# Fourier decay, Renewal theorem and Spectral gaps for random walks on split semisimple Lie groups

# Jialun LI

## Abstract

We establish an exponential error term for the renewal theorem in the context of products of random matrices, which is surprising compared with classical abelian cases. A key tool is the Fourier decay of the Furstenberg measures on the projective spaces, which is a higher dimensional generalization of a recent work of Bourgain-Dyatlov.

**Résumé**: On établit un terme d'erreur exponentiel dans le théorème de renouvellement dans le cadre de produits de matrices aléatoires, qui est inattendu par rapport au cas classique abélien. L'outil clef est le décroissance de Fourier de mesures de Furstenberg sur les espaces projectifs, qui est une généralisation en dimension supérieure d'un travail récent de Bourgain-Dyatlov.

# 1 Introduction

Let V be a finite dimensional irreducible representation of a split semisimple Lie group G (For example  $G = \operatorname{SL}_{m+1}(\mathbb{R})$  and  $V = \mathbb{R}^{m+1}$ ). Let  $X = \mathbb{P}V$  be the real projective space of V, which is the set of lines of V. Then we have a group action of G on X. Let  $\mu$  be a Borel probability measure on G and let  $\Gamma_{\mu}$  be the subgroup generated by the support of  $\mu$ . We call  $\mu$  Zariski dense if  $\Gamma_{\mu}$  is a Zariski dense subgroup of G. This means that the measure  $\mu$  does not concentrate on any proper algebraic subgroup of G. We also need the hypothesis of finite exponential moment. If G is a subgroup of a matrix group, the definition of exponential moment is that there exists  $\epsilon$  positive such that

$$\int_G \|g\|^{\epsilon} \mathrm{d}\mu(g) < \infty.$$

For the general case, please see Definition 2.49. From now on, we always suppose that the measure  $\mu$  is Zariski dense with a finite exponential moment.

We can give a random walk on  $X = \mathbb{P}V$  induced by  $\mu$ . Fix a point x in X. At each step, we go to a random point gx, where g is a random element in G with law  $\mu$ . By a theorem of Furstenberg, this random walk has a unique stationary measure  $\nu$  on X, called the Furstenberg measure or the  $\mu$ -stationary measure. That is to say, the measure  $\nu$  satisfies  $\nu = \mu * \nu := \int_G g_* \nu d\mu(g)$ , where  $g_* \nu$  is the pushforward of  $\nu$  by the action of g on X. This measure was introduced by Furstenberg when he established the law of large numbers for products of random matrices. The properties of the  $\mu$ -stationary measure are also important in other limit theorems for products of random matrices.

For any natural number n, let  $\mu^{*n}$  be the *n*-times convolution of the measure  $\mu$ . If  $X_1, \dots, X_n$  are i.i.d. random variables in G with the same distribution  $\mu$ , then  $\mu^{*n}$  is the distribution of the product  $X_1X_2\cdots X_n$ . We are interested in the limit distribution of the product.

# **Renewal theorem**

Let  $\|\cdot\|$  be a norm on V. For g in G, let  $\|g\|$  be its operator norm. For a positive bounded Borel function f on  $\mathbb{R}$  and a real number t, we define the renewal sum for norm by

$$R_P f(t) := \sum_{n=0}^{+\infty} \int_G f(\log \|g\| - t) \mathrm{d}\mu^{*n}(g).$$

Because of the positivity of f, this sum is well defined. It is natural to try to relate the limit law for norms to the limit law for cocycles. We define the cocycle function  $\sigma : G \times X \to \mathbb{R}$  by, for  $x = \mathbb{R}v$  in X and g in G,  $\sigma(g, x) = \log \frac{||gv||}{||v||}$ . For a positive bounded Borel function f on  $\mathbb{R}$ , the renewal sum for cocycles is defined by

$$Rf(x,t) := \sum_{n=0}^{+\infty} \int_G f(\sigma(g,x) - t) d\mu^{*n}(g), \text{ for } x \in X \text{ and } t \in \mathbb{R}$$

The renewal theorem was first introduced by Blackwell and in our situation by Kesten [Kes74]. The main result (due to Guivarc'h and Le Page [GLP16]) is that when time t tends to infinite, the renewal sum Rf(x,t) tends to  $\frac{1}{\sigma_{\mu}} \int f$ , where  $\sigma_{\mu}$  is the Lyapunov constant defined by  $\sigma_{\mu} := \int_{G \times X} \sigma(g, x) d\mu(g) d\nu(x)$ . From the definition, we see that the Lyapunov constant  $\sigma_{\mu}$  is an average of the cocycle function  $\sigma(g, x)$  with respect to the measure  $\mu \otimes \nu$ . The renewal theorem gives us a phenomenon of equidistribution when the time t is large enough.

**Theorem 1.1.** Let **G** be a connected algebraic semisimple Lie group defined and split over  $\mathbb{R}^1$ and let  $G = \mathbf{G}(\mathbb{R})$  be its group of real points. Let  $\mu$  be a Zariski dense Borel probability measure on G with a finite exponential moment. Let V be an irreducible representation of G with a norm. There exists  $\epsilon > 0$  such that for  $f \in C_c^{\infty}(\mathbb{R})$  and  $t \in \mathbb{R}$ , we have

$$Rf(x,t) = \frac{1}{\sigma_{\mu}} \int_{-t}^{\infty} f(u) dLeb(u) + O_f(e^{-\epsilon|t|}),$$

and if the norm is good

$$R_P f(t) = \frac{1}{\sigma_{\mu}} \int_{-t}^{\infty} f(u) \mathrm{d}Leb(u) + O_f(e^{-\epsilon|t|}),$$

where  $O_f$  depends on the support and some Sobolev norm of f.

**Remark.** For limit law of norms we need an additional hypothesis that the norm is good. For example, when  $G = SL_{m+1}(\mathbb{R})$  and  $V = \mathbb{R}^{m+1}$ , any euclidean norm on  $\mathbb{R}^{m+1}$  is a good norm. For the definition, please see Definition 2.8.

We should compare this result with the renewal theorem on  $\mathbb{R}$  (the commutative case). If  $\mu$  is a measure on  $\mathbb{R}$  whose support is finite, then the error term in the renewal theorem is never exponential.

Our result improves a result of Boyer [Boy16], where the error term is polynomial on t. We hope this type of result can give some exponential error terms in the orbital counting problem of higher rank. Given a discrete subgroup  $\Gamma$  of  $\mathrm{SL}_{m+1}(\mathbb{R})$ , we are interested in the asymptotic growth of  $\#\{\gamma \in \Gamma | d(\gamma o, o) \leq R\}$ , where o is the base point in  $\mathrm{SL}_{m+1}(\mathbb{R})/\mathrm{SO}(m+1)$ . See for instance Lalley [Lal89], Quint [Qui05] and Sambarino [Sam15]. This type of error term is always connected with some spectral gap property. Before going to more technical results, we give a simple example of Fourier decay, which is the key tool used in the proof of spectral gap.

<sup>&</sup>lt;sup>1</sup>For example,  $\mathbf{G} = \mathbf{SL}_{m+1}, \ m \ge 1.$ 

# Fourier decay

Before stating our main result of Fourier decay, we introduce another property of the stationary measure. Guivarc'h established the Hölder regularity of stationary measures, which means that there exist C, c positive such that for every r positive, the r neighbourhood of any hyperplane in X has  $\nu$  measure less than  $Cr^c$ . This implies that the stationary measure  $\nu$  has positive dimension. This also says that  $\nu$  does not concentrate on some hyperplane, which is reasonable due to the hypothesis of Zariski density of  $\mu$ .

Let us see the example  $G = \mathrm{SL}_2(\mathbb{R})$  and  $X = \mathbb{P}(\mathbb{R}^2)$ . Fix the identification of  $\mathbb{P}(\mathbb{R}^2)$  with the circle  $\mathbb{T} \simeq \mathbb{R}/\pi\mathbb{Z}$ , given by the transitive action of the group PSO<sub>2</sub>. We can define the Fourier coefficients of the stationary measure  $\nu$  by

$$\hat{\nu}(k) = \int_{\mathbb{T}} e^{2ikx} \mathrm{d}\nu(x), \text{ for } k \in \mathbb{Z}.$$

**Theorem 1.2.** Let  $\mu$  be a Zariski dense Borel probability measure on  $SL_2(\mathbb{R})$  with a finite exponential moment. Let  $\nu$  be the  $\mu$ -stationary measure on  $\mathbb{T}$ . Then there exists  $\epsilon$  positive such that

$$|\hat{\nu}(k)| = O(|k|^{-\epsilon}).$$
 (1.1)

In other words, the Fourier coefficients of the stationary measure have polynomial decay. By a general argument, the polynomial decay of Fourier coefficients implies Guivarc'h's regularity. But the regularity is also a crucial ingredient in the proof. This generalizes the recent work of Bourgain-Dyatlov [BD17]. For more similar results, please see [LI18a] and its references.

A similar Fourier decay for the Lie group  $SL_2(\mathbb{C})$  is established in [LNP19] for Patterson-Sullivan measures, which cannot be treated by our method due to the non splitness of  $SL_2(\mathbb{C})$ . It would also be interesting to establish a similar Fourier decay for the group  $SL_2(\mathbb{Q}_p)$  and the stationary measure on  $\mathbb{P}^1_{\mathbb{Q}_p}$ .

# Spectral gap

Equip  $\mathbb{P}V$  with a Riemannian distance. For  $\gamma$  positive, let  $C^{\gamma}(\mathbb{P}V)$  be the space of  $\gamma$ -Hölder functions. We introduce the transfer operator, which is an analogue of the characteristic function in our case.

**Definition.** For z in  $\mathbb{C}$  with the real part  $|\Re z|$  small enough, let  $P_z$  be the operator on the space of continuous functions, which is given by

$$P_z f(x) = \int_G e^{z\sigma(g,x)} f(gx) d\mu(g), \text{ for } x \in \mathbb{P}V.$$

We keep the assumption that  $\mu$  is a Zariski dense Borel probability measure on G with a finite exponential moment. The use of this transfer operator on the products of random matrices has been introduced by Guivarc'h and Le Page. Due to the property of exponential moment, when  $|\Re z|$  is small enough, the operator  $P_z$  preserves the Banach space  $C^{\gamma}(\mathbb{P}V)$  for  $\gamma > 0$  small enough. Due to the contracting action of G on X, for z in a small ball centred at 0, the spectral radius of  $P_z$  on  $C^{\gamma}(\mathbb{P}V)$  is less than 1 except at 0. Due to the non-arithmeticity of  $\Gamma_{\mu}$ , on the imaginary line, the operator  $P_z$  also has spectral radius less than 1 except at 0. These were used to give limit theorems for products of random matrices by Le Page and Guivarc'h (Please see [LP82] and [BQ16]).

**Theorem 1.3.** Let **G** be a connected algebraic semisimple Lie group defined and split over  $\mathbb{R}$ and let  $G = \mathbf{G}(\mathbb{R})$  be its group of real points. Let  $\mu$  be a Zariski dense Borel probability measure on G with finite exponential moment. Let V be an irreducible representation of G with a norm. For every  $\gamma > 0$  small enough, there exists  $\delta > 0$  such that for all |b| > 1 and |a| small enough the spectral radius of  $P_{a+ib}$  acting on  $C^{\gamma}(\mathbb{P}V)$  satisfies

$$\rho(P_{a+ib}) < 1 - \delta.$$

Even in the case of  $SL_2(\mathbb{R})$ , the result is new and only known in some special case. When  $\mu$  is supported on a finite number of elements of  $SL_2(\mathbb{R})$  and these elements generate a Schottky semi group, this result is due to Naud [Nau05]. When  $\mu$  is absolutely continuous with respect to the Haar measure on  $SL_2(\mathbb{R})$ , this result can be obtained directly using high oscillations.

It is interesting that the three objects, the Fourier decay, the Renewal theorem and the spectral gap are roughly equivalent. In [L118a], we use the Renewal theorem to prove the Fourier decay. In this manuscript, we use the Fourier decay to prove the spectral gap, and then use the spectral gap to prove the Renewal theorem. They are analogue with similar objects for convex cocompact surfaces. In this more geometric setting, the Fourier decay was recently studied by Bourgain-Dyatlov; the spectral gap can be interpreted as the zero free region of the Selberg zeta function or the gap of the eigenvalues of the Laplace operator on the surface; the renewal theorem is replaced by the counting problem of the lattice points or the primitive closed geodesics. Please see Borthwick [Bor07] and the references there.

This result should be compared with similar results for random walks on  $\mathbb{R}$ . Let  $\mu$  be a Borel probability measure on  $\mathbb{R}$  with finite support. Then

$$\liminf_{|b| \to \infty} |1 - \hat{\mu}(ib)| = 0,$$

which is totally different from our case and where  $\hat{\mu}(z)$  is the Laplace transform of the measure  $\mu$ , given by

$$\hat{\mu}(z) = \int_{\mathbb{R}} e^{zx} \mathrm{d}\mu(x).$$

The proof is direct. Let  $\{x_1, \ldots, x_l\}$  be the support of  $\mu$ . Then  $\hat{\mu}(ib) = \sum_{1 \leq j \leq l} \mu(x_j) e^{ibx_j}$ , and we only need to find b such that all the terms are uniformly near 1. Using the fact that  $\liminf_{b\to\infty} d_{\mathbb{R}^l}(b(x_1, \ldots, x_l), 2\pi\mathbb{Z}^l) = 0$ , we have the claim.

An analogous result is valid if we replace the projective space  $\mathbb{P}V$  by the flag variety  $\mathscr{P}$ . Let  $\mathscr{P}$  be the full flag variety of G and let  $\mathfrak{a}$  be a Cartan subspace of the Lie algebra  $\mathfrak{g}$  of G. For  $g \in G$  and  $\eta \in \mathscr{P}$ , let  $\sigma(g, \eta)$  be the Iwasawa cocycle, which takes values in  $\mathfrak{a}$ . We fix a Riemannian distance on  $\mathscr{P}$ . We can similarly define the space of  $\gamma$ -Hölder functions  $C^{\gamma}(\mathscr{P})$ . Let  $\varpi, \vartheta$  be in  $\mathfrak{a}^*$ . For a continuous function f on  $\mathscr{P}$  and  $|\varpi|$  small enough, the transfer operator  $P_{\varpi+i\vartheta}$  on the flag variety is defined by

$$P_{\varpi+i\vartheta}f(\eta) = \int_G e^{(\varpi+i\vartheta)\sigma(g,\eta)} f(g\eta) \mathrm{d}\mu(g).$$

**Theorem 1.4** (Spectral gap). Let **G** be a connected algebraic semisimple Lie group defined and split over  $\mathbb{R}$  and let  $G = \mathbf{G}(\mathbb{R})$  be its group of real points. Let  $\mu$  be a Zariski dense Borel probability measure on G with finite exponential moment. For every  $\gamma > 0$  small enough, there exists  $\delta > 0$  such that for all  $\vartheta, \varpi$  in  $\mathfrak{a}^*$  with  $|\vartheta| > 1$  and  $|\varpi|$  small enough the spectral radius of  $P_{\varpi+i\vartheta}$  acting on  $C^{\gamma}(\mathscr{P})$  satisfies

$$\rho(P_{\varpi+i\vartheta}) < 1 - \delta.$$

# Main technical result

The key ingredient of the proof of the above results is the following Fourier decay property of the  $\mu$ -stationary measure on the flag variety  $\mathscr{P}$ . We start with the case for  $SL_2(\mathbb{R})$ .

**Theorem 1.5.** Let  $\mu$  be a Zariski dense Borel probability measure on  $SL_2(\mathbb{R})$  with a finite exponential moment. Let  $X = \mathbb{P}(\mathbb{R}^2)$  and let  $\nu$  be the  $\mu$ -stationary measure on X.

For every  $\gamma > 0$ , there exist  $\epsilon_0 >, \epsilon_1 > 0$  depending on  $\mu$  such that the following holds. For any pair of real functions  $f \in C^2(X)$ ,  $r \in C^{\gamma}(X)$  and  $\xi > 0$  such that  $|\varphi'| \ge \xi^{-\epsilon_0}$  on the support of r,  $||r||_{\infty} \le 1$  and

$$\|\varphi\|_{C^2} + c_\gamma(r) \le \xi^{\epsilon_0}$$

then

$$\left|\int e^{i\xi\varphi(x)}r(x)\mathrm{d}\nu(x)\right|\leq\xi^{-\epsilon_1}\quad\text{ for all }\xi\text{ large enough}.$$

**Remark 1.6.** Theorem 1.2 is a corollary Theorem 1.5. This is also a generalization of the same theorem for the Patterson-Sullivan measures as in [BD17].

Theorem 1.5 is a particular case of a more general result: Theorem 1.7 below. In order to state the Fourier decay on the flag variety, we need to introduce a special condition. Let r be a continuous function on  $\mathscr{P}$  and let C > 0. For a  $C^2$  function  $\varphi$  on  $\mathscr{P}$ , we say  $\varphi$  is (C, r) good if it satisfies some assumptions on the Lipschitz norm and derivative, which will be defined later (Definition 4.1). Recall that for a  $\gamma$ -Hölder function f, we have defined  $c_{\gamma}(f) =$  $\sup_{x \neq x'} \frac{|f(x) - f(x')|}{d(x,x')^{\gamma}}$ . Due to some technical problem, we will only prove a simply connected case in Section 4. (For example the group  $\mathbf{SL}_{m+1}$  is simply connected but  $\mathbf{PGL}_{m+1}$  is not.) The general case will be proved in Appendix 5.1 by a covering argument.

**Theorem 1.7** (Fourier decay). Let  $\mathbf{G}$  be a connected  $\mathbb{R}$ -split reductive group whose semisimple part is simply connected and let  $G = \mathbf{G}(\mathbb{R})$  be its group of real points. Let  $\mu$  be Zariski dense Borel probability measure on G with finite exponential moment. Let  $\nu$  be the  $\mu$ -stationary measure on the flag variety  $\mathscr{P}$ .

For every  $\gamma > 0$ , there exist  $\epsilon_0 > 0, \epsilon_1 > 0$  depending on  $\mu$  such that the following holds. For any pair of real functions  $\varphi \in C^2(\mathscr{P}), r \in C^{\gamma}(\mathscr{P})$  and  $\xi > 0$  such that  $\varphi$  is  $(\xi^{\epsilon_0}, r)$  good,  $\|r\|_{\infty} \leq 1$  and  $c_{\gamma}(r) \leq \xi^{\epsilon_0}$ , then

$$\left| \int e^{i\xi\varphi(\eta)} r(\eta) \mathrm{d}\nu(\eta) \right| \le \xi^{-\epsilon_1} \quad \text{for all } \xi \text{ large enough.}$$
(1.2)

**Remark 1.8.** The decay rate only depends on the constants in the large deviation principles and the regularity of stationary measures. This should be compared with [BD17], where the spectral gap and the decay rate only depend on the dimension of the Patterson-Sullivan measure.

When  $G = \mathrm{SL}_2(\mathbb{R})$ , the (C, r) goodness is exactly the assumption of  $\varphi$  in Theorem 1.5, which is natural for having a Fourier decay. Theorem 1.7 clearly implies Theorem 1.5.

The proof of Theorem 1.7 follows the similar strategy as in [BD17]. But in this higher dimension and higher rank case, new difficulties appear and we need new ideas to overcome these difficulties.

