| r^{\lambda} .
$$

Proof. - Since we assume that we are in the standard case, by [16, Lemma 6.1] we have

$$
b_{\lambda, \mu}=N_{\mu} \int_{\mathbb{U}} m_{\lambda}(u) \overline{P_{\mu}(u)} d \pi(u) .
$$

Using (1.5) and the techniques used to derive [16, (5.2)] we see that

$$
\begin{aligned}
b_{\lambda, \mu} & =\frac{W_{0}\left(q^{-1}\right)}{W_{0 \mu}\left(q^{-1}\right)\left|W_{0}\right|} r^{\mu} \int_{\mathbb{U}} \sum_{w \in W_{0}} \frac{m_{\lambda}(w u) u^{-w \mu} \overline{c(w u)}}{|c(w u)|^{2}} d u \\
& =\frac{W_{0}\left(q^{-1}\right)}{W_{0 \mu}\left(q^{-1}\right)} r^{\mu} \int_{\mathbb{U}} m_{\lambda}(u) \frac{u^{-\mu}}{c(u)} d u
\end{aligned}
$$

and the first claim follows. The final claim follows from Proposition 3.3(iv).
Let $\Upsilon=\left\{u \in \operatorname{Hom}\left(P, \mathbb{C}^{\times}\right):\left|u^{w \lambda}\right| \leqslant r^{\lambda}\right.$ for all $\lambda \in P^{+}$and all $\left.w \in W_{0}\right\}$.
Theorem A.3. - The algebra homomorphism $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ is bounded if and only if $u \in \Upsilon$.

Proof. - In the exceptional case this follows from [10, Theorem 4.7.1], as remarked at the beginning of this appendix. Suppose we are in the standard case. If $u \in \Upsilon$ is nonsingular, then by (1.6) we have

$$
\left|h_{u}\left(A_{\lambda}\right)\right| \leqslant \frac{1}{W_{0}\left(q^{-1}\right)} \sum_{w \in W_{0}}|c(w u)| \quad \text { for all } \lambda \in P^{+}
$$

Thus $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ is bounded, and so by $\left[6\right.$, Theorem I.2.5], $\left|h_{u}\left(A_{\lambda}\right)\right| \leqslant 1$ for all $\lambda \in P^{+}$.

If $u \in \Upsilon$ is singular, it is clear that there exists a sequence $\left(u_{(k)}\right)_{k \in \mathbb{N}}$ in $\Upsilon$ such that each $u_{(k)}$ is nonsingular and $u_{(k)} \rightarrow u$. By the above we have $\left|h_{u_{(k)}}\left(A_{\lambda}\right)\right| \leqslant 1$ for all $\lambda \in P^{+}$, and since each $h_{u_{(k)}}$ is a Laurent polynomial (in the variables $\left\{u_{(k)}^{\lambda_{i}}\right\}_{i \in I_{0}}$ ) it follows that $\left|h_{u}\left(A_{\lambda}\right)\right| \leqslant 1$ for all $\lambda \in P^{+}$. Thus $h_{u}$ is bounded for all $u \in \Upsilon$.

Suppose now that $h_{u}: \mathscr{A} \rightarrow \mathbb{C}$ is bounded (so $\left|h_{u}\left(A_{\lambda}\right)\right| \leqslant 1$ for all $\lambda \in P^{+}$). Then for all $\lambda \in P^{+}$, by (A.2) and Lemma A.2,

$$
\left|m_{\lambda}(u)\right| \leqslant \sum_{\mu \preceq \lambda}\left|b_{\lambda, \mu}\right| \leqslant C\left|\Pi_{\lambda}\right| r^{\lambda}
$$

Thus fixing $\lambda$ and considering $m_{k \lambda}(u)$ for $k \in \mathbb{N}$ gives

$$
\left|\sum_{\mu \in W_{0} \lambda}\left(r^{-\lambda} u^{\mu}\right)^{k}\right| \leqslant p(k)
$$

where $p(k)$ is a polynomial. It is elementary that this implies that $\left|r^{-\lambda} u^{\mu}\right| \leqslant$ 1 for all $\mu \in W_{0} \lambda$, hence the result.

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