The main difficulty comes from higher dimension, that is the verification of a subtle nonconcentration hypothesis, which is also the main difficulty that prevents Bourgain-Dyatlov from generalizing their result to higher dimension [BD17, Page 4]. For applying the higher dimensional discretized sum-product estimate, Proposition 3.17 (a generalized version of a result of Bourgain [Bou10]), we need to verify certain measures on  $\mathbb{R}^m$  are not concentrated on any affine subspaces of  $\mathbb{R}^m$ . The key idea is to transfer the problem to an estimate of volume, which gives nonconcentration for all affine subspaces simultaneously (Corollary 3.6). Then we use representation theory and Guivarc'h regularity to verify the condition in Section 3.

Another difficulty comes from the higher rank. In rank one case, the action of group G on the flag variety  $\mathscr{P}$  is conformal. But in higher rank case, this is not true. Using root systems, we are able to find the directions of slowest contraction speed on the tangent bundle of the flag variety  $\mathscr{P}$  under the action of G, explained in Section 2.

# Notation

We will make use of some classical notation: for two real functions f and g, we write  $f = O(g), f \ll g$  or  $g \gg f$  if there exists a constant C > 0 such that  $|f| \leq Cg$ , where C only depends on the ambient group G and the measure  $\mu$ . We write  $f \asymp g$  if  $f \ll g$  and  $g \ll f$ . We write  $f = O_{\epsilon}(g), f \ll_{\epsilon} g$  or  $g \gg_{\epsilon} f$  if the constant C depends on an extra parameter  $\epsilon > 0$ .

We always use  $0 < \delta < 1$  to denote an error term and  $0 < \beta < 1$  to denote the magnitude. The quantity  $\beta^{-1}$  is supposed to be greater than  $\delta^{-O(1)}$ . If  $\delta^{O(1)}f \leq g \leq \delta^{-O(1)}f$ , then we say that f and g are of the same size.

# 2 Random walks on Reductive groups

The representation theory of algebraic groups is more clear than the representation theory of Lie groups. We will use the vocabulary of algebraic groups. In this manuscript, without further assumption, we assume **G** is a connected  $\mathbb{R}$ -split reductive  $\mathbb{R}$ -group. From Section 2.3 to Section 4, we add the assumption that the semisimple part is simply connected. Please see [Hel79], [Bor90] and [BQ16] for more details.

We write **G** for an algebraic group, and  $G = \mathbf{G}(\mathbb{R})$  for its group of real points, equipped with the Lie group topology.

# 2.1 Reductive groups and representations

# **Reductive groups**

Let **G** be a connected  $\mathbb{R}$ -split reductive  $\mathbb{R}$ -group. Let **A** be a maximal  $\mathbb{R}$ -split torus in **G**. Because **G** is  $\mathbb{R}$ -split, the group **A** is also the maximal torus of **G** and the centralizer of **A** in **G** is **A**. Let **C** be the connected component of the centre of **G**, which is contained in the maximal torus **A**. The semisimple part of **G** is the derived group  $\mathscr{D}\mathbf{G} = [\mathbf{G}, \mathbf{G}]$ . Let **B** be the subtorus of **A** given by  $\mathbf{A} \cap \mathscr{D}\mathbf{G}$ . The dimension of **A** and **B** are called the reductive rank and the semisimple rank of **G**, respectively. We write r and m for the reductive rand and the semisimple rank.

Because we are dealing with real groups, we will use transcendental methods to describe the structure of **G**. Let G, A, B and C be the group of real points of **G**, **A**, **B** and **C**. Let  $\theta$  be a Cartan involution of G which satisfies  $\theta(A) = A$  and such that the set of fixed points  $K = \{g \in G | \ \theta(g) = g\}$  is a maximal compact subgroup of G. Let  $\mathfrak{g}, \mathfrak{k}, \mathfrak{a}, \mathfrak{b}$  and  $\mathfrak{c}$  be the Lie algebra of G, K, A, B and C, respectively. Then  $\mathfrak{a} = \mathfrak{b} \oplus \mathfrak{c}$  due to  $\mathfrak{g} = \mathscr{D}\mathfrak{g} \oplus \mathfrak{c}$ . We write exp for the exponential map from  $\mathfrak{a}$  to A. We also write  $\theta$  for the differential of the Cartan involution, whose set of fixed points is  $\mathfrak{k}$  and which equals -id on  $\mathfrak{a}$ .

For X, Y in  $\mathfrak{g}$ , the Killing form is defined as

$$K(X,Y) = tr(\mathrm{ad}X\mathrm{ad}Y).$$

The Killing form is positive definite on  $\mathfrak{b}$  and negative definite on  $\mathfrak{k}$ . Endowed with the Killing form, the Lie algebra  $\mathfrak{b}$  and its dual  $\mathfrak{b}^*$  become Euclidean spaces.

## Root systems and the Weyl group

The spaces  $\mathfrak{b}^*$  and  $\mathfrak{c}^*$  are seen as subspaces of  $\mathfrak{a}^*$ , which takes value zero on  $\mathfrak{c}$  and  $\mathfrak{b}$ , respectively. Let R be the root system of  $\mathfrak{g}$  with respect to  $\mathfrak{a}$ , that is the set of nontrivial weights of the adjoint action of  $\mathfrak{a}$  on  $\mathfrak{g}$ . It is actually a subset of  $\mathfrak{b}^*$ . Because  $\mathfrak{c}$  is in the centre of  $\mathfrak{g}$ , its adjoint action on  $\mathfrak{g}$  is trivial. Fix a choice of positive roots  $R^+$ . Let  $\Pi$  be the collection of primitive simple roots of  $R^+$ . Let  $\mathfrak{a}^+$  be the Weyl chamber defined by  $\{X \in \mathfrak{a} | \alpha(X) \ge 0, \forall \alpha \in \Pi\}$ . Let  $\mathfrak{a}^{++}$  be the interior of Weyl chamber defined by  $\{X \in \mathfrak{a} | \alpha(X) \ge 0, \forall \alpha \in \Pi\}$ . Using the root system, we have a decomposition of  $\mathfrak{g}$  into eigenspaces of  $\mathfrak{a}$ ,

$$\mathfrak{g} = \mathfrak{z} \oplus \bigoplus_{\alpha \in R} \mathfrak{g}^{lpha},$$

where  $\mathfrak{z}$  is the centralizer of  $\mathfrak{a}$  and  $\mathfrak{g}^{\alpha}$  is the eigenspace given by

$$\mathfrak{g}^{\alpha} = \{ X \in \mathfrak{g} | \ [Y, X] = \alpha(Y) X \text{ for all } Y \in \mathfrak{a} \}.$$

Since the group **G** is split, we know  $\mathfrak{a} = \mathfrak{z}$  and that  $\mathfrak{g}^{\alpha}$  are of dimension 1.

Recall that for every root  $\alpha$  in R, there is an orthogonal symmetry  $s_{\alpha}$  which preserves Rand  $s_{\alpha}(\alpha) = -\alpha$ . For  $\alpha \in R$ , let  $H_{\alpha}$  be the unique element in  $\mathfrak{b}$  such that  $s_{\alpha}(\alpha') = \alpha' - \alpha'(H_{\alpha})\alpha$ for  $\alpha' \in \mathfrak{b}^*$ . The set  $\{H_{\alpha} | \alpha \in R\}$  is called the set of duals roots in  $\mathfrak{b}$ . Since the Cartan involution  $\theta$  equals -id on  $\mathfrak{a}$ , this implies  $\theta \mathfrak{g}^{\alpha} = \mathfrak{g}^{-\alpha}$  for  $\alpha \in R$ . Using the Killing form, we can prove that  $[\mathfrak{g}^{\alpha}, \mathfrak{g}^{-\alpha}] = \mathbb{R}H_{\alpha}$ . (See [Ser66, Cha. 4, Theorem 2] for more details) Hence, there is a unique choose (up to sign)  $X_{\alpha} \in \mathfrak{g}^{\alpha}, Y_{\alpha} \in \mathfrak{g}^{-\alpha}$  such that

$$[X_{\alpha}, Y_{\alpha}] = H_{\alpha} \text{ and } \theta(X_{\alpha}) = -Y_{\alpha}.$$

Let  $K_{\alpha} = X_{\alpha} - Y_{\alpha}$ . Due to  $\theta K_{\alpha} = K_{\alpha}$ , the element  $K_{\alpha}$  is in  $\mathfrak{k}$ .

Let W be the Weyl group of R. Then the group W acts simply transitively on the set of Weyl chambers. Let  $w_0$  be the unique element in W which sends the Weyl chamber  $\mathfrak{a}^+$  to the Weyl chamber  $-\mathfrak{a}^+$ . Let  $\iota = -w_0$  be the opposition involution. The Weyl group also acts on  $\mathfrak{a}^*$ by the dual action. Let  $N_G(A)$  be the normalizer of A in G. An element in  $N_G(A)/A$  induces an automorphism on the tangent space  $\mathfrak{a}$ . This gives an isomorphism from  $N_G(A)/A$  to the Weyl group W. Hence  $w_0$  can be realized as an element in G/A and its action on  $\mathfrak{a}$  is given by conjugation.

# The Iwasawa cocycle

Let  $\mathbf{n} = \bigoplus_{\alpha \in R^+} \mathbf{g}^{\alpha}$  and  $\mathbf{n}^- = \bigoplus_{\alpha \in R^+} \mathbf{g}^{-\alpha}$ . They are nilpotent Lie algebras. Let **N** be the connected algebraic subgroup of **G** with Lie algebra  $\mathbf{n}$ . The group **N** is normalized by **A**. Let  $\mathbf{P} = \mathbf{A} \ltimes \mathbf{N}$  be a minimal parabolic subgroup. The flag variety  $\mathscr{P}$  is defined to be the set of conjugations of P under the action of G. Since the normalizer of P in G is itself, we have an isomorphism

$$G/P \to \mathscr{P}$$

We write  $\eta_o$  for the subgroup P seen as a point in  $\mathscr{P}$ . Let M be the subgroup of A, whose element has order at most two. Since A is isomorphic to  $(\mathbb{R}^*)^r$ , we know that  $M \simeq (\mathbb{Z}/2\mathbb{Z})^r$  and  $A = M \times A_e$ , where  $A_e = \exp(\mathfrak{a})$  is the analytical connected component of A and  $A_e \simeq (\mathbb{R}_{>0})^r$ . Let  $G^o$  be the connected component of the identity element in G.

# **Lemma 2.1.** Let **G** be a connected $\mathbb{R}$ -split reductive $\mathbb{R}$ -group. Then we have $K = K^{\circ}M$ .

*Proof.* By Matsumoto's theorem [Mat64], we have  $G = G^{\circ}A = G^{\circ}M$ . Hence the group M intersect each connected component of G. We claim that  $K = K^{\circ}M$ .

We know that  $K \supset K^o M$ , because the group M equals to  $A \cap K$  due to [Ben05, Lemme 4.2]. This can be proved directly by considering the action of the Cartan involution on M. The group  $K^o M$  intersects each connected component of G and the intersection with  $G^o$  contains  $K^o$ . By maximality of K, we conclude that  $K = K^o M$ .

We have an Iwasawa decomposition of G given by

$$G = KAN.$$

The action of K on  $\mathscr{P}$  is transitive. Hence  $\mathscr{P}$  is a compact manifold. By Lemma 2.1, we have  $G = KA_eMN = KA_eN$ . This is a bijection between G and  $K \times A_e \times N$ . Then we can define the Iwasawa cocycle  $\sigma$  from  $G \times \mathscr{P}$  to  $\mathfrak{a}$ . Let  $\eta$  be in  $\mathscr{P}$  and g be in G. By the transitivity of K, there exists  $k \in K$  such that  $\eta = k\eta_o$ . By the Iwasawa decomposition, there exists a unique element  $\sigma(g, \eta)$  in  $\mathfrak{a}$  such that

$$gk \in K \exp(\sigma(g,\eta))N.$$

We can verify that this is well defined and  $\sigma$  is an additive cocycle, that is for g, h in G and  $\eta$  in  $\mathscr{P}$ 

$$\sigma(gh,\eta) = \sigma(g,h\eta) + \sigma(h,\eta).$$

Due to the direct sum  $\mathfrak{a} = \mathfrak{b} \oplus \mathfrak{c}$ , we can decompose the Iwasawa decomposition into the semisimple part and the central part of the cocycle, that is

$$\sigma(g,\eta) = \sigma_{ss}(g,\eta) + c(g),$$

where  $\sigma_{ss}$  lies in  $\mathfrak{b}$  and c(g) in  $\mathfrak{c}$ . The central part c(g) does not dependent on  $\eta$ , because the map

$$G \to G/\mathscr{D}G$$

kills the semisimple part and the restriction on  $C_e = \exp(\mathfrak{c})$  is injective. Moreover, since the Iwasawa cocycle is additive, the central part is additive. That is for g, h in G

$$c(gh) = c(g) + c(h).$$

#### The Cartan decomposition

The Cartan decomposition says that  $G = KA^+MK = KA^+K$ , where  $A^+$  is the image of the Weyl chamber  $\mathfrak{a}^+$  under the exponential map. For g in G, by Cartan decomposition, we can write  $g = k_g a_g \ell_g$  with  $k_g, \ell_g$  in K and  $a_g$  in  $A^+$ . The element  $a_g$  is unique and there is a unique element  $\kappa(g)$  in  $\mathfrak{a}^+$  such that  $a_g = \exp(\kappa(g))$ . We call  $\kappa(g)$  the Cartan projection of g. Then  $\kappa(g^{-1}) = \iota \kappa(g)$ , where  $\iota$  is the opposition involution. Since A is contained in P, we can define  $\zeta_o = w_0 \eta_o$ , where the element  $w_0$  in the Weyl group is seen as an element in G/A. (As an element in  $\mathscr{P}, \zeta_o$  is the opposite parabolic group with respect to P and A) Let  $\eta_g^M = k_g \eta_o$  and  $\zeta_g^m = \ell_g^{-1} \zeta_o$ . When  $\kappa(g)$  is in  $\mathfrak{a}^{++}$ , they are uniquely defined, independently of the choice of  $k_g$ and  $\ell_g$ .

We can also define a unique decomposition of  $\kappa(g)$  into semisimple part and central part. Due to  $\kappa(g) = \sigma(g, \ell_q^{-1} \eta_o) = \sigma_{ss}(g, \ell_q^{-1} \eta_o) + c(g)$ , we have

$$\kappa(g) = \kappa_{ss}(g) + c(g).$$

# **Dominant** weights

Here, we need the hypothesis that  $\mathscr{D}\mathbf{G}$  is simply connected.

Let  $X(\mathbf{A})$  and  $X(\mathbf{B})$  be the character groups of  $\mathbf{A}$  and  $\mathbf{B}$ , respectively. We will identify  $X(\mathbf{A})$  and  $X(\mathbf{B})$  as discrete subgroups of  $\mathfrak{a}^*$  and  $\mathfrak{b}^*$  by taking differential. The elements of  $\mathfrak{a}^*$  in  $X(\mathbf{A})$  are called weights. All the roots are weights, because they come from adjoint action of A on  $\mathfrak{g}^{\alpha}$ .

Since  $\{H_{\alpha}\}_{\alpha\in\Pi}$  is a basis of  $\mathfrak{b}$ , let  $\{\tilde{\omega}_{\alpha}\}_{\alpha\in\Pi}$  be the dual basis, it is called the fundamental weights. Since the derived group  $\mathscr{D}\mathbf{G}$  is simply connected, we have

$$\mathbf{X}(\mathbf{B}) = \bigoplus_{\alpha \in \Pi} \mathbb{Z} \tilde{\omega}_{\alpha}.$$
 (2.1)

Since **B** is a closed subgroup of a split torus **A**, every character on **B** extends to a character on **A**. Hence for  $\alpha \in \Pi$  there exist

$$\chi_{\alpha} \in \mathcal{X}(\mathbf{A}) \text{ such that } \chi_{\alpha}|_{\mathfrak{b}} = \tilde{\omega}_{\alpha}.$$
 (2.2)

We fix this choice of  $\chi_{\alpha}$ . We write  $\omega_{\alpha}$  for the element in  $\mathfrak{a}^*$  which is another extension of  $\tilde{\omega}_{\alpha}$  and vanishes on  $\mathfrak{c}$ , that is

$$\omega_{\alpha}|_{\mathfrak{b}} = \tilde{\omega}_{\alpha} \text{ and } \omega_{\alpha}|_{\mathfrak{c}} = 0.$$
 (2.3)

The element  $\omega_{\alpha}$  is not always a character of **A**, but a multiple of  $\omega_{\alpha}$  will be a character of **A**. Because  $\omega_{\alpha}$  can be expressed as a linear combination of simple roots with rational coefficients.

Recall that a weight is a dominant weight, if for every w in the Weyl group W, the difference  $\chi - w(\chi)$  is a sum of positive roots.

**Lemma 2.2.** If  $\mathscr{D}\mathbf{G}$  is simply connected. For every  $\alpha \in \Pi$ , the weight  $\chi_{\alpha}$  is a dominant weight.

*Proof.* The action of the Weyl group on  $\mathfrak{c}^*$ , the linear functionals vanishing on  $\mathfrak{b}$ , is trivial. Because  $\tilde{\omega}_{\alpha}$  is a fundamental weight, we know that

$$\chi_{\alpha} - w(\chi_{\alpha}) = \omega_{\alpha} - w(\omega_{\alpha})$$

equals a sum of positive roots. The proof is complete.

# **Representations and highest weight**

Let  $(\rho, V)$  be a linear finite dimensional algebraic representation of G. We only consider finite dimensional representations here. The set of restricted weights  $\Sigma(\rho)$  of the representation is the set of elements  $\omega$  in  $\mathfrak{a}^*$  such that the eigenspace

$$V^{\omega} = \{ v \in V | \forall X \in \mathfrak{a}, \ \mathrm{d}\rho(X)v = \omega(X)v \}$$

is nonzero, where  $d\rho$  is the tangent map of  $\rho$  from  $\mathfrak{g}$  to End(V). From definition, we see that  $\omega$  is the differential of a character on  $\mathbf{A}$ , which is a weight. We define a partial order on the restricted weights: For  $\omega_1, \omega_2$  in  $\Sigma(\rho)$ ,

 $\omega_1 \geq \omega_2 \Leftrightarrow \omega_1 - \omega_2$  is a sum of positive roots.

If  $\omega$  is in  $\Sigma(\rho)$ , then we say that  $\omega$  is a weight of V and a vector v in  $V^{\omega}$  is said to have weight  $\omega$ . We call  $\rho$  proximal if there exists  $\chi$  in  $\Sigma(\rho)$  which is greater than the other restricted weights and such that  $V^{\chi}$  is of dimension 1. We should pay attention that a proximal representation is not supposed to be irreducible. The advantage of the splitness of G is that all the irreducible representations are proximal, which will be extensively used later on.

Suppose that  $(\rho, V)$  is an irreducible representation. Let  $\chi \in \mathfrak{a}^*$  be the highest weight of  $(\rho, V)$ . We write  $V_{\chi,\eta} = \rho(g)V^{\chi}$  for  $\eta = g\eta_o$ , which is well defined because the parabolic subgroup P fixes the subspace  $V^{\chi}$ . This gives a map from  $\mathscr{P}$  to  $\mathbb{P}V$  by

$$\mathscr{P} \to \mathbb{P}V, \ \eta \mapsto V_{\chi,\eta}.$$
 (2.4)

In the case of split reductive groups, for a character  $\chi$  on **A**, there exists a irreducible algebraic representation with highest weight  $\chi$  if and only if  $\chi$  is a dominant weight [Tit71]. Let

$$\Theta_{\rho} = \{ \alpha \in \Pi : \ \chi - \alpha \text{ is a weight of } \rho \}.$$

By Lemma 2.2, we have

**Lemma 2.3.** If  $\mathscr{D}\mathbf{G}$  is simply connected. There exists a family of representations  $(\rho_{\alpha}, V_{\alpha})_{\alpha \in \Pi}$ such that the highest weight of  $\rho_{\alpha}$  is  $\chi_{\alpha}$ . Furthermore,  $\Theta_{\rho_{\alpha}} = \{\alpha\}$ . The product of the maps given by (2.4)

$$\mathscr{P} \longrightarrow \prod_{\alpha \in \Pi} \mathbb{P} V_{\alpha}, \quad \eta \mapsto (V_{\chi_{\alpha},\eta})_{\alpha \in \Pi},$$

is an embedding of  $\mathcal{P}$  to the product of projective spaces.

**Lemma 2.4.** Let  $(\rho, V)$  be an irreducible representation of G with highest weight  $\chi$ . Then  $\Theta_{\rho} = \{\alpha\}$  is equivalent to say that  $\chi(H_{\alpha}) > 0$  for only one simple root  $\alpha$ .

*Proof.* Consider the representation of the Lie algebra  $\mathfrak{s}_{\alpha} = \langle H_{\alpha}, X_{\alpha}, Y_{\alpha} \rangle$  on v of highest weight. By the classification of the representations of  $\mathfrak{sl}_2$ , we know that  $Y_{\alpha}v \neq 0$  if and only if  $\chi(H_{\alpha}) > 0$ . The vector  $Y_{\alpha}v$  is the only way to obtain a vector of weight  $\chi - \alpha$  by [Ser66, Chapter 7, Proposition 2]. The proof is complete.

This lemma explains for  $\rho_{\alpha}$ , we have  $\Theta_{\rho_{\alpha}} = \{\alpha\}$ , due to  $\chi_{\alpha}(H_{\beta}) = \tilde{\omega}_{\alpha}(H_{\beta}) = \delta_{\alpha\beta}$ . This family of representation will be fixed from now on until Section 4.3.

For  $\alpha \in \Pi$ , let

$$\tilde{\chi}_{\alpha} = n_{\alpha}\omega_{\alpha}, \text{ where } n_{\alpha} \in \mathbb{N} \text{ and } \tilde{\chi}_{\alpha} \text{ is a dominant weight.}$$
(2.5)

This gives another family of representations  $\tilde{V}_{\alpha}$ , which will be used only in Section 4.3. The main difference with  $\chi_{\alpha}$  is that  $\tilde{\chi}_{\alpha}$  vanishes on  $\mathfrak{c}$ . For semisimple case, the elements  $\omega_{\alpha}, \chi_{\alpha}, \tilde{\chi}_{\alpha}$  are the same.

**Definition 2.5** (Super proximal representation). Let  $(\rho, V)$  be an irreducible representation of G with highest weight  $\chi$ . We call V super proximal if the exterior square  $\wedge^2 V$  is also proximal. This is equivalent to  $\Theta_{\rho} = \{\alpha\}$ , and  $V^{\chi-\alpha}$  is of dimension 1.

**Lemma 2.6.** If the highest weight  $\chi$  of an irreducible representation satisfies  $\chi(H_{\alpha}) > 0$  for only one simple root  $\alpha$ , then this representation is super proximal.

*Proof.* Because the central part of G preserves eigenspaces of A. It is also an irreducible representation of the semisimple part. It will be thus sufficient to prove the semisimple case.

Let  $\alpha$  be the simple root. Let v be a nonzero vector with highest weight  $\chi$ . By [Ser66, Chapter 7, Proposition 2], the representation V is generated by vectors  $Y_{\beta_1} \cdots Y_{\beta_k} v$ , where  $\beta_1, \ldots, \beta_k$  are positive roots. Hence a vector of weight  $\chi - \alpha$  can only be obtained by  $Y_{\alpha}v$ . The dimension of  $V^{\chi-\alpha}$  is no greater than 1. Since  $\chi - \alpha$  is a weight due to Lemma 2.4, the proof is complete.

For  $\chi \in \mathfrak{a}^*$ , if it is a weight, we will use  $\chi^{\sharp}$  to denote its corresponding algebraic character in  $X(\mathbf{A})$ . By the definition of eigenspace  $V^{\chi}$ , we have

**Lemma 2.7.** Let  $(\rho, V)$  be an irreducible representation of G. Let  $\chi^{\sharp}$  be an algebraic character of A. For a in A and  $v \in V^{\chi}$ , we have

$$\rho(a)v = \chi^{\sharp}(a)v.$$

This lemma will be used to determiner the sign in Section 2.5.

#### **Representations and good norms**

**Definition 2.8.** Let  $\|\cdot\|$  be an euclidean norm on a representation  $(\rho, V)$  of G. We call  $\|\cdot\|$  a good norm if  $\rho(A)$  is symmetric and  $\rho(K)$  preserves the norm.

By [Hel79], [BQ16, Lemma 6.33], good norms exist on every representation of G. One advantage of good norm is that for v, u in V and g in G

$$\langle \rho(g)v, u \rangle = \langle v, \rho(\theta(g^{-1})u) \rangle$$

where  $\theta$  is the Cartan involution. The above equation is true because it is true for g in A and K. This means that for good norm we have

$${}^{t}\rho(g) = \rho(\theta(g^{-1})).$$
 (2.6)

The application (2.4) enables us to get information on  $\mathscr{P}$  by the representations. For an element g in GL(V), let ||g|| be its application norm.

**Lemma 2.9.** Let G be a connected reductive  $\mathbb{R}$ -group. Let  $(\rho, V)$  be an irreducible linear representation of G with good norm. Let  $\chi$  be the highest weight of V. For  $\eta$  in  $\mathscr{P}$  and a non zero vector  $v \in V_{\chi,\eta}$ , we have

$$\frac{\|\rho(g)v\|}{\|v\|} = \exp(\chi\sigma(g,\eta)),\tag{2.7}$$

$$\|\rho(g)\| = \exp(\chi\kappa(g)). \tag{2.8}$$

Please see [BQ16, Lemma 8.17] for the proof.

# Examples

For the group  $\mathbf{GL}_{m+1}$ , the maximal torus A can be taken as the diagonal subgroup and the Lie algebra  $\mathfrak{a}$  is the set of diagonal matrices. The Lie algebra  $\mathfrak{b}$  is the subset of  $\mathfrak{a}$  with trace zero. The Lie algebra  $\mathfrak{c} = \{X \in \mathfrak{a} | x_1 = x_2 = \cdots = x_{m+1}\}$ . For X in  $\mathfrak{a}$ , we write  $X = \operatorname{diag}(x_1, \ldots, x_{m+1})$  with  $x_i \in \mathbb{R}$ . Let  $\lambda_i$  in  $\mathfrak{a}^*$  be the linear map given by  $\lambda_i(X) = x_i$ . The root system R is given by

$$R = \{\lambda_i - \lambda_j | i \neq j, \text{ and } i, j \in \{1, \dots, m+1\}\}.$$

A choice of positive roots is  $\lambda_i - \lambda_j$  with i < j. The set of simple roots is  $\Pi = \{\lambda_i - \lambda_{i+1} | i = 1, \ldots, m\}$ . Let  $\alpha_i = \lambda_i - \lambda_{i+1}$ . The Weyl chamber is

$$\mathfrak{a}^+ = \{ X \in \mathfrak{a} | x_1 \ge x_2 \ge \cdots \ge x_{m+1} \}.$$

The fundamental weights are  $\tilde{\omega}_{\alpha_i} = \lambda_1 + \cdots + \lambda_i$  for  $i = 1, \ldots, m$  on  $\mathfrak{b}$ . The weights  $\chi_{\alpha_i}$  has the same form as  $\tilde{\omega}_{\alpha_i}$  in  $\mathfrak{a}$ . The weights  $\omega_{\alpha_i} = \chi_{\alpha_i} - \frac{i}{m+1}(\lambda_1 + \cdots + \lambda_{m+1})$ . The representations  $V_{\alpha_i}$  are given by  $V_{\alpha_i} = \wedge^i \mathbb{R}^{m+1}$  for  $i = 1, \ldots, m$ . The maximal compact subgroup K is O(m+1) and the parabolic group P is the upper triangular subgroup and N is the subgroup of P with all the diagonal entries equal to 1. The flag variety  $\mathscr{P}$  is the set of all flags

$$W_1 \subset W_2 \subset \cdots \subset W_m,$$

where  $W_i$  is a subspace of  $\mathbb{R}^{m+1}$  of dimension *i*.

Let  $\epsilon_{i,j}$  be the square matrix of dimension m + 1 with the only nonzero entry at the *i*-th row and *j*-th column, which equals 1. The element  $H_{\alpha_i}$  is  $\epsilon_{i,i} - \epsilon_{i+1,i+1}$ . The element  $X_{\alpha_i}, Y_{\alpha_i}$  are given by  $\epsilon_{i,i+1}, \epsilon_{i+1,i}$ . The Cartan involution  $\theta$  is the additive inverse of the transpose, that is  $\theta(X) = -^t X$  for X in  $\mathfrak{a}$ .

The Weyl group W is isomorphic to the symmetric group  $\mathscr{S}_{m+1}$ . The action on  $\mathfrak{a}$  is simply given by the permutation of coordinates and the element  $w_0$  sends  $X = \text{diag}(x_1, \ldots, x_{m+1})$  to  $w_0 X = \text{diag}(x_{m+1}, \ldots, x_1)$ .

# 2.2 Linear actions on vector spaces

Let V be a vector space with euclidean norm. Then we have an induced norm on its dual space  $V^*$ , exterior powers  $\wedge^j V$  and tensor products  $\otimes^j V$ .

For  $x = \mathbb{R}v, x' = \mathbb{R}v'$  in  $\mathbb{P}V$ , we define the distance between x, x' by

$$d(x, x') = \frac{\|v \wedge v'\|}{\|v\| \|v'\|}.$$
(2.9)

This distance has the advantage that it behaves well under the action of GL(V). See for example Lemma 2.11. For  $y = \mathbb{R}f$  in  $\mathbb{P}V^*$ , let  $y^{\perp} = \mathbb{P}(\ker f) \subset \mathbb{P}V$  be a hyperplane in  $\mathbb{P}V$ . For  $x = \mathbb{R}v$  in  $\mathbb{P}V$ , we define the distance of x to  $y^{\perp}$  by

$$\delta(x, y) = \frac{|f(v)|}{\|f\| \|v\|},$$

which is explained by  $\delta(x, y) = d(x, y^{\perp}) = \min_{x' \in y^{\perp}} d(x, x')$ . Let  $K_V$  be the compact group which preserves the norm. Let  $A_V^+$  be the set of diagonal elements such that  $\{a = \operatorname{diag}(a_1, \cdots, a_d) | a_1 \ge a_2 \ge \cdots \ge a_d\}$ , under the basis  $\{e_1, \cdots, e_d\}$ . Let  $A_V^{++}$  be the interior of  $A_V^+$ . For g in GL(V), by the Cartan decomposition we can choose

$$g = k_g a_g \ell_g$$
, where  $a_g \in A_V^+$  and  $k_g, \ell_g \in K_V$ . (2.10)

Let  $x_g^M = \mathbb{R}k_g e_1$  and  $y_g^m = \mathbb{R}^t g e_1^*$  be the density points of g on  $\mathbb{P}V$  and tg on  $\mathbb{P}V^*$ , which is unique and independent of the choice of basis when  $a_g$  is in  $A_V^{++}$ . For r > 0 and g in GL(V), let

$$\begin{split} b^M_{V,g}(r) &= \{ x \in \mathbb{P}V | d(x, x^M_g) \leq r \}, \\ B^m_{V,g}(r) &= \{ x \in \mathbb{P}V | \delta(x, y^m_g) \geq r \}. \end{split}$$

These two sets play important role when we want to get some ping-pong property. The elements in set  $B_{V,g}^m$  have distance at least r to the hyperplane determined by  $y_g^m$ . For g in GL(V), let  $\gamma_{1,2}(g) := \frac{\|\wedge^2 g\|}{\|g\|^2}$  be the gap of g.

# Distance and norm

We start with general g in GL(V), where V is a finite dimensional vector space with euclidean norm. We need some technical control of distance. These are quantitative versions of the same controls in [Qui02, Lemma 2.5, 4.3, 6.5].

For g in GL(V) and  $x = \mathbb{R}v \in \mathbb{P}V$ , we define an additive cocycle  $\sigma_V : GL(V) \times \mathbb{P}V \to \mathbb{R}$  by

$$\sigma_V(g, x) = \log \frac{\|gv\|}{\|v\|}.$$
(2.11)

This is called cocycle, because for g, h in G, we have

$$\sigma_V(gh, x) = \sigma_V(g, hx) + \sigma_V(h, x).$$

We fix the operator norm  $\|\cdot\|$  on GL(V).

**Lemma 2.10.** For any g in GL(V) and x in  $\mathbb{P}V$ , we have

$$\delta(x, y_g^m) \le \frac{\|gv\|}{\|g\| \|v\|} \le 1.$$
(2.12)

Please see [BQ16, Lem 14.2] for the proof.

**Lemma 2.11.** Let  $\delta > 0$ . For g in GL(V), if  $\beta = \gamma_{1,2}(g) \leq \delta^2$ , then

• the action of g on  $B^m_{V,q}(\delta)$  is  $\beta \delta^{-2}$ -Lipschitz and

$$gB^m_{V,g}(\delta) \subset b^M_{V,g}(\beta\delta^{-1}) \subset b^M_{V,g}(\delta),$$

• the restriction of the real valued function  $\sigma_V(g, \cdot)$  on  $B^m_{V,g}(\delta)$  is  $2\delta^{-1}$ -Lipschitz.

*Proof.* Due to [BQ16, Lem 14.2],

$$d(gx, x_g^M)\delta(x, y_g^m) \le \gamma_{1,2}(g) = \beta.$$

Hence

$$d(gx, x_g^M) \le \beta \delta(x, y_g^m)^{-1} \le \beta \delta^{-1},$$

which implies the inclusion.

For  $x = \mathbb{R}v$  and  $x' = \mathbb{R}v'$  in  $B^m_{V,q}(\delta)$ , by (2.12), we have

$$d(gx, gx') = \frac{\|gv \wedge gv'\|}{\|v \wedge v'\|} \frac{\|v \wedge v'\|}{\|v\| \|v'\|} \frac{\|v\| \|v'\|}{\|gv\| \|gv'\|} \le \gamma_{1,2}(g) d(x, x') \delta^{-2},$$

which implies the Lipschitz property of g.

For the Lipschitz property of  $\sigma_V(g, \cdot)$ , please see [BQ16, Lemma 17.11].

For two different points  $x = \mathbb{R}v$  and  $x' = \mathbb{R}v'$  in  $\mathbb{P}V$ , we write  $x \wedge x' = \mathbb{R}(v \wedge v') \in \mathbb{P}(\wedge^2 V)$ . Lemma 2.12. For any g in GL(V) and two different points  $x = \mathbb{R}v, x' = \mathbb{R}v'$  in  $\mathbb{P}V$ , we have

$$\gamma_{1,2}(g)\delta(x \wedge x', y^m_{\wedge^2 g}) \le \frac{d(gx, gx')}{d(x, x')}.$$
(2.13)

*Proof.* By definition and (2.12), we have

$$d(gx, gx') = \frac{\|gv \wedge gv'\|}{\|v \wedge v'\|} \frac{\|v \wedge v'\|}{\|v\| \|v'\|} \frac{\|v\| \|v'\|}{\|gv\| \|gv'\|} \ge \gamma_{1,2}(g)\delta(x \wedge x', y^m_{\wedge^2 g})d(x, x').$$
  
is complete.

The proof is complete.

# 2.3 Actions on Flag varieties

# **Representations and Density points**

Now, suppose that V is a representation of G with a good norm. Recall that  $V^{\chi}$  is the eigenspace of the highest weight. Let  $V^*$  be the dual space of V. The representation of G on  $V^*$  is the dual representation given by: for  $g \in G$  and  $f \in V^*$ , let  $\rho^*(g)f = {}^t\rho(g^{-1})f$ . This definition gives

$$\langle \rho^*(g)f, \rho(g)v \rangle = \langle {}^t\rho(g^{-1})f, \rho(g)v \rangle = \langle f, v \rangle,$$

for f in  $V^*$  and v in V. Then the highest weight of  $V^*$  is  $\iota \chi$ . The following results explain the relation between different definitions by using combinatoric information on root systems and representations.

**Lemma 2.13.** We claim that for every irreducible representation V and weight  $\chi$ ,

$$V_{\chi,\zeta_o} = V^{w_0\chi}.\tag{2.14}$$

*Proof.* This can be verified as follows: For X in  $\mathfrak{a}$  and v in  $V^{\chi}$ ,

$$d\rho(X)\rho(w_0)v = w_0 d\rho(w_0 X)v = \chi(w_0 X)w_0 v = (w_0 \chi)(X)w_0 v.$$

The proof is complete.

**Lemma 2.14.** Let V be a proximal representation of G. Then we have

$$x_{\rho(g)}^{M} = \rho(k_g) V^{\chi} \text{ and } y_{\rho(g)}^{m} = {}^{t} \rho(\ell_g) (V^*)^{-\chi}.$$
 (2.15)

If V is irreducible, then we have

$$x_{\rho(g)}^{M} = V_{\chi,\eta_{g}^{M}} \text{ and } y_{\rho(g)}^{m} = V_{\iota\chi,\zeta_{g}^{m}}^{*}.$$

*Proof.* Let  $\{e_1, \ldots, e_d\}$  be an orthonormal basis of V composed of eigenvectors of  $\rho(A)$  such that  $e_1 \in V^{\chi}$ . Then  $\rho(A)$  is diagonal. For  $g = \exp(X) \in A^+$ , since  $\chi$  is the highest weight, we have

$$a_1 = \exp(\chi(X)) \ge a_2, \dots, a_d.$$

By the definition of a good norm,  $\rho(K)$  preserves the norm. Hence for g in G, the formula  $\rho(g) = \rho(k_g)\rho(a_g)\rho(\ell_g)$  is a decomposition which satisfies (2.10) in the previous paragraph with some permutation of  $\{e_2, \ldots, e_d\}$ . But these permutations do not change the density points. Hence we have  $x_{\rho(g)}^M = \mathbb{R}\rho(k_g)e_1 = \rho(k_g)V^{\chi}$ . If V is irreducible we have  $x_{\rho(g)}^M = V_{\chi,\eta_g^M}$ . In the dual space, we can verify that  $e_1^*$  has weight  $-\chi$ , which is the lowest weight in weights

In the dual space, we can verify that  $e_1^*$  has weight  $-\chi$ , which is the lowest weight in weights of  $V^*$ . By the same argument as in  $\mathbb{P}V$ , we have

$$y_{\rho(g)}^{m} = \mathbb{R}^{t} \rho(\ell_{g}) e_{1}^{*} = {}^{t} \rho(\ell_{g}) (V^{*})^{-\chi}$$

We also have a map from  $\mathscr{P}$  to  $\mathbb{P}V^*$ . Hence by (2.14) with representation  $V^*$  and weight  $\iota\chi$ , we know  $V^*_{\iota\chi,\zeta_o} = (V^*)^{w_0\iota\chi} = (V^*)^{-\chi}$ . For  $\zeta = g\zeta_o$  in  $\mathscr{P}$ , by definition,

$$V_{\iota\chi,\zeta}^* = g V_{\iota\chi,\zeta_o}^* = g (V^*)^{-\chi}.$$
(2.16)

Since V is irreducible, by (2.16) we have  $y_{\rho(g)}^m = {}^t \rho(\ell_g)(V^*)^{-\chi} = \rho^*(\ell_g^{-1})(V^*)^{-\chi} = V^*_{\iota\chi,\zeta_g^m}.$ 

## **Distance on Flag varieties**

For  $\alpha$  in  $\Pi$ , we abbreviate  $V_{\chi_{\alpha},\eta}, V^*_{\iota\chi_{\alpha},\zeta}$  to  $V_{\alpha,\eta}, V^*_{\alpha,\zeta}$ . For g in G, by Lemma 2.14, we find  $x^M_{\rho_{\alpha}(g)} = V_{\alpha,\eta^M_g}$  and  $y^m_{\rho_{\alpha}(g)} = V^*_{\alpha,\zeta^m_g}$ . For  $\eta, \eta'$  in  $\mathscr{P}$ , let

$$d_{\alpha}(\eta, \eta') = d(V_{\alpha, \eta}, V_{\alpha, \eta'})$$

be its distance between their images in  $\mathbb{P}V_{\alpha}$ . We define a distance on the flag variety. It is the maximal distance induced by projections,

$$d(\eta, \eta') = \max_{\alpha \in \Pi} d(V_{\alpha, \eta}, V_{\alpha, \eta'}).$$
(2.17)

We have another embedding of the flag variety

$$\mathscr{P} \to \prod_{\alpha \in \Pi} \mathbb{P}(V_{\alpha}^*).$$

For  $\zeta = k\zeta_o \in \mathscr{P}$ , by definition, we have  $V_{\alpha,\zeta}^* = kV_{\alpha,\zeta_o}^*$ . For  $\eta \in \mathscr{P}$  and  $\zeta \in \mathscr{P}$ , we set

$$\delta(\eta,\zeta) = \min_{\alpha \in \Pi} \delta(V_{\alpha,\eta}, V_{\alpha,\zeta}^*).$$

In particular, because the images of  $\eta_o, \zeta_o$  in  $\mathbb{P}V_\alpha, \mathbb{P}V_\alpha^*$  are  $V^{\chi_\alpha}, (V^*)^{-\chi_\alpha}$ , we know  $\delta(V_{\alpha,\eta_o}, V_{\alpha,\zeta_o}^*) = \delta(V^{\chi_\alpha}, (V^*)^{-\chi_\alpha}) = 1$ , and then

$$\delta(\eta_o, \zeta_o) = 1. \tag{2.18}$$

We write

$$b_{V_{\alpha},g}^{M}(r) = \{ x \in \mathbb{P}V_{\alpha} | d(x, x_{\rho_{\alpha}(g)}^{M}) \leq r \}, \\ B_{V_{\alpha},g}^{m}(r) = \{ x \in \mathbb{P}V_{\alpha} | \delta(x, y_{\rho_{\alpha}(g)}^{m}) \geq r \}.$$

They are subsets of  $\mathbb{P}V_{\alpha}$ . Write

$$b_g^M(r) = \{\eta \in \mathscr{P} | \forall \alpha \in \Pi, \ V_{\alpha,\eta} \in b_{V_\alpha,g}^M(r)\} = \{\eta \in \mathscr{P} | d(\eta, \eta_g^M) \le r\},\$$
$$B_g^m(r) = \{\eta \in \mathscr{P} | \forall \alpha \in \Pi, \ V_{\alpha,\eta} \in B_{V_\alpha,g}^m(r)\} = \{\eta \in \mathscr{P} | \delta(\eta, \zeta_g^m) \ge r\}.$$

They are subsets of  $\mathscr{P}$ .

# Distance and norms

We need a multidimensional version of the lemmas in Section 2.2. They are about the similar quantities on flag varieties. The idea is to use all the representations  $\rho_{\alpha}$ . For an element X in  $\mathfrak{b}$ , we have

$$\sup_{\alpha \in \Pi} |\chi_{\alpha}(X)| \le ||X|| \ll \sup_{\alpha \in \Pi} |\chi_{\alpha}(X)|.$$
(2.19)

Using Lemma 2.9, (2.19) and  $\sigma(g,\eta) - \kappa(g) \in \mathfrak{b}$ , we deduce the following two lemmas from Lemma 2.10 and Lemma 2.11

**Lemma 2.15.** For g in G and  $\eta$  in  $\mathscr{P}$ ,

 $\|\sigma(g,\eta) - \kappa(g)\| \ll |\log \delta(\eta,\zeta_g^m)|.$ 

For g in G and  $\alpha \in \Pi$ , by Lemma 2.9,

$$\gamma_{1,2}(\rho_{\alpha}(g)) = \frac{\|\wedge^2 \rho_{\alpha}(g)\|}{\|\rho_{\alpha}(g)\|^2} = e^{(2\chi_{\alpha} - \alpha - 2\chi_{\alpha})\kappa(g)} = e^{-\alpha\kappa(g)}.$$

Let

$$\gamma(g) = \sup_{\alpha \in \Pi} e^{-\alpha \kappa(g)}.$$
(2.20)

We call it the gap of g.

**Lemma 2.16.** Let  $\delta > 0$ . For g in G, if  $\beta = \gamma(g) = \sup_{\alpha \in \Pi} \exp(-\alpha \kappa(g)) \le \delta^2$ , then

• the action of g on  $B^m_q(\delta)$  is  $\beta \delta^{-2}$ -Lipschitz and

$$gB_g^m(\delta) \subset b_g^M(\beta\delta^{-1}) \subset b_g^M(\delta),$$

• the restriction of the  $\mathfrak{a}$ -valued function  $\sigma(g, \cdot)$  on  $B_g^m(\delta)$  is  $O(\delta^{-1})$ -Lipschitz.

These properties tell us that the action of an element g on a large set of the flag variety  $\mathscr{P}$  behaves like uniformly contracting map.

We also need to compare the distance on the projective space and the flag variety. Recall the map from  $\mathscr{P}$  to  $\mathbb{P}V$  defined in (2.4).

**Lemma 2.17.** Let  $(\rho, V)$  be an irreducible representation of G with highest weight  $\chi$ . There exists a constant C > 0 depending on the chosen norm such that for  $\eta, \eta'$  in  $\mathscr{P}$ ,

$$d(V_{\chi,\eta}, V_{\chi,\eta'}) \le Cd(\eta, \eta'). \tag{2.21}$$

The intuition is that a differentiable map between two compact Riemannian manifolds is Lipschitz. For more details, please see Corollary 5.9 in Appendix 5.2.

# 2.4 Actions on the tangent bundle of the Flag variety

In this section, we will study the action of G on the tangent bundle of  $\mathscr{P}$ . Recall that  $\mathscr{P} \simeq G/P$  is the flag variety and P = AN is a parabolic subgroup.

We first study the tangent bundle of the homogeneous space

$$\mathcal{P}_0 = G/A_e N$$

Recall that  $A_e$  is the analytical connected component of A, given by  $\exp(\mathfrak{a})$ . Note that the left action of K on  $\mathscr{P}_0$  is simply transitive (due to the Iwasawa decomposition in split case). Let  $z_o$  be the base point  $A_eN$  in  $\mathscr{P}_0$ . We can identify the left K-invariant vector fields as

$$T_{z_o}\mathscr{P}_0 = T_{z_o}(G/A_eN) \simeq \mathfrak{g}/\mathfrak{p}.$$

Hence the tangent bundle of  $\mathscr{P}_0$  has an isomorphism

$$T\mathscr{P}_0\simeq \mathscr{P}_0\times \mathfrak{g}/\mathfrak{p},$$

that is because we can identify the tangent space at  $z_o$  and  $z = kz_o$  by the left action of k. We denote by (z, Y) a point of  $T\mathscr{P}_0$  where z is in  $\mathscr{P}_0$  and Y is in  $\mathfrak{g}/\mathfrak{p}$ . We use elements in  $\mathfrak{n}^- = \bigoplus_{\alpha \in R^+} \mathfrak{g}^{-\alpha}$  as representative elements in  $\mathfrak{g}/\mathfrak{p}$ .

Then we describe the left action of G on  $T\mathscr{P}_0$ . Take Y in  $\mathfrak{g}^{-\alpha}$  and  $z = kz_o$  in  $\mathscr{P}_0$ . For g in G, by the Iwasawa decomposition we have a unique k' in K and a unique  $\sigma(g,k)$  in  $\mathfrak{a}$  such that  $gk = k'p \in k' \exp(\sigma(g,k))N$ , where  $p \in A_eN$ . Here  $\sigma(g,k)$  is understood as  $\sigma(g,k\eta_o)$ . Due to

$$gk \exp(tY)z_o = k'p \exp(tY)z_o = k' \exp(tAd_pY)z_o,$$

by taking derivative at t = 0, the left action of g on the tangent vector (z, Y) satisfies

$$L_g(z, Y) = (z', \operatorname{Ad}_p Y),$$

where  $z' = k' \eta_o$  and Ad is the adjoint action of P on  $\mathfrak{g}/\mathfrak{p}$ .

Now we restrict our attention to simple roots. Let  $\alpha$  be a simple root. Due to  $Y \in \mathfrak{g}^{-\alpha}$ , we have  $\operatorname{Ad}_N Y \subset Y + \mathfrak{a} + \mathfrak{n}$ , which implies that the unipotent part N acts trivially on  $(\mathfrak{g}^{-\alpha} + \mathfrak{p})/\mathfrak{p}$ . By  $p \in \exp(\sigma(g, k))N$ , we have

$$\operatorname{Ad}_{p}Y = \exp(-\alpha\sigma(g,k))Y \text{ on } (\mathfrak{g}^{-\alpha} + \mathfrak{p})/\mathfrak{p}.$$
(2.22)

This means that the line bundle  $\mathscr{P}_0 \times \mathfrak{g}^{-\alpha}$  is stable under the left action of G, and we call it the  $\alpha$ -bundle.

The flag variety  $\mathscr{P}$  is a quotient of  $\mathscr{P}_0$  by the right action of group M, due to  $A = MA_e$ . We use  $\pi$  to denote the quotient map. The right action of M also induces an action on the tangent bundle. For (z, Y) in  $T\mathscr{P}_0$  and m in M, by  $k \exp(tY)mz_o = km \exp(t\mathrm{Ad}_{m^{-1}}Y)z_o$ , we have

$$R_m(kz_o, Y) = (kmz_o, \mathrm{Ad}_{m^{-1}}Y).$$
 (2.23)

Descending to the quotient implies the tangent bundle of  $\mathscr{P}$  satisfies

$$T\mathscr{P}\simeq\mathscr{P}_0\times_M\mathfrak{g}/\mathfrak{p},$$

which is the quotient space of  $\mathscr{P}_0 \times \mathfrak{g}/\mathfrak{p}$  by the equivalence relation generated by the action of M, (2.23). Due to M < A, its adjoint action preserves the line  $\mathfrak{g}^{-\alpha}$  in  $\mathfrak{g}/\mathfrak{p}$ . Hence the  $\alpha$ -bundle on  $\mathscr{P}_0$  descends to a line bundle on  $\mathscr{P}$ , and we call it  $\mathscr{P}_{\alpha}$ , a subbundle of the tangent bundle. The integral curves of  $\alpha$ -bundle on  $\mathscr{P}_0$  are closed, and we call them  $\alpha$ -circles on  $\mathscr{P}_0$ . At a point  $z = kz_o$  in  $\mathscr{P}_0$ , it is given by

$$\gamma_{\alpha} : \mathbb{R} \to \mathscr{P}_0, \ t \mapsto k \exp(tK_{\alpha}) z_o. \tag{2.24}$$

This can be verified directly, because the tangent vector of the curve at time t is  $(\gamma_{\alpha}(t), K_{\alpha}) = (\gamma_{\alpha}(t), Y_{\alpha})$ , due to the definition of  $\mathfrak{g}/\mathfrak{p}$ , which belongs to the  $\alpha$ -bundle. The one parameter subgroup  $\{\exp(tK_{\alpha}) : t \in \mathbb{R}\}$  is a compact subgroup of G, which is isomorphic to SO(2). We call it  $O_{\alpha}$ .

Under the right action of M, the  $\alpha$ -circles on  $\mathscr{P}_0$  descends to the  $\alpha$ -circles on  $\mathscr{P}$ .

**Lemma 2.18.** Under the map (2.4), the image of the  $\alpha$ -circle containing  $\eta = k\eta_o$  in  $\mathbb{P}V_{\alpha}$  is the projective line generated by  $\rho_{\alpha}(k)V^{\chi_{\alpha}}$  and  $\rho_{\alpha}(k)V^{\chi_{\alpha}-\alpha}$ .

Let  $\chi$  be a dominant weight such that  $\chi(H_{\alpha}) = 0$ . Then the image of an  $\alpha$ -circle in  $\mathbb{P}V_{\chi}$  is a point.

*Proof.* Since  $\alpha$ -bundle is left K-invariant, the set of  $\alpha$ -circles are also left K-invariant. It is sufficient to consider the  $\alpha$ -circle containing  $\eta_o$ . Let  $(\rho, V)$  be an irreducible representation of highest weight  $\chi$ . By (2.24) and (2.4), the image of  $\alpha$ -circle is given by  $\rho(O_{\alpha})V^{\chi}$ .

Consider the Lie algebra  $\mathfrak{s}_{\alpha}$  generated by  $H_{\alpha}, X_{\alpha}, Y_{\alpha}$ , which is isomorphic of  $\mathfrak{sl}_2$ . For v in  $V^{\chi}$ , we have  $d\rho(H_{\alpha})v = \chi(H_{\alpha})v$ . Due to the classification of the irreducible representation of  $\mathfrak{sl}_2$ , the irreducible representation  $V_1$  of  $\mathfrak{s}_{\alpha}$  generating by  $V^{\chi}$  is of dimension  $\chi(H_{\alpha}) + 1$ .

When  $\chi$  satisfies  $\chi(H_{\alpha}) = 0$ , the above argument implies that  $V_1$  is a trivial representation and  $\rho(O_{\alpha})$  acts trivially on  $V_1$ . Hence the image of the  $\alpha$ -circle is a point.

When  $\chi = \chi_{\alpha}$ , the same argument implies  $V_1$  is of dimension 2. Another eigenspace of  $V_1$  is  $V^{\chi_{\alpha}-\alpha}$ . The group  $\rho(O_{\alpha})$  acts as SO(2) on  $V_1$ , which implies the result.

**Remark 2.19.** If we introduce the partial flag variety  $\mathscr{P}_{\Pi-\{\alpha\}}$ , then  $\alpha$ -circle is simply the fibre of the quotient map  $\mathscr{P} \to \mathscr{P}_{\Pi-\{\alpha\}}$ . This point of view also implies Lemma 2.18.

Generally, the  $\alpha$ -bundle on  $\mathscr{P}$  is non trivial in the sense of line bundle.

**Example 2.20.** Let G be  $SL_3(\mathbb{R})$ . Recall that

$$\mathfrak{a} = \{ X = \operatorname{diag}(x_1, x_2, x_3) | x_1 + x_2 + x_3 = 0, x_1, x_2, x_2 \in \mathbb{R} \},\$$

and  $\alpha_1, \alpha_2$  are two simple roots given by  $\alpha_1 = \lambda_1 - \lambda_2$  and  $\alpha_2 = \lambda_2 - \lambda_3$ . The group M is  $\{e, \operatorname{diag}(1, -1, -1), \operatorname{diag}(-1, 1, -1), \operatorname{diag}(-1, -1, 1)\} \simeq (\mathbb{Z}/2\mathbb{Z})^2$ . We have

$$\operatorname{Ad}_{\operatorname{diag}(1,-1,-1)}Y_{\alpha_1} = \alpha_1^{\sharp}(\operatorname{diag}(1,-1,-1))Y_{\alpha_1} = -Y_{\alpha_1}$$

In this case the action of M is nontrivial and it is not a normal subgroup of K = SO(3). The  $\alpha$ -bundle on  $\mathscr{P}$  restricted to an  $\alpha$ -circle is roughly a Möbius band.

In this case,  $\alpha_1$ -circles are given by  $\{W_1 \subset W_2\}$ , where  $W_2$  is a fixed two dimensional subspace of  $\mathbb{R}^3$  and  $W_1$  varies in one dimensional subspaces of  $W_2$ . On the contrary,  $\alpha_2$ -circles are given by  $\{W_1 \subset W_2\}$  with  $W_1$  fixed and  $W_2$  varying in two planes which contain  $W_1$ . From this description, we can easily see the G invariance of the set of  $\alpha$  circles.

It is better to work on  $\mathscr{P}_0$ , where the  $\alpha$ -bundle is trivial. One difficulty is that in the covering space  $\mathscr{P}_0$ , we need to capture the missing information of group M. More precisely, for h in G and z, z' in  $\mathscr{P}_0$  if  $h\pi(z), h\pi(z')$  are close, we do not know whether hz, hz' are close or not. This will be answered at the end of Section 2.5.

**Remark 2.21.** In an abstract language as in [BQ14, Lemma 4.8], we have a principal bundle  $M \to \mathscr{P}_0 \to \mathscr{P}$ , where the action of M on  $\mathscr{P}_0$  is a right action. We also have a left action of a semigroup  $\Gamma$  in G on  $\mathscr{P}_0$  and  $\mathscr{P}$  ( $\Gamma$  will be taken as  $\Gamma_{\mu}$  in our case). Suppose that we have a  $\Gamma$ -minimal set  $\Lambda_{\Gamma}$  in  $\mathscr{P}$ . The lifting of  $\Lambda_{\Gamma}$  to  $\mathscr{P}_0$  has different possibilities. Let  $\eta$  be a point in  $\Lambda_{\Gamma}$  and  $z = kz_o$  be a lifting in  $\mathscr{P}_0$ . Let  $M_z = \{m \in M | \overline{\Gamma k}m = \overline{\Gamma k}\}$ . Then we have a nice equivalence

 $\{\Gamma - minimal \text{ orbit in } \mathscr{P}_0\} \longleftrightarrow M_z \backslash M.$ 

In particular, if  $\Gamma$  is a semigroup of matrices of positive entries, then  $M_z = \{e\}$  and  $\Gamma$  has the maximal number of minimal orbits in  $\mathcal{P}_0$ .

# 2.5 The sign group

Recall the notation for reductive groups and Lie algebras. Let  $N^-$  be the subgroup with Lie algebra  $\mathfrak{n}^-$ . We have a Bruhat decomposition of the reductive group G([Bor90, 21.15]), where the main part is given by

$$N^- \times M \times A_e \times N \to G.$$

The image U is a Zariski open subset of G and the map is injective. For elements in U, we can define a map m to the group M, mapping an element g to the part of M in the Bruhat decomposition.

A part of M is given by the different connected components of G. Let

$$M_0 = M \cap G^o \text{ and } M_1 = M/M_0,$$

the quotient group. Let  $\pi_0(X)$  be the set of connected components of a topological space X and let  $\#\pi_0(X)$  be its number.

**Lemma 2.22.** Let **G** be a connected  $\mathbb{R}$ -split reductive  $\mathbb{R}$ -group. If  $\mathscr{D}$ **G** is simply connected, then  $M \cap B = M_0$ .

*Proof.* Recall that  $B = A \cap \mathscr{D}G$ . Due to  $\mathscr{D}G \supset K^o$ ,

$$M \cap B = M \cap (A \cap \mathscr{D}G) = M \cap \mathscr{D}G \supset M \cap K^o.$$

Since M is a subset of K, we see that  $M_0 = M \cap G^o = M \cap K \cap G^o = M \cap K^o$ . On the other hand, since  $\mathscr{D}\mathbf{G}$  is simply connected, the group of real points  $\mathscr{D}G$  is connected in the Lie group topology. Therefore

$$M \cap K^o = M_0 = M \cap G^o \supset M \cap \mathscr{D}G = M \cap B.$$

The proof is complete.

Let  $\mathscr{P}_1 = G/A_e BN$ .

**Lemma 2.23.** The homogeneous space  $\mathscr{P}_1$  has the same number of connected components as  $\mathscr{P}_0$ , that is  $\#\pi_0(\mathscr{P}_0) = \#\pi_0(\mathscr{P}_1) = \#\pi_0(M_1)$ , and each connected component of  $\mathscr{P}_1$  is isomorphic to  $\mathcal{P}$  as topological spaces.

*Proof.* Since  $\mathscr{D}G$  is connected, we know  $A_eBN \subset A_e\mathscr{D}G \subset G^o$ . The number of connected components of  $\mathscr{P}_1$  equals to  $\#\pi_0(G) = \#\pi_0(\mathscr{P}_0)$ .

The degree of covering  $\mathscr{P}_1 \to \mathscr{P}$  equals to

$$#(A/A_eB) = #(M/M \cap B).$$

By Lemma 2.22, we have  $M \cap B = M_0$ . Hence

$$\#(A/A_eB) = \#(M/M_0) = \#\pi_0(M_1) = \#\pi_0(G) = \#\pi_0(\mathscr{P}_1).$$

Since  $\mathscr{P}$  is connected, the result follows.

Hence, the  $M_1$  part can be determined by seeing in which connected component of G the element g is. Later, we want to know for two near elements g, g' in G, whether we have m(q) = m(q') or not. The connected component is easy to determine and in later proof we will skip the step for verifying the connected component.

In order to study the  $M_0$  part, we will use representations defined in Lemma 2.3. This is in the same spirit as the treatment of the sign group M in [Ben05]. Let  $v_{\alpha}$  be a non zero eigenvector with highest weight  $\chi_{\alpha}$  in  $V_{\alpha}$ . Let sg be the sign function on  $\mathbb{R}$ .

Lemma 2.24. For g in U, we have

$$\operatorname{sg}\langle v_{\alpha}, \rho_{\alpha}(g)v_{\alpha}\rangle = \chi_{\alpha}^{\sharp}(\mathbf{m}(g)),$$

where  $\chi^{\sharp}_{\alpha}$  is the corresponding algebraic character on A of the weight  $\chi_{\alpha}$ .

*Proof.* Since  $v_{\alpha}$  is N-invariant and the Cartan involution  $\theta$  maps  $N^{-}$  to N, by (2.6)

$$\langle v_{\alpha}, \rho_{\alpha}(N^{-}MA_{e}N)v_{\alpha} \rangle = \langle {}^{t}\rho_{\alpha}(N^{-})v_{\alpha}, \rho_{\alpha}(MA_{e}N)v_{\alpha} \rangle = \langle \rho_{\alpha}(\theta(N^{-}))v_{\alpha}, \rho_{\alpha}(MA_{e}N)v_{\alpha} \rangle$$
$$= \langle \rho_{\alpha}(N)v_{\alpha}, \rho_{\alpha}(MA_{e}N)v_{\alpha} \rangle = \langle v_{\alpha}, \rho_{\alpha}(MA_{e})v_{\alpha} \rangle.$$

The action of  $A_e$  does not change the sign, hence by Lemma 2.7 we have

$$\operatorname{sg}\langle v_{\alpha}, \rho_{\alpha}(g)v_{\alpha}\rangle = \operatorname{sg}\langle v_{\alpha}, \rho_{\alpha}(\mathbf{m}(g))v_{\alpha}\rangle = \chi_{\alpha}^{\sharp}(\mathbf{m}(g)).$$

The proof is complete.

In the simply connected case, we have  $X(\mathbf{B}) = \bigoplus_{\alpha \in \Pi} \mathbb{Z} \tilde{\omega}_{\alpha}$ . Due to  $M_0 = M \cap B$ , we know that  $\chi^{\sharp}_{\alpha}(m) = \tilde{\omega}^{\sharp}_{\alpha}(m)$  for m in  $M_0$ . Therefore

**Lemma 2.25.** The function  $\Pi_{\alpha \in \Pi} \chi_{\alpha}^{\sharp} : M_0 \to \mathbb{R}^m$  given by

$$\Pi_{\alpha\in\Pi}\,\chi^{\sharp}_{\alpha}(m) = (\chi^{\sharp}_{\alpha}(m))_{\alpha\in\Pi} \quad for \ m \in M_0,$$

is injective.

**Definition 2.26.** We define the sign function from  $G \times G$  to  $M \cup \{0\}$  by

$$\mathbf{m}(g,g') = \begin{cases} \mathbf{m}(\theta(g^{-1})g') & \text{if } \theta(g^{-1})g' \in U, \\ 0 & \text{if not,} \end{cases}$$

where g, g' are in G and  $\theta$  is the Cartan involution.



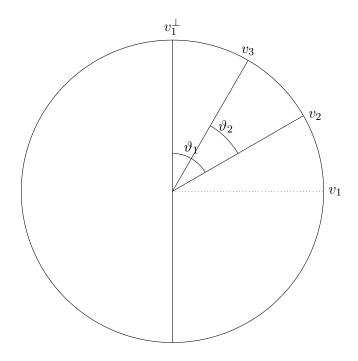


Figure 1: Angle

This definition exploits the relation between g and g'. More precisely, for u, v in  $V_{\alpha}$  we have  $\langle v, \rho_{\alpha}(\theta(g^{-1})g')u \rangle = \langle \rho_{\alpha}gv, \rho_{\alpha}g'u \rangle$ , which explains the definition. Due to  $\theta(N) = N^{-}$ , the sign function m factors through  $G/A_eN \times G/A_eN = \mathscr{P}_0 \times \mathscr{P}_0$ .

We now explain the sign function for the case m = 1, that is  $\operatorname{GL}_2(\mathbb{R})$ . We only need to consider the representation of  $\operatorname{GL}_2(\mathbb{R})$  on  $\mathbb{R}^2$ . Let  $v_0 = \begin{pmatrix} 1 \\ 0 \end{pmatrix}$  be a vector with highest weight in  $\mathbb{R}^2$ . Then

$$\langle v_0, \theta(g^{-1})g'v_0 \rangle = \langle gv_0, g'v_0 \rangle,$$

which is the inner product of the first column of g and g'. The sign function is used to determine whether these two vectors  $gv_0, g'v_0$  have an acute angle and whether g and g' are in the same connected component.

By the Bruhat decomposition, we have the following lemma.

**Lemma 2.27.** For g, g' in G and m in M, we have

$$\mathbf{m}(g, g'm) = \mathbf{m}(gm, g') = \mathbf{m}(g, g')m.$$

**Lemma 2.28.** Take a Cartan decomposition of g, that is  $g = k_g a_g \ell_g \in KA^+K$ . Then for h in G,

$$\mathbf{m}(k_g, gh) = \mathbf{m}(\ell_q^{-1}, h).$$

The key observation here is that the sign function is locally constant. Recall that  $\zeta_o$  is point in  $\mathscr{P}$  and its image in  $\mathbb{P}V_{\alpha}^*$  is the linear functional on  $V_{\alpha}$  which vanishes on the hyperplane perpendicular to  $V^{\chi_{\alpha}}$ . Recall that  $\delta(\eta, \zeta) = \min_{\alpha \in \Pi} \delta(V_{\alpha,\eta}, V_{\alpha,\zeta}^*)$  and  $d(\eta, \eta') = \max_{\alpha \in \Pi} d(V_{\alpha,\eta}, V_{\alpha,\eta'})$ .

**Lemma 2.29.** For  $k_1, k_2, k_3$  in K, if  $\delta(k_2\eta_o, k_1\zeta_o) > d(k_2\eta_o, k_3\eta_o)$ , then

$$m(k_1, k_2) = m(k_1, k_3)m(k_2, k_3)$$

*Proof.* By Lemma 2.27, it is sufficient to consider  $k_i \in K^o$ . By definition, we have  $\delta(k_2\eta_o, k_1\zeta_o) = \delta(k_1^{-1}k_2\eta_o, \zeta_o)$  and  $m(k_1, k_2) = m(id, k_1^{-1}k_2)$ . Hence, we can suppose that  $k_1 = e$ , the identity element in K. Lemma 2.24 and Lemma 2.25 imply that it is sufficient to prove that if  $\delta(k_2\eta_o, \zeta_o) > d(k_2\eta_o, k_3\eta_o)$  and  $m(k_2, k_3) = e$ , then for every simple root  $\alpha$ , we have

$$\operatorname{sg}\langle v_{\alpha}, \rho_{\alpha}(k_2)v_{\alpha}\rangle = \operatorname{sg}\langle v_{\alpha}, \rho_{\alpha}(k_3)v_{\alpha}\rangle.$$

Fix a simple root  $\alpha$  in  $\Pi$ . Abbreviate  $v_{\alpha}, \rho_{\alpha}(k_2)v_{\alpha}, \rho_{\alpha}(k_3)v_{\alpha}$  to  $v_1, v_2, v_3$ . Let  $\vartheta_1$  be the angle between the vector  $v_2$  and the hyperplane  $v_1^{\perp}$  and let  $\vartheta_2$  be the angle between  $v_2$  and  $v_3$ . Due to  $m(k_2, k_3) = e$ , this implies

$$0 < \langle v_1, k_2^{-1} k_3 v_1 \rangle = \langle k_2 v_1, k_3 v_1 \rangle = \langle v_2, v_3 \rangle,$$

the angle  $\vartheta_2$  is acute. The image of  $\zeta_0$  in  $\mathbb{P}V^*_{\alpha}$  is given by  $\mathbb{R}\langle v_1, \cdot \rangle$ . The hypothesis  $\delta(k_2\eta_o, \zeta_o) > d(k_2\eta_o, k_3\eta_o)$  implies that

$$\sin\vartheta_1 = \langle v_1, v_2 \rangle > \|v_2 \wedge v_3\| = \sin\vartheta_2.$$

Hence  $\vartheta_2 < \vartheta_1$  and  $v_2, v_3$  are in the same side of the hyperplane  $v_1^{\perp}$ , which implies  $\operatorname{sg}\langle v_1, v_2 \rangle = \operatorname{sg}\langle v_1, v_3 \rangle$ . Please see figure 1.

We state a consequence of Lemma 2.29 which will be used in Section 4.2 to get independence of certain measures  $\lambda_j$ .

**Lemma 2.30.** Let  $\delta < 1/2$ , let g, h be in G and k, k' in K. If h, k, k' satisfy

$$d(k\eta_o, k'\eta_o) < \delta, k\eta_o, k'\eta_o \in B_h^m(\delta), \eta_h^M \in B_g^m(3\delta) \text{ and } \gamma(h) < \delta^2$$

then

$$\mathbf{m}(k_g, ghk) = \mathbf{m}(\ell_q^{-1}, hk')\mathbf{m}(k, k')$$

*Proof.* By Lemma 2.27, it is sufficient to prove the case that m(k, k') = e and k, k' in  $K^o$ . By Lemma 2.28,

$$\mathbf{m}(k_g, ghk) = \mathbf{m}(\ell_q^{-1}, hk). \tag{2.25}$$

Denote  $k\eta_o, k'\eta_o$  by  $\eta, \eta'$ . Then by Lemma 2.16, we have  $h\eta, h\eta' \in b_h^M(\delta) \subset B_g^m(2\delta)$ . Hence by  $d(h\eta, h\eta') < 2\delta \leq \delta(h\eta, \zeta_q^m) = \delta(h\eta, \ell_q^{-1}\zeta_o)$  and Lemma 2.29, we have

$$m(\ell_g^{-1}, hk) = m(\ell_g^{-1}, hk')m(hk, hk').$$
(2.26)

The main point here is to prove the following lemma.

Lemma 2.31. Under the same assumption as in Lemma 2.30, we have

$$\mathbf{m}(hk, hk') = \mathbf{m}(k, k').$$

Combined with (2.25) and (2.26), the proof is complete.

Proof of Lemma 2.31. Without loss of generality, suppose that m(k, k') = e. Due to  $k\eta_o \in B_h^m(\delta)$ , we can chose a  $\ell_h$  in the Cartan decomposition  $h = k_h a_h \ell_h$  such that  $m(\ell_h^{-1}, k) = e$ . By Lemma 2.29, the hypothesis that  $\delta(k\eta_o, \ell_h^{-1}\zeta_o) > \delta > d(k\eta_o, k'\eta_o)$  implies  $m(\ell_h^{-1}, k') = m(\ell_h^{-1}, k) = e$ . By Lemma 2.28, we conclude that  $e = m(k_h, hk) = m(\ell_h^{-1}, k) = m(\ell_h^{-1}, k') = m(k_h, hk')$ . Here we need a distance  $d_0$  on  $\mathscr{P}_0$ , which is defined in Appendix 5.2. Let  $z = kz_o$  and  $z' = k'z_o$ . By Lemma 5.8,

$$d_0(hz, hz') \le d_0(hz, z_h) + d_0(z_h, hz') \le d(hk\eta_o, \eta_h^M) + d(\eta_h^M, hk'\eta_o).$$
(2.27)

Hence by (2.27), we have  $d_0(hz, hz') \leq 2\delta < 1$ , which implies m(hk, hk') = e due to Lemma 5.8.

The proof of Lemma 2.31 also says that if z, z' are close and away from the bad subvariety defined by h, the gap of h is large, then hz, hz' are also close.

# 2.6 Derivative

Let  $\varphi$  be a  $C^1$  function on  $\mathscr{P}_0$ . We will give some property of the directional derivative of  $\varphi$ . We write  $\partial_{\alpha}\varphi$  for the directional derivative  $\partial_{Y_{\alpha}}\varphi$ , where  $\alpha$  is a simple root. It turns out later that these directions are the major directions when we consider the action of G on  $\mathscr{P}_0$ .

**Definition 2.32** (Arc length). Let  $z_1, z_2$  be two points in the same  $\alpha$ -circle in  $\mathscr{P}_0$ . If  $m(z_1, z_2) = e$ , we define the arc length distance between  $z_1, z_2$  by

$$d_A(z_1, z_2) := \arcsin d(\pi z_1, \pi z_2)$$

**Remark 2.33.** This is a restriction of left K-invariant distance, which can be induced by the K-invariant Riemann metric  $d_2$  in the appendix.

**Lemma 2.34** (The Newton-Leibniz formula). Let  $z_1, z_2$  be two points in the same  $\alpha$ -circle on  $\mathscr{P}_0$  such that  $m(z_1, z_2) = e$ . Let  $u = d_A(z_1, z_2)$  and let  $\gamma : [0, u] \to \mathscr{P}_0$  be the curve in the  $\alpha$ -circle connecting  $z_1, z_2$  with unit speed (in the sense of arc length). Then for g in G

$$\varphi(gz_1) - \varphi(gz_2) = \pm \int_0^u \partial_\alpha \varphi_{g\gamma(s)} e^{-\alpha \sigma(g,\gamma(s))} \mathrm{d}s, \qquad (2.28)$$

where the sign depends on the direction of  $\gamma$ .

**Remark 2.35.** The  $\alpha$ -circle already has an orientation given by  $Y_{\alpha}$ . The sign is negative if the curve  $\gamma$  is negatively oriented.

*Proof.* Without loss of generality, suppose that  $\gamma$  is positively oriented. Recall that  $K_{\alpha} = Y_{\alpha} - X_{\alpha}$  for  $\alpha \in \Pi$ . The images of  $K_{\alpha}$  and  $Y_{\alpha}$  coincide in  $\mathfrak{g}/\mathfrak{p}$ . Then  $k_2 = k_1 \exp(uK_{\alpha})$  and  $\gamma(s) = k_1 \exp(sK_{\alpha})z_o$  for  $s \in [0, u]$ . By the Newton-Leibniz formula and (2.22) we have

$$\varphi(gz_2) - \varphi(gz_1) = \int_0^u \mathrm{d}\varphi_{g\gamma(s)} \mathrm{d}g_{\gamma(s)} K_\alpha \mathrm{d}s = \int_0^u \mathrm{d}\varphi_{g\gamma(s)} \mathrm{d}g_{\gamma(s)} Y_\alpha \mathrm{d}s$$
$$= \int_0^u \mathrm{d}\varphi_{g\gamma(s)} \exp(-\alpha\sigma(g,\gamma(s))) Y_\alpha \mathrm{d}s = \int_0^u \partial_\alpha \varphi_{g\gamma(s)} e^{-\alpha\sigma(g,\gamma(s))} \mathrm{d}s.$$

The proof is complete.

For m in M and  $\alpha$  in  $\Pi$ , by Lemma 2.7 with the adjoint representation of G on  $\mathfrak{g}$ , due to  $Y_{\alpha} \in \mathfrak{g}^{-\alpha}$ , we have  $\operatorname{Ad}_{m}Y_{\alpha} = (-\alpha)^{\sharp}(m)Y_{\alpha} = \alpha^{\sharp}(m)^{-1}Y_{\alpha} = \alpha^{\sharp}(m)Y_{\alpha}$ . The last equality is due to  $\alpha^{\sharp}(m) \in \{\pm 1\}$ . Thanks to (2.23), we have

**Lemma 2.36.** Let m be in M and let  $\varphi$  be a  $C^1$  function on  $\mathscr{P}_0$  which is right M-invariant. We have for  $z = kz_o$  in  $\mathscr{P}_0$ 

$$\partial_{\alpha}\varphi_{kmz_o} = \alpha^{\sharp}(m)\partial_{\alpha}\varphi_z.$$

We say a function  $\varphi$  on  $\mathscr{P}_0$  is the lift of a function on  $\mathbb{P}V_{\alpha}$ , if there exists a function  $\varphi_1$  on  $\mathbb{P}V_{\alpha}$  such that for  $z = kz_o \in \mathscr{P}_0$ 

$$\varphi(z) = \varphi_1(V_{\alpha,k\eta_o}).$$

By Lemma 2.18, we have

**Lemma 2.37.** If  $\varphi$  is a  $C^1$  function on  $\mathscr{P}_0$ , which is the lift of a  $C^1$  function on  $\mathbb{P}V_{\alpha}$ , then

$$\partial_{\alpha'}\varphi = 0 \text{ for } \alpha' \neq \alpha, \alpha' \in \Pi.$$

# 2.7 Changing Flags

This part is trivial for  $SL_2(\mathbb{R})$ , where the flag variety  $\mathbb{P}(\mathbb{R}^2)$  is a single  $\alpha$ -orbit. In this section, we suppose that the semisimple rank m is no less than two.

On the flag variety, we have many directions in the tangent space. Roughly speaking, the action of g is contracting and the contraction speed on  $Y_{\alpha}$  is given by  $e^{-\alpha\kappa(g)}$ ,  $\alpha \in R^+$ . Due to  $\kappa(g)$  being in the Weyl chamber  $\mathfrak{a}^+$ , the slowest directions are given by simple roots. Other directions are negligible. The main result Lemma 2.45 is a quantitative version of this intuition.

We have already seen that if two points  $\eta, \eta'$  are in the same  $\alpha$ -circle, then we have a nice formula for the difference of the value of a real function  $\varphi$  at  $g\eta$  and  $g\eta'$ , where  $g \in G$ . We want to do this for  $\eta, \eta'$  in general position. For this purpose, we need to change the point according to g. This is a key new observation in higher rank.

If we are on the euclidean space  $\mathbf{E}^n$  and we are only allowed to move along the directions of coordinate vectors. For any two points x, x', we can walk from x to x' with at most n moves. But this is not true for the flag variety  $\mathscr{P}$ . Suppose that we are only allowed to move along  $\alpha$ circles with  $\alpha \in \Pi$ . Then for two general points  $\eta, \eta'$  in  $\mathscr{P}$ , it takes more than  $m = \#\Pi$  moves to walk from one point to the other point. We try to move in each  $\alpha$  circle at most one time and to make the resulting points as close as possible.

Recall that V is a finite dimensional vector space with euclidean norm. Let  $l = \mathbb{R}(v_1 \wedge v_2)$  be a point in  $\mathbb{P}(\wedge^2 V)$ , which is also a line in  $\mathbb{P}V$ .

**Lemma 2.38.** Let  $x = \mathbb{R}w_1$  be a point in  $\mathbb{P}V$  and  $l = \mathbb{R}(v_1 \wedge v_2)$  be a line in  $\mathbb{P}V$ . Then we have

$$d(l,x) := \min_{x' \in l} d(x',x) = \frac{\|v_1 \wedge v_2 \wedge w_1\|}{\|v_1 \wedge v_2\| \|w_1\|}.$$

*Proof.* The geometric meaning of  $||v_1 \wedge v_2 \wedge w_1||$  is the volume of the parallelepiped generated by three vectors  $v_1, v_2, w_1$ . This volume can also be calculated as the product of the area of the parallelogram generated by  $v_1$  and  $v_2$ , that is  $||v_1 \wedge v_2||$ , and the distance of  $w_1$  to the plane generated by  $v_1$  and  $v_2$ , that is  $d(w_1, \text{Span}(v_1, v_2))$ . Hence, we have the formula

$$||v_1 \wedge v_2 \wedge w_1|| = ||v_1 \wedge v_2||d(w_1, \operatorname{Span}(v_1, v_2)).$$
(2.29)

The distance  $d(w_1, \text{Span}(v_1, v_2))$  equals  $||w_1|| d(l, x)$ , because the geometric sense of d(l, x) is the sine of the angle between the vector  $w_1$  and the plane  $\text{Span}(v_1, v_2)$ . Together with (2.29), we have the result .

**Lemma 2.39.** Let x be a point in  $\mathbb{P}V$  and l be a line in  $\mathbb{P}V$ . If  $g \in GL(V)$  satisfies that  $\delta(x, y_g^m), \delta(l, y_{\wedge^2 q}^m) > \delta$ , then

$$d(gl,gx) \le \delta^{-2} \gamma_{1,3}(g) d(l,x),$$

where  $\gamma_{1,3}(g) = \frac{\|\wedge^3 g\|}{\|\wedge^2 g\|\|g\|}.$ 

Compared with Lemma 2.12, with more degree of freedom the contracting speed is significantly greater.

*Proof.* By definition and  $l = \mathbb{R}(v_1 \wedge v_2), x = \mathbb{R}w_1$ , we have

$$d(gl,gx) = \frac{\|\wedge^2 g(v_1 \wedge v_2) \wedge gw_1\|}{\|\wedge^2 g(v_1 \wedge v_2)\| \|gw_1\|} \le \frac{\|\wedge^3 g\| \|v_1 \wedge v_2 \wedge w_1\|}{\|\wedge^2 g(v_1 \wedge v_2)\| \|gw_1\|},$$

Then by Lemma 2.10, we have

$$d(gl,gx) \leq \frac{\|\wedge^3 g\| \|v_1 \wedge v_2 \wedge w_1\|}{\delta^2 \|\wedge^2 g\| \|v_1 \wedge v_2\| \|g\| \|w_1\|} = \frac{\|\wedge^3 g\|}{\delta^2 \|\wedge^2 g\| \|g\|} d(l,x).$$

22

The proof is complete.

Lemma 2.39 can also be understood that there exists a point  $x' = \mathbb{R}v' \in l$  such that  $v' \wedge w_1$  is orthogonal to the vector of highest weight in  $\wedge^2 V$ . Then the distance between gx' and gx will be roughly  $\gamma_{1,3}(g)$ .

We will start to change the flags. Recall that for  $\alpha \in \Pi$  and  $\eta, \eta'$  in  $\mathscr{P}$ , the function  $d_{\alpha}(\eta, \eta')$  is the distance between the images of  $\eta$  and  $\eta'$  in  $\mathbb{P}V_{\alpha}$ . If one wants to change a flag in the  $\alpha$ -circle in  $\mathscr{P}$ , there are some constraints from the structure of flags. We introduce the following definition which explains the constraint.

By Lemma 2.18, we have

**Lemma 2.40.** The image of the  $\alpha$ -circle of  $\eta$  in  $\mathbb{P}V_{\alpha}$  is a projective line and we call it  $l_{\alpha,\eta}$ . Seen as an element in  $\mathbb{P}(\wedge^2 V_{\alpha})$ , the element  $l_{\alpha,\eta}$  is actually in  $\mathbb{P}V_{2\chi_{\alpha}-\alpha} \subset \mathbb{P}(\wedge^2 V_{\alpha})$ .

**Example 2.41.** If  $G = SL_{m+1}(\mathbb{R})$ . Let

$$\eta = \{W_1 \subset W_2 \subset \cdots \subset W_{m+1} = \mathbb{R}^{m+1}\}$$

be a flag in  $\mathscr{P}$ . Recall that  $W_r$  are r-dimensional subspaces of  $\mathbb{R}^{m+1}$ . Take  $W_0 = \{0\}$ . Let  $i_r$  be the natural embedding of the Grassmannian to projective spaces, that is  $\mathbb{G}_r(\mathbb{R}^{m+1}) \to \mathbb{P}(\wedge^r \mathbb{R}^{m+1})$ . In this case, we see that

$$l_{\alpha_r,\eta} = i_r (W_{r+1} \supset W'_r \supset W_{r-1}),$$

being a line in  $\mathbb{P}(\wedge^r \mathbb{R}^{m+1})$ , which is the image of all the r dimensional subspace  $W'_r$  of  $\mathbb{R}^{m+1}$ such that  $W_{r-1} \subset W'_r \subset W_{r+1}$ .

**Definition 2.42.** Let  $(\eta_0, \eta_1, \ldots, \eta_k)$  be a sequence of points in  $\mathscr{P}$ . We call it a chain if any consecutive elements  $\eta_i, \eta_{i+1}$  are in the same  $\alpha$ -circle for some  $\alpha \in \Pi$ , and we write  $\alpha(\eta_i, \eta_{i+1})$  for this simple root.

**Lemma 2.43.** Let  $(\eta_0, \ldots, \eta_l)$  be a chain and let  $\alpha$  be a simple root. If the set of simple roots appearing in the chain does not contain  $\alpha$ , then the image of the chain in  $\mathbb{P}V_{\alpha}$  is a single point, that is

$$V_{\alpha,\eta_i} = V_{\alpha,\eta_0}, \ \forall j = 1, \dots, l.$$

If the set of simple roots appearing in the chain also does not contain  $\alpha'$  such that  $\alpha + \alpha'$  is a root, then

$$l_{\alpha,\eta_j} = l_{\alpha,\eta_0}, \ \forall j = 1, \dots, l.$$

*Proof.* The first equality is direct consequence of Lemma 2.18 and the relation  $\chi_{\alpha}(H_{\alpha'}) = \delta_{\alpha\alpha'}$ .

For the second equality, let  $\alpha'$  be a simple root such that  $\alpha + \alpha'$  is not a root. The projective line  $l_{\alpha,\eta}$  in  $\mathbb{P}V_{\alpha}$  is uniquely determined by the image of  $\eta$  in  $\mathbb{P}V_{2\chi_{\alpha}-\alpha}$ . Hence we only need to understand the image of  $\alpha'$ -circle in  $\mathbb{P}V_{2\chi_{\alpha}-\alpha}$ . By definition,

$$(2\chi_{\alpha} - \alpha)(H_{\alpha'}) = 2\delta_{\alpha\alpha'} - \alpha(H_{\alpha'}).$$

Since  $\alpha + \alpha'$  is not a root, we know  $\alpha(H_{\alpha'}) = 0$  and  $(2\chi_{\alpha} - \alpha)(H_{\alpha'}) = 0$ . By Lemma 2.18, the image of  $\alpha'$ -circle in  $\mathbb{P}V_{2\chi_{\alpha}-\alpha}$  is point. The proof is complete.

The Coxeter diagram of an irreducible root system is a tree, module the multiplicities of edges. We can find a disjoint union  $\Pi_1$  and  $\Pi_2$  of vertices such that there is no edge whose two endpoints are in the same  $\Pi_i$ . In the Coxeter diagram, two simple roots  $\alpha$ ,  $\alpha'$  are connected by an edge if and only if  $\alpha + \alpha'$  is a root. Hence, we have

**Lemma 2.44.** We can separate  $\Pi$  into a disjoint union  $\Pi_1$  and  $\Pi_2$  such that for  $\alpha, \alpha'$  in the same atom  $\Pi_j$ ,

 $\alpha + \alpha'$  is not a root.

Let  $l_1 = \#\Pi_1$  and  $l_2 = \#\Pi_2$ .

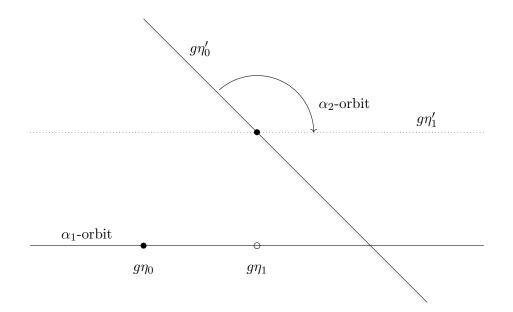


Figure 2: Changing Flag for  $SL_3(\mathbb{R})$ 

Now, we state our main result of this part, which will be used in the main approximation (Proposition 4.12).

**Lemma 2.45.** Let  $\eta, \eta'$  be two points in  $\mathscr{P}$  and let g be in G. If for  $\alpha \in \Pi_1$ ,

$$\delta(V_{\alpha,\eta'}, y^m_{\rho_\alpha(q)}), \delta(l_{\alpha,\eta}, y^m_{\wedge^2\rho_\alpha(q)}) > \delta,$$

for  $\alpha \in \Pi_2$ ,

$$\delta(V_{\alpha,\eta}, y^m_{\rho_{\alpha}(g)}), \delta(l_{\alpha,\eta'}, y^m_{\wedge^2 \rho_{\alpha}(g)}) > \delta.$$

Then we can find two chains  $(\eta = \eta_0, \eta_1, \dots, \eta_{l_1})$  and  $(\eta' = \eta'_0, \eta'_1, \dots, \eta'_{l_2})$  such that

$$d(g\eta_j, g\eta_{j+1}) = d_\alpha(g\eta_j, g\eta_{j+1}) = d_\alpha(g\eta, g\eta') + O(\delta^{-2}\beta e^{-\alpha\kappa(g)}),$$
(2.30)

where  $\alpha = \alpha(\eta_j, \eta_{j+1}) \in \Pi_1$  and different j correspond to different roots; similarly for  $\eta'$ . We also have that for all  $\alpha \in \Pi$ 

$$d_{\alpha}(g\eta_{l_1}, g\eta'_{l_2}) \le \beta e^{-\alpha\kappa(g)} \delta^{-2}, \tag{2.31}$$

where  $\beta$  is the gap of g, that is  $\beta = \gamma(g) = \max_{\alpha \in \Pi} \{e^{-\alpha \kappa(g)}\}.$ 

The point is that the contraction speed  $\beta$  implies that the term  $\delta^{-2}\beta e^{-\alpha\kappa(g)}$  is of smaller magnitude than  $e^{-\alpha\kappa(g)}$ . The objective is to walk from  $g\eta$  to  $g\eta'$  only through  $\alpha$  circles and to preserve information of distance. Since we can neglect error term, it is simpler to walk from  $g\eta$ to  $g\eta_{l_1}$  through some  $\alpha$  circles and to walk from  $g\eta'$  to  $g\eta'_{l_2}$  through the other  $\alpha$  circles, which means the corresponding simple roots are different from the first walk. After this operation, the distance between  $g\eta_{l_1}$  and  $g\eta'_{l_2}$  is negligible, due to (2.31). The distance of the move in the  $\alpha$ circle is approximately the distance between the images of  $g\eta$  and  $g\eta'$  in  $\mathbb{P}V_{\alpha}$ , due to (2.30).

Proof of Lemma 2.45. If we have already found  $(\eta_0, \ldots, \eta_j)$  and  $j < l_1$ , we want to find  $\eta_{j+1}$ . Let  $\alpha \in \Pi_1$  be a root that does not appear in the chain. Hence by Lemma 2.43,

$$V_{\alpha,\eta_j} = V_{\alpha,\eta_0} = V_{\alpha,\eta}.$$
(2.32)

Due to Lemma 2.44 and Lemma 2.43, we have further

$$l_{\alpha,\eta_j} = l_{\alpha,\eta_0} = l_{\alpha,\eta}.$$
(2.33)

We are in the situation of Lemma 2.39 with  $V = V_{\alpha}$ ,  $x = V_{\alpha,\eta'}$  and  $l = l_{\alpha,\eta}$ . Due to the hypothesis, Lemma 2.39 and Lemma 2.40, we can find  $\eta_{i+1}$  in the same  $\alpha$ -circle of  $\eta_i$  such that

$$d_{\alpha}(g\eta_{j+1},g\eta') = d(\rho_{\alpha}gV_{\alpha,\eta_{j+1}},\rho_{\alpha}gV_{\alpha,\eta'}) \le \delta^{-2}\gamma_{1,3}(\rho_{\alpha}g) \le \delta^{-2}\beta e^{-\alpha\kappa(g)}.$$
(2.34)

Hence by (2.32) and (2.34),

$$d_{\alpha}(g\eta_{j+1},g\eta_j) = d_{\alpha}(g\eta_{j+1},g\eta) = d_{\alpha}(g\eta,g\eta') + O(\delta^{-2}\beta e^{-\alpha\kappa(g)}),$$

which is (2.30). Please see Figure 2, where an element in the flag variety is represented by a projective line with a point.

We need to verify the distance between  $g\eta_{l_1}$  and  $g\eta'_{l_2}$ . Without loss of generality, suppose that  $\alpha \in \Pi_1$ . Then by Lemma 2.43, the construction and (2.34),

$$d_{\alpha}(g\eta_{l_1},g\eta'_{l_2}) = d_{\alpha}(g\eta_{l_1},g\eta') = d_{\alpha}(g\eta_{j+1},g\eta') \le \delta^{-2}\beta e^{-\alpha\kappa(g)},$$

where j is the unique number such that  $\alpha(\eta_j, \eta_{j+1}) = \alpha$ .

**Remark 2.46.** In the case of  $SL_3(\mathbb{R})$ , we know that  $\wedge^2 V_{\alpha_1}$  and  $\wedge^2 V_{\alpha_2}$  are isomorphic to  $V_{\alpha_2}$  and  $V_{\alpha_1}$ , respectively. The condition in Lemma 2.45 is equivalent to  $\eta, \eta'$  in  $B_g^m(\delta)$ .

In the case of  $\operatorname{SL}_{m+1}(\mathbb{R})$ , the representations  $V_r = \wedge^r \mathbb{R}^{m+1}$  are fundamental representation. Since  $\operatorname{SL}_{m+1}(\mathbb{R})$  is split,  $\wedge^2 V_r$  is again proximal, but may not be irreducible. In Lemma 2.65, we will proceed to give a control on  $y^m_{\wedge^2(\wedge^r q)}$ .

The condition of Lemma 2.45 is not really important, what we need is that the condition is true with a loss of exponentially small measure when we consider the random walks on G.

**Lemma 2.47.** With the same assumption and construction in Lemma 2.45, if we also have  $\eta, \eta' \in B_g^m(\delta)$ , then  $g\eta_j, g\eta'_l$  are in  $b_g^M(\beta\delta^{-O(1)})$  for  $1 \leq j \leq l_1$  and  $1 \leq l \leq l_2$ .

*Proof.* By hypothesis, Lemma 2.16 implies that  $g\eta, g\eta' \in b_q^M(\beta\delta^{-1})$ . By (2.30),

$$d(g\eta_j, g\eta_{j+1}) \le 2\beta\delta^{-1} + O(\delta^{-2}\beta e^{-\alpha\kappa(g)}) \le \beta\delta^{-O(1)}$$

Hence by induction, we have  $g\eta_j \in b_g^M(\beta \delta^{-O(1)})$  for all j. Similarly the results hold for  $g\eta'_l$ .  $\Box$ 

# 2.8 Random walks and Large deviation principles

The study of random walks on projective spaces and flag varieties are connected by representation theory.

Let X be  $\mathscr{P}$  or  $\mathbb{P}V$ , where V is an irreducible representation of G. There is a natural group action of G on X. Let  $\mu$  be a Borel probability measure on G. Then a Borel probability measure  $\nu$  on X is called  $\mu$ -stationary if

$$\nu = \mu * \nu := \int_G g_* \nu \mathrm{d} \mu(g),$$

where  $g_*\nu$  is the pushforward measure of  $\nu$  under the action of g on X.

**Lemma 2.48** (Furstenberg). Let  $\mu$  be a Zariski dense Borel probability measure on G. There exists a unique  $\mu$ -stationary probability measure  $\nu$  on the flag variety and its images in the projective spaces  $\mathbb{P}V$  are the unique  $\mu$ -stationary probability measures when V is an irreducible representation of G.

See [Fur73], [BQ16, Proposition 10.1] for more details. In order to distinguish stationary measures on different spaces, we use  $\nu_V$  to denote a  $\mu$ -stationary measure on  $\mathbb{P}V$ .

**Definition 2.49.** Let  $\mu$  be a Zariski dense Borel probability measure on G. The measure  $\mu$  has a finite exponential moment if there exists  $t_0 > 0$  such that

$$\int_G e^{t_0 \|\kappa(g)\|} \mathrm{d}\mu(g) < \infty.$$

**Remark.** This definition coincides with the definition given in the introduction for matrix groups, because in that case  $\log ||g|| = \chi \kappa(g)$  where  $\chi$  is the highest weight of a faithful representation. This  $\chi$  is in the dual cone of  $\mathfrak{a}^+$  and  $\chi(X) \gg ||X||$  for X in  $\mathfrak{a}^+$ .

**Definition 2.50.** Let  $\mu$  be a Zariski dense Borel probability measure with exponential moment on G. The Lyapunov constant  $\sigma_{\mu}$  is defined as the average of the Iwasawa cocycle

$$\sigma_{\mu} := \int_{G \times \mathscr{P}} \sigma(g, \eta) \mathrm{d}\mu(g) \mathrm{d}\nu(\eta).$$

**Lemma 2.51.** Let  $\mu$  be a Zariski dense Borel probability measure with exponential moment on G. Then the Lyapunov constant  $\sigma_{\mu}$  is in  $\mathfrak{a}^{++}$ , the interior of the Weyl chamber. Equivalently, for any simple root  $\alpha$ , we have  $\alpha(\sigma_{\mu}) > 0$ .

The maximal positivity of Lyapunov constant in Lemma 2.51 is due to Guivarc'h-Raugi [GR85] and Goldsheid-Margulis [GM89]. See [BQ16, Corollary 10.15] for more details. Lemma 2.51 will be used to show that the action of G on  $\mathscr{P}$  is contracting in Section 4.2, where the contraction speed is give by  $\beta = \sup_{\alpha \in \Pi} \{e^{-\alpha \sigma_{\mu}}\}.$ 

In following proposition, we give the large deviation principle for the Cartan projection. We keep the assumption that  $\mu$  is a Zariski dense Borel probability measure on G with a finite exponential moment.

**Proposition 2.52.** For every  $\epsilon > 0$  there exist C, c > 0 such that for all  $n \in \mathbb{N}$  and  $\eta \in \mathscr{P}$  we have

$$\mu^{*n}\{g \in G | \|\kappa(g) - n\sigma_{\mu}\| \ge n\epsilon\} \le Ce^{-c\epsilon n},$$
(2.35)

See [BQ16, Thm 13.17] for more details.

**Proposition 2.53.** If  $(\rho, V)$  is an irreducible representation of G, then for every  $\epsilon > 0$  there exist C, c such that for all x in  $\mathbb{P}V$  and y in  $\mathbb{P}V^*$  and  $n \ge 1$  we have

$$\mu^{*n} \{ g \in G | \ \delta(x, y_g^m) \le e^{-n\epsilon} \} \le C e^{-c\epsilon n},$$

$$\mu^{*n} \{ g \in G | \ \delta(x_g^M, y) \le e^{-n\epsilon} \} \le C e^{-c\epsilon n}.$$

$$(2.36)$$

See [BQ16, Prop 14.3] for more details. Attention, we need  $\rho$  to be proximal in Proposition 2.53. Here the representation is automatically proximal due to the splitness of G.

**Proposition 2.54.** For every  $\epsilon > 0$  there exist C, c such that for all  $\eta$ ,  $\eta'$  in  $\mathscr{P}$  and  $n \ge 1$  we have

$$\mu^{*n}\{g \in G | \ \delta(\eta_a^M, \zeta) \le e^{-n\epsilon}\} \le Ce^{-c\epsilon n},\tag{2.37}$$

$$\mu^{*n}\{g \in G | \ \delta(\eta, \zeta_g^m) \le e^{-n\epsilon}\} \le Ce^{-c\epsilon n}, \tag{2.38}$$

Proposition 2.54 is a multidimensional version of Proposition 2.53.

**Proposition 2.55** (Hölder regularity). If  $(\rho, V)$  is an irreducible representation of G, then there exist constants C > 0, c > 0 such that for every y in  $\mathbb{P}V^*$  and r > 0 we have

$$\nu_V(\{x \in \mathbb{P}V | \ \delta(x, y) \le r\}) \le Cr^c.$$
(2.39)

The proximality of the representation is also needed in Proposition 2.55. This result is due to Guivarc'h [Gui90]. See [BQ16, Thm 14.1] for more details. As a corollary of Proposition 2.55, we have the following.

**Corollary 2.56.** If  $(\rho, V)$  is an irreducible representation of G with highest weight  $\chi$ , then there exist constants C > 0, c > 0 such that for every y in  $\mathbb{P}V^*$  and r > 0 we have

$$\nu(\{\eta \in \mathscr{P} \mid \delta(V_{\chi,\eta}, y) \le r\}) \le Cr^c.$$
(2.40)

*Proof.* By Lemma 2.48, we have

$$\nu(\{\eta \in \mathscr{P} \mid \delta(V_{\chi,\eta}, y) \le r\}) = \nu_V(\{x \in \mathbb{P}V \mid \delta(x, y) \le r\}).$$

Hence Corollary 2.56 follows from Proposition 2.55.

All the results in this section mean that the quantities considered here are really flexible. We can always image that things happen as wished in a large probability, a very positive expectation. Bad things are near some algebraic subvariety and have exponential small measures. For later convenience, we introduce the following definition.

**Definition 2.57** (Good element). For  $n \in \mathbb{N}, \epsilon > 0$  and  $\eta, \zeta \in \mathscr{P}$ , we say that an element h is  $(n, \epsilon, \eta, \zeta)$  good if

$$\|\kappa(h) - n\sigma_{\mu}\| \le \epsilon n/C_A \text{ and } \delta(\eta, \zeta_h^m), \delta(\eta_h^M, \zeta) > 2e^{-\epsilon n/C_A}, \tag{2.41}$$

where  $C_A$  is a constant greater than 2, which is only depend on the whole group and will be determined in Lemma 2.59.

**Lemma 2.58.** We have that h is  $(n, \epsilon, \eta, \zeta)$  good outside an exponentially small set, that is to say there exist C > 0, c > 0 such that

$$\mu^{*n}$$
{h is not  $(n, \epsilon, \eta, \zeta)$  good.}  $\leq Ce^{-c\epsilon n}$ .

*Proof.* This is due to the large deviation principle (2.35), (2.37) and (2.38).

**Lemma 2.59.** Let  $\delta = e^{-\epsilon n}$  and  $\beta = \max_{\alpha \in \Pi} e^{-\alpha \sigma_{\mu} n}$ . Suppose that  $\epsilon$  is small enough such that  $\beta < \delta^3$ . If h is  $(n, \epsilon, \eta, \zeta_q^m)$  good, then

$$\gamma(h) \le \beta \delta^{-1} \le \delta^2$$
 and  $\|\sigma(gh, \eta) - \kappa(g) - n\sigma_{\mu}\| \le \epsilon n.$ 

Proof. By hypothesis,

$$\gamma(h) = \max_{\alpha \in \Pi} e^{-\alpha \kappa(h)} = \sup_{\alpha \in \Pi} e^{-\alpha n \sigma_{\mu}} e^{\alpha (n \sigma_{\mu} - \kappa(h))} \le \beta \delta^{-1},$$

if we take  $C_A$  large enough such that for all simple roots  $\alpha$  and X in  $\mathfrak{a}$ , we have  $|\alpha(X)| \leq C_A ||X||$ .

By Lemma 2.16, we have  $h\eta \in b_h^M(\gamma(h)/\delta) \subset b_h^M(\delta) \subset B_g^m(\delta)$ . Hence by Lemma 2.15

$$\begin{aligned} \|\sigma(gh,\eta) - \kappa(g) - n\sigma_{\mu}\| &= \|\sigma(g,h\eta) - \kappa(g) + \sigma(h,\eta) - n\sigma_{\mu}\| \\ &\ll |\log \delta(h\eta,\zeta_{q}^{m})| + |\log \delta(\eta,\zeta_{h}^{m})| + \|\kappa(h) - n\sigma_{\mu}\| \ll \epsilon n/C_{A}. \end{aligned}$$

Hence if  $C_A$  is large enough depending on the whole group, the inequality holds.

For later usage in Section 3, we will define another notation of goodness.

**Definition 2.60.** For  $n \in \mathbb{N}$ ,  $\epsilon > 0$  and  $\zeta \in \mathscr{P}$ , we say that an element h is  $(n, \epsilon, \zeta)$  good if

$$\|\kappa(h) - n\sigma_{\mu}\| \le \epsilon n/C_A \text{ and } \delta(\eta_h^M, \zeta) > 2e^{-\epsilon n/C_A}.$$
(2.42)

**Lemma 2.61.** Let  $\delta = e^{-\epsilon n}$  and  $\beta = \max_{\alpha \in \Pi} e^{-\alpha \sigma_{\mu} n}$ . There exists a flag  $\eta_{\alpha}$  in  $\mathscr{P}$  which is different from  $\eta_{o}$  only in its image in  $\mathbb{P}V_{\alpha}$  and

$$V_{\alpha,\eta_{\alpha}} = V^{\chi_{\alpha}-\alpha}.$$
(2.43)

If h is  $(n, \epsilon, \zeta_q^m)$  good, then for  $\eta = l_h^{-1} \eta_{\alpha}$ , we have

$$e^{\omega_{\alpha'}(\sigma(gh,\eta)-\kappa(g)-n\sigma_{\mu})} \in [\delta,\delta^{-1}] \text{ for } \alpha' \neq \alpha \text{ and } e^{\omega_{\alpha}(\sigma(gh,\eta)-\kappa(g)-n\sigma_{\mu})} \leq \beta\delta^{-1}.$$
 (2.44)

*Proof.* The existence of  $\eta_{\alpha}$  is guaranteed by Lemma 2.18. In the  $\alpha$  circle of  $\eta_o$ , there exists a point  $\eta_{\alpha}$  whose image in  $\mathbb{P}V_{\alpha}$  is exactly  $V^{\chi_{\alpha}-\alpha}$ . This is the  $\eta_{\alpha}$  that we are looking for.

Without loss of generality, we can suppose that  $l_h = e$ . The image of  $\eta_{\alpha}$  in  $\mathbb{P}V_{\alpha'}$  is the same as  $\eta_o$  if  $\alpha' \neq \alpha$ . Hence by (2.7), we have  $\omega_{\alpha'}\sigma(gh,\eta_{\alpha}) = \omega_{\alpha'}\sigma(gh,\eta_o)$  for  $\alpha' \neq \alpha$ . By (2.18), that is  $\delta(\eta_o,\zeta_o) = 1$ , the element h is  $(n,\epsilon,\eta_o,\zeta_g^m)$  good. By Lemma 2.59, we obtain the first part of (2.44).

The image of  $\eta_{\alpha}$  in  $\mathbb{P}V_{\alpha}$  is  $V^{\chi_{\alpha}-\alpha}$ , whose weight is  $\chi_{\alpha}-\alpha$ . Hence by (2.7),

$$\chi_{\alpha}\sigma(h,\eta_{\alpha}) = \log\frac{\|hv\|}{\|v\|} = \log\frac{\|\exp(\kappa(h))v\|}{\|v\|} = (\chi_{\alpha} - \alpha)\kappa(h).$$
(2.45)

By (2.7) and (2.8), we have  $\chi_{\alpha}(\sigma(g,h\eta) - \kappa(g)) \leq 0$ . Together with (2.45),

$$\chi_{\alpha}(\sigma(gh,\eta) - \kappa(g) - n\sigma_{\mu}) = \chi_{\alpha}(\sigma(g,h\eta) - \kappa(g)) + \chi_{\alpha}(\sigma(h,\eta) - n\sigma_{\mu})$$
  
$$\leq (\chi_{\alpha} - \alpha)\kappa(h) - n\chi_{\alpha}\sigma_{\mu} = -n\alpha\sigma_{\mu} + (\chi_{\alpha} - \alpha)(\kappa(h) - n\sigma_{\mu}).$$

By (2.42) and  $\chi_{\alpha} - \omega_{\alpha} \in \mathfrak{c}^*$ , the proof is complete.

This lemma tells us that by changing the image of  $\eta$  in one projective space, the value of Iwasawa cocycle only changes in that space. There is some independence of the value of Iwasawa cocycle with respect to  $\eta$ .

**Example 2.62.** In the case of  $SL_{m+1}(\mathbb{R})$ ,

$$\eta_{\alpha_d} = \{ \mathbb{R}e_1 \subset \cdots \subset \mathbb{R}e_1 \oplus \cdots \oplus \mathbb{R}e_{d-1} \subset \mathbb{R}e_1 \oplus \cdots \oplus \mathbb{R}e_{d-1} \oplus \mathbb{R}e_{d+1} \subset \cdots \},\$$

and its image in  $\wedge^d(\mathbb{R}^{m+1})$  is  $\mathbb{R}(v_1 \wedge \cdots \wedge v_{d-1} \wedge v_{d+1})$ .

Let V be a representation of G. Let  $\mathbb{G}_2(V) := \{2\text{-planes in } V\}$  be the Grassmannian variety of V. Let  $q_{\lambda} : \wedge^2 V \to \wedge^2 V$  be the G-equivalent projection on the sum of all the irreducible subrepresentations of  $\wedge^2 V$  with highest weight equal to  $\lambda$ .

**Lemma 2.63.** Let V be an irreducible representation of G with highest weight  $\chi$ . For a simple root  $\alpha$ , let  $q_{2\chi-\alpha}$  be the G-equivalent projection from  $\wedge^2 V$  to  $\wedge^2 V$ . There exists c > 0 such that for all v, v' in V

$$\sum_{\alpha \in \Pi} \|q_{2\chi-\alpha}(v \wedge v')\| \ge c \|v \wedge v'\|.$$

*Proof.* By Lemma 2.64, we know that  $\frac{\sum_{\alpha \in \Pi} \|q_{2\chi-\alpha}(v \wedge v')\|}{\|v \wedge v'\|}$ :  $\mathbb{G}_2(V) \to \mathbb{R}_{\geq 0}$  is a positive continuous function. Since  $\mathbb{G}_2(V)$  is a compact space, on which positive continuous function has a lower bound, the result follows.

The following lemma is similar to [BQ12, Lemma 3.3].

**Lemma 2.64.** With the same assumption as in Lemma 2.63, then  $\bigcap_{\alpha \in \Pi} q_{2\chi-\alpha}$  does not contain any pure wedge.

*Proof.* Let W' be the intersection of all the kernels, that is  $W' = \bigcap_{\alpha \in \Pi} q_{2\chi-\alpha}$ . The two sets  $\mathbb{G}_2(V)$  and  $\mathbb{P}W'$  are closed subset of  $\mathbb{P}(\wedge^2 V)$  and G invariant. Therefore their intersection is again a G invariant closed subvariety which is complete. Let B be the Borel subgroup of G, which is solvable. By [Bor90, Thm.10.4], the action of a solvable algebraic connected group on a complete variety has fixed points. We claim that the fixed points of B on  $\mathbb{G}_2(V)$  are the lines with the highest weight. Then the result follows by the fact that these lines do not belong to W'.

Suppose that there exit v, u in V such that  $v \wedge u$  is B invariant. We can decompose v, u as a sum  $v = \sum_{\lambda} v_{\lambda}$  and  $u = \sum_{\lambda} u_{\lambda}$ . Since we can replace v, u by bv, bu for b in B, we can suppose that the component of highest weight  $v_{\chi}$  is non zero. Since the dimension of  $V^{\chi}$  is 1, we can suppose that  $u_{\chi} = 0$ . Let  $\rho \neq \chi$  be a highest weight such that  $u_{\rho}$  is nonzero. The B invariance of  $\mathbb{R}(v \wedge u)$  also implies that the action of  $X_{\alpha}$ , for  $\alpha$  simple roots, fixes the line. Hence  $X_{\alpha}(v \wedge u) = X_{\alpha}v \wedge u + v \wedge X_{\alpha}u \in \mathbb{R}v \wedge u$ . The weight  $\chi + \rho + \alpha$  is higher than all the weights appear in  $v \wedge u$ , hence  $v_{\chi} \wedge X_{\alpha}u_{\rho} = 0$  for all simple roots  $\alpha$ . This implies that  $\rho = \chi - \alpha$  for some simple root  $\alpha$ . Therefore  $v \wedge u$  contains  $v_{\chi} \wedge u_{\chi-\alpha}$ . Since  $v \wedge u = v_{\chi} \wedge u_{\chi-\alpha}$  which is a vector of highest weight in  $\wedge^2 V$ .

We want to prove a large deviation principle for a special reducible representation. This lemma will be used in Lemma 4.11 to control  $y^m_{\wedge^2 q}$  in Lemma 2.12 and Lemma 2.45.

**Lemma 2.65.** Let V be a super proximal representation of G (Definition 2.5). For  $\epsilon > 0$  there exist C, c > 0 such that the following holds. For  $x = \mathbb{R}v, x' = \mathbb{R}v' \in \mathbb{P}V$  with  $x \neq x'$ , we have

$$\mu^{*n}\{g \in G | \delta(x \wedge x', y^m_{\wedge^2 \rho(g)}) < e^{-\epsilon n}\} \le C e^{-\epsilon \epsilon n}$$

Due to Definition 2.5, there is only one simple root  $\alpha$  such that  $q_{2\chi-\alpha}(\wedge^2 V)$  is non zero. Write  $\wedge^2 V = W \oplus W'$ , where W is the irreducible representation generated by the vector corresponding to the highest weight in  $\wedge^2 V$ , and W' is the G-invariant complementary subspace. Then  $q_{2\chi-\alpha}(\wedge^2 V) = W$ , and we write  $Pr_W = q_{2\chi-\alpha}$ .

Proof of Lemma 2.65. By (2.15), we see that a non zero vector in  $y^m_{\wedge^2 g}$  vanishes on W' and  $y^m_{\wedge^2 g}$  can be seen as an element in  $\mathbb{P}W^*$ . We only need to consider the projection of  $v \wedge v'$  onto W and use large deviation principle (2.36). By Lemma 2.63,

$$\delta(x \wedge x', y^m_{\wedge^2 g}) = \frac{|f(v \wedge v')|}{\|v \wedge v'\|} = \frac{|f(Pr_W(v \wedge v'))|}{\|Pr_W(v \wedge v')\|} \frac{\|Pr_W(v \wedge v')\|}{\|v \wedge v'\|} \ge c\delta(Pr_W(x \wedge x'), y^m_{\wedge^2 g}),$$

where f is a unit vector in  $y^m_{\wedge^2 q}$ . The proof is complete.

# **3** Non concentration condition

We want to verify the main input for the sum-product estimate, the non concentration condition. If we want to get the non concentration directly, then this becomes an effective local limit estimate, which is difficult due to the lack of spectral gap. Hence, we transfer it to the Hölder regularity of stationary measure.

For the first time read, the reader can neglect g in the left of h and think the semisimple case  $\mathrm{SL}_{m+1}(\mathbb{R})$ . The main idea of the proof is already there. Adding g is a technical step, which is needed in its application. (We only need an additional condition on  $\eta_h^M$  to control  $\kappa(gh)$ .)

## 3.1 Projective, Weak and Strong non concentration

Recall that m is the semisimple rank of G and  $\chi_1, \dots, \chi_m$  are fixed weights, where we change the subscript from  $\alpha \in \Pi$  to  $i \in \{1, \dots, m\}$ . The set  $\{\omega_i\}_{1 \leq i \leq m}$  are the extension of fundamental weights  $\tilde{\omega}_i$  to  $\mathfrak{a}$  which vanishes on  $\mathfrak{c}$  and the restriction of  $\omega_i$  and  $\chi_i$  to  $\mathfrak{b}$  coincides with  $\tilde{\omega}_i$ . Recall that  $\alpha_1, \dots, \alpha_m$  are the simple roots of  $\mathfrak{a}^*$ .

In order to distinguish different objects, we will use capital letter X to denote functions or random variables and use small letter x to denote vectors or indeterminates.

Let *L* be the  $d \times d$  square matrix which changes the basis  $(\omega_1, \dots, \omega_m)$  of  $\mathfrak{b}^*$  to the basis  $(-\alpha_1, \dots, -\alpha_m)$ , that is  $L_{ij} = -\alpha_i(H_j)$ . Then *L* is an integer matrix. Hence, we can define  $E_d$ , a rational map from  $(\mathbb{R}^*)^m$  to  $(\mathbb{R}^*)^d$ , which is given by  $y = E_d(x)$  for  $x \in (\mathbb{R}^*)^m$  where

$$y_i = \prod_{1 \le j \le m} x_j^{L_{ij}}.$$

Fix an element g in G. Let

$$X_g(n,h,\eta) = (e^{\omega_1(\sigma(gh,\eta) - \kappa(g) - n\sigma_\mu)}, \dots, e^{\omega_m(\sigma(gh,\eta) - \kappa(g) - n\sigma_\mu)}),$$
  
$$Y_g^n(h,\eta) = (e^{-\alpha_1(\sigma(gh,\eta) - \kappa(g) - n\sigma_\mu)}, \dots, e^{-\alpha_m(\sigma(gh,\eta) - \kappa(g) - n\sigma_\mu)})$$

for  $\eta$  in  $\mathscr{P}$  and h in G. By definition,  $E_d X_g(n, h, \eta)$  is the vector which is composed of the first d components of  $Y_g^n(h, \eta)$ , that is

$$p_d Y_q^n(h,\eta) = E_d X_g(n,h,\eta), \tag{3.1}$$

where  $p_d : \mathbb{R}^m \to \mathbb{R}^d$  is the map which takes a vector x of  $\mathbb{R}^m$  to the vector of  $\mathbb{R}^d$  composed of the first d components of x. In the following argument g is fixed or g equals identity. Hence we will abbreviate  $X_g, Y_q^n, Y_e^n$  to  $X, Y^n, Y_0^n$ .

will abbreviate  $X_g, Y_g^n, Y_e^n$  to  $X, Y^n, Y_0^n$ . We define an affine determinant  $A_d$  on  $(\mathbb{R}^d)^{d+1}$ . For d+1 vectors  $y^1, \dots, y^{d+1}$  in  $\mathbb{R}^d$ , let  $A_d$  be the determinant of the  $(d+1) \times (d+1)$  matrix  $\begin{pmatrix} y^1 & \cdots & y^{d+1} \\ 1 & \cdots & 1 \end{pmatrix}$ , which is the volume of the d+1-dimensional parallelogram generated by vectors  $(y^i, 1)$  for  $i = 1, \dots, d+1$ . Let  $e_i$  be the vector in  $\mathbb{R}^d$  with only i-th coordinate nonzero and equal to 1. By identifying  $e_1 \wedge \cdots \wedge e_d$  with number 1, we can also define  $A_d$  by

$$A_d(y^1, \cdots, y^{d+1}) = \sum_{1 \le i \le d+1} (-1)^{i+d+1} y^1 \wedge \cdots \wedge \widehat{y^i} \wedge \cdots \wedge y^{d+1}.$$

For d+1 vectors  $x^1, \dots, x^{d+1}$  in  $\mathbb{R}^m$ , let  $B_d$  be a rational function defined by

$$B_d(x^1, \cdots, x^{d+1}) = A_d(E_d x^1, \cdots, E_d x^{d+1}).$$

We introduce the notation

$$\mathbf{h}_{d+1} = (h_1, \ldots, h_{d+1}),$$

which is an element in  $G^{\times (d+1)}$ . Let

$$A_d^n(\mathbf{h}_{d+1},\eta) := B_d(X(n,h_1,\eta),\dots,X(n,h_{d+1},\eta)).$$

**Definition 3.1.** We say that  $\mu$  satisfies the projective non concentration (PNC) on dimension d, if for every  $\epsilon > 0$  there exist c, C > 0 such that for all n in  $\mathbb{N}$ ,  $\eta$  in  $\mathscr{P}$  and g in G

$$\sup_{a \in \mathbb{R}, v \in \mathbb{S}^{d-1}} \mu^{*n} \{ h \in G | | \langle v, Y^n(h, \eta) \rangle - a | \le e^{-\epsilon n} \} \le C e^{-\epsilon \epsilon n},$$

where v is regarded as a vector in  $\mathbb{R}^d \times \{0\}^{m-d} \subset \mathbb{R}^m$ .

More geometrically, this is equivalent to say that the measure of  $Y^n(h,\eta)$  close to an affine hyperplane is exponentially small.

**Definition 3.2.** We say that  $\mu$  satisfies the weak non concentration (WNC) on dimension d, if for every  $\epsilon > 0$  there exist c, C > 0 such that for all n in  $\mathbb{N}$ ,  $\eta$  in  $\mathscr{P}$  and g in G

 $(\mu^{*n})^{\otimes (d+2)} \{ (\mathbf{h}_{d+1}, \ell) \in G^{\times (d+2)} || A_d^n(\mathbf{h}_{d+1}, \ell\eta)| \le e^{-\epsilon n} \} \le C e^{-\epsilon \epsilon n}.$ 

**Definition 3.3.** We say that  $\mu$  satisfies the strong non concentration (SNC) on dimension d, if for every  $\epsilon > 0$  there exist c, C > 0 such that for all n in  $\mathbb{N}$ ,  $\eta$  in  $\mathscr{P}$  and g in G

 $(\mu^{*n})^{\otimes (d+1)} \{ \mathbf{h}_{d+1} \in G^{\times (d+1)} || A_d^n(\mathbf{h}_{d+1}, \eta)| \le e^{-\epsilon n} \} \le C e^{-\epsilon \epsilon n}.$ 

We will proceed by induction. When d = 0, we make the convention that  $A_0^d = 1$  and it is trivial that SNC holds. Then

- SNC on dimension  $d \Rightarrow$  WNC on dimension d (By definition)
- PNC on dimension  $d \Leftrightarrow \text{SNC}$  on dimension d (Lemma 3.7)
- WNC on dimension  $d \Rightarrow PNC$  on dimension d (Lemma 3.9)
- SNC on dimension  $d 1 \Rightarrow$  WNC on dimension d (Lemma 3.10).

In the above implications, the constants C, c will change. We can conclude

**Proposition 3.4.** Let  $\mu$  be a Zariski dense Borel probability measure on G with exponential moment. Then  $\mu$  satisfies PNC on dimension m.

#### 3.2Away from affine hyperplanes

We need a lemma of linear algebra, which relates different non concentrations. This lemma is already known from [EMO05, Lemma 7.5]. Recall that for two subsets A, B of a metric space (X, d), the distance between A and B is defined as

$$d(A,B) = \inf_{x \in A, y \in B} d(x,y).$$

**Lemma 3.5.** Let C > 0, c > 0. Let  $u_1, \dots, u_{d+1}$  be vectors in  $\mathbb{R}^d$  with length less than C. Consider the following conditions:

*i.* There exists an affine hyperplane l such that for i = 1, ..., d + 1,

$$d(u_i, l) \le c$$

*ii.* We have

$$\left\|\sum_{1\leq i\leq d+1}(-1)^{i}u_{1}\wedge\cdots\wedge\widehat{u_{i}}\wedge\cdots\wedge u_{d+1}\right\| < c,$$

where  $\hat{u_i}$  means this term is not in the wedge product.

....

*iii.* There exists i in  $\{1, \ldots, d\}$  such that

$$d(u_i, \operatorname{Span}_{\operatorname{aff}}(u_{d+1}, u_1, \dots, u_{i-1})) < c,$$

where  $\text{Span}_{\text{aff}}$  is the affine subspace generated by the elements in the bracket.

Then  $i(c) \Rightarrow ii(2^{d+1}C^{d-1}c)$ ,  $ii(c) \Rightarrow iii(c^{1/d})$  and  $iii(c) \Rightarrow i(c)$ .

*Proof.* We first transfer the affine problem to a linear problem. Let  $v_i = u_i - u_{d+1}$  for  $i = 1, \ldots, d$ . Then  $v_i$  are vectors with length less than 2C. The above three conditions are equivalent to (with change of constants in i)

i'. There exists a linear subspace l of codimension 1 such that for  $i = 1, \ldots, d$ 

$$d(v_i, l) \le c$$

ii'. We have

$$\|v_1 \wedge \cdots \wedge v_d\| < c.$$

iii'. There exists i such that

$$d(v_i, \operatorname{Span}(v_1, \dots, v_{i-1})) < c,$$

where Span is the linear subspace generated by the elements in the bracket.

 $iii'(c) \Rightarrow i'(c)$ : Let the hyperplane l be  $\operatorname{Span}(v_1, \cdots, \hat{v_i}, \cdots, v_d)$ . Then i'(c) follows from iii'(c).

 $i'(c) \Rightarrow ii'(2^d C^{d-1}c)$ : Due to i', the volume of the parallelogram generated by  $\{v_i\}_{1 \le i \le d}$  is less than  $(2C)^{d-1}2c$ , which is ii'.

 $ii'(c) \Rightarrow iii'(c^{1/d})$ : Due to the same argument as in Lemma 2.38, we have a formula of volume,

$$\|v_1 \wedge \dots \wedge v_d\| = \prod_{1 \le i \le d} d(v_i, \operatorname{Span}(v_1, \dots, v_{i-1})),$$

from which the result follows.

As a corollary, we have the following lemma, which is general and deals with random variables.

**Corollary 3.6.** Let  $X_1, \ldots, X_{d+1}$  be i.i.d. random vectors in  $\mathbb{R}^d$  bounded by C > 0. Let l be an affine hyperplane in  $\mathbb{R}^d$ . Then for any c > 0, we have

$$\mathbb{P}\{d(X_1, l) < c\}^{d+1} \le \mathbb{P}\{\|\sum_{i=1}^{d} (-1)^i X_1 \wedge \dots \wedge \hat{X}_i \wedge \dots \wedge X_{d+1}\| < 2^{d+1} C^{d-1} c\},$$
(3.2)

and

$$\mathbb{P}\{\|\sum_{1\leq i\leq d} (-1)^{i}X_{1}\wedge\cdots\wedge\hat{X}_{i}\wedge\cdots\wedge X_{d+1}\| < c\} \\
\leq \sum_{1\leq i\leq d} \mathbb{P}\{d(X_{i},\operatorname{Span}_{\operatorname{aff}}(X_{d+1},X_{1},\cdots,X_{i-1})) < c^{1/d}\}.$$
(3.3)

**Lemma 3.7.** PNC on dimension d is equivalent to SNC on dimension d.

*Proof.* Let  $X_i = E_d X(n, h_i, \eta)$  for  $i = 1, \dots, d+1$ , where  $h_i$  has distribution  $\mu^{*n}$ . Due to Lemma 2.59, with a loss of exponentially small measure, we can suppose that  $X_i$  are bounded by  $C = e^{\epsilon_2 n}$ , where  $\epsilon_2 = \epsilon/(2d)$ .

Due to (3.1), we have  $\langle v, Y^n(h,\eta) \rangle = \langle p_d v, E_d X(n,h,\eta) \rangle$ . PNC asks exactly that the probability that  $E_d X$  is close to a hyperplane is small. By (3.2), PNC on dimension d follows from SNC on dimension d.

By (3.3), SNC on dimension d follows from PNC on dimension d.

**Remark 3.8.** We explain that SNC implies the stronger form of SNC, which will be used later. Let O(d) be the orthogonal group in dimension d. The stronger form of SNC says that for any  $(\rho_1, \cdots, \rho_{d+1}) \in \mathcal{O}(d)^{\times (d+1)}$ , we have

$$(\mu^{*n})^{\otimes (d+1)} \{ \mathbf{h}_{d+1} \in G^{\times (d+1)} || A_d(\rho_1 E_d X(n, h_1, \eta), \dots, \rho_{d+1} E_d X(n, h_{d+1}, \eta)| \le e^{-\epsilon n} \} \le C e^{-\epsilon \epsilon n}.$$

By Lemma 3.7, SNC implies PNC. We adopt the notation in the proof of Lemma 3.7. By (3.3) and the fact that O(d) preserves the distance,

$$\mathbb{P}\{\|\sum_{1 \le i \le d} (-1)^i \rho_1 X_1 \wedge \cdots \widehat{\rho_i X_i} \dots \wedge \rho_{d+1} X_{d+1}\| < c\} \\ \le \sum_{1 \le i \le d} \mathbb{P}\{d(\rho_i X_i, l_i) < c^{1/d}\} = \sum_{1 \le i \le d} \mathbb{P}\{d(X_i, \rho_i^{-1} l_i) < c^{1/d}\},$$

where  $l_i = \text{Span}_{\text{aff}}(\rho_{d+1}X_{d+1}, \rho_1X_1, \cdots, \rho_{i-1}X_{i-1})$ . Therefore SNC implies the stronger form of SNC.

Lemma 3.9. WNC on dimension d implies PNC on dimension d.

WNC is weaker than SNC, because WNC is not uniform on position  $\eta$ . Let  $f(\eta)$  be  $(\mu^{*n})^{\otimes (d_2)} \{...\eta\}$  in SNC (Definition 3.3). Then WNC only asks that  $\int f(\ell\eta) d\mu^{*n}(\ell)$  is small, whereas SNC asks that  $f(\eta)$  is small for every  $\eta$ . The cocycle property is the key point to obtain an estimate uniform on position from an estimate not uniform on position.

Proof of Lemma 3.9. Let  $\delta = e^{-\epsilon n}$ . We first prove the result for 2n. Recall that h is a random variable which takes values in G with the distribution  $\mu^{*2n}$ . Let  $h = \ell_1 \ell$  such that  $\ell_1$  and  $\ell$  have distribution  $\mu^{*n}$ . Then the cocycle property implies  $Y^n(h, \eta) = Y^n(\ell_1 \ell, \eta) = Y^n(\ell_1, \ell \eta) Y_0^n(\ell, \eta)$ . Fubini's theorem implies

$$E := \sup_{a,v} \mu^{*2n} \{ h | \langle v, Y^{2n}(h,\eta) \rangle \in B(a,\delta) \}$$
  
$$\leq \int_{G} \sup_{a,v} \mu^{*n} \{ \ell_1 | \langle v, Y^n(\ell_1,\ell\eta) Y_0^n(\ell,\eta) \rangle \in B(a,\delta) \} d\mu^{*n}(\ell) .$$

The cocycle property is crucial here. Fix  $\ell$  and fix a, v. We can write

$$\langle v, Y^n(\ell_1, \ell\eta) Y_0^n(\ell, \eta) \rangle = R \langle v', Y^n(\ell_1, \ell\eta) \rangle$$

where  $R = ||v \cdot Y_0^n(\ell, \eta)|| \ge \min_{1 \le j \le d} |Y_0^n(\ell, \eta)_j|$ . Here v' is a vector of norm 1, defined by  $v' = v \cdot Y_0^n(\ell, \eta)/R$ , depending on v, l and  $\eta$ . By Lemma 2.58 and Lemma 2.59, for  $\ell$  outside an exponentially small set independent of a, v, we have  $R \ge \delta^{1/2}$ . Therefore

$$E \le \int_{G} \sup_{a,v} \mu^{*n} \{\ell_1 | \langle v, Y^n(\ell_1, \ell\eta) \rangle \in B(a, \delta^{1/2}) \} \mathrm{d}\mu^{*n}(\ell) + O_{\epsilon}(\delta^c), \tag{3.4}$$

where c > 0 comes from the large deviation principle (Lemma 2.58). By Hölder's inequality,

$$\int_{G} \sup_{a,v} \mu^{*n} \{\ell_1 | \langle v, Y^n(\ell_1, \ell\eta) \rangle \in B(a, \delta^{1/2}) \} d\mu^{*n}(\ell) 
\leq \left( \int (\sup_{a,v} \mu^{*n} \{\ell_1 | \langle v, Y^n(\ell_1, \ell\eta) \rangle \in B(a, \delta^{1/2}) \})^{d+1} d\mu^{*n}(\ell) \right)^{1/(d+1)}.$$
(3.5)

By the same argument as in Lemma 3.7

 $\sup_{a,v} \mu^{*n} \{\ell_1 | \langle v, Y^n(\ell_1, \ell\eta) \rangle \in B(a, \delta^{1/2}) \}^{d+1} \le \mu^{*(d+1)n} \{ (\mathbf{h}_{d+1}) | |A_d^n(\mathbf{h}_{d+1}, \ell\eta)| \le 2\delta^{1/4} \} + O_{\epsilon}(\delta^c).$ 

Therefore, by (3.4) and (3.5), we have

$$E^{d+1} \le \mu^{*(d+2)n} \{ (\mathbf{h}_{d+1}, \ell) || A_d^n(\mathbf{h}_{d+1}, \ell\eta) | \le 2\delta^{1/4} \} + O_{\epsilon}(\delta^c).$$

The proof for 2n ends by Definition 3.2.

It remains to prove the same result for 2n + 1. Let  $h = \ell \ell$  such that  $\ell$  has distribution  $\mu^{*(n+1)}$  and  $\ell_1$  has distribution  $\mu^{*n}$ . Following the same argument, we have

 $E^{d+1} \le \mu^{*(d+1)n+(n+1)}\{(\mathbf{h}_{d+1},\ell)||A_d^n(\mathbf{h}_{d+1},\ell\eta)| \le 2\delta^{1/4}\} + O_{\epsilon}(\delta^c).$ 

Since  $\ell$  only changes the position  $\eta$ , the uniformity of WNC implies that

$$\mu^{*(d+1)n+(n+1)}\{(\mathbf{h}_{d+1},\ell)||A_d^n(\mathbf{h}_{d+1},\ell\eta)| \le 2\delta^{1/4}\} = \int_{l_3 \in G} \mu^{*(d+2)n}\{(\mathbf{h}_{d+1},l_2)||A_d^n(\mathbf{h}_{d+1},l_2(l_3\eta))| \le 2\delta^{1/4}\} d\mu(l_3) \ll_{\epsilon} \delta^c.$$

The proof is complete.

# 3.3 Hölder regularity

In this section, we will prove

**Lemma 3.10.** SNC on dimension d - 1 implies WNC on dimension d.

Using other representations, we can get more information on the Iwasawa cocycle. This idea has already been used in [Aou13] for problem concerning transience of algebraic subvariety of split real Lie groups. It is also used in the work of Bourgain-Gamburd on the spectral gap of dense subgroups in SU(n), for establishing transience of subgroups.

The key tool is the following estimate. See [BQ16, Proposition 14.3] or [Gui90] for example.

**Lemma 3.11.** Let V be an irreducible representation of G. Let  $\mu$  be a Zariski dense Borel probability measure on G with exponential moment. For every  $\epsilon > 0$  there exist c, C > 0 such that for v in V and f in V<sup>\*</sup> we have

$$\mu^{*n}\{\ell \in G | |f(\ell v)| \le ||f|| ||\ell||e^{-\epsilon n}\} \le Ce^{-c\epsilon n}.$$

The intuition is that if a function f is not small at some point, then it is robustly large for almost all points.

In this part, we write  $V_j = V_{\chi_j}$  for the fixed representation in Lemma 2.3 and we write  $V_{j,\eta}$  for the image of  $\eta \in \mathscr{P}$  in  $\mathbb{P}V_j$  for  $j = 1, \ldots, m$ . Let  $v^j$  be a nonzero vector in  $V_{j,\eta}$ . For  $\ell$  in G, we abbreviate  $\rho_j(\ell)v^j$  to  $\ell v^j$ . Since  $v^j$  lives in  $V_j$ , we use the same symbol  $\|\cdot\|$  for norms on different  $V_j$ , which makes no confusion. For a vector x in  $\mathbb{R}^m$ , we denote by  $x_i$  the *i*-th coordinate. We use upper script to denote different vectors. We want to replace  $\omega_j$  by  $\chi_j$ , because  $\chi_j \sigma(g, \eta)$  has a nice interpretation using representations (2.8). Let  $\chi_j^c = \chi_j - \omega_j$ , which vanishes on  $\mathfrak{b}$ .

Before proving Lemma 3.10, we introduce some linear algebras. We want to construct a linear form. Recall that  $E_d$  is a rational map,  $A_d$  is the affine determinant,  $B_d$  is the composition of  $A_d$  and  $E_d$  and

$$A_d^n(\mathbf{h}_{d+1},\eta) := B_d(X(n,h_1,\eta),\dots,X(n,h_{d+1},\eta)),$$

where

$$X(n,h,\eta) = (e^{\omega_j(\sigma(gh,\eta) - \kappa(g) - n\sigma_\mu)})_{1 \le j \le m} = \left(\frac{e^{\chi_j^c(-c(h) + n\sigma_\mu)}}{e^{\chi_j(\kappa(g) + n\sigma_\mu)}} \frac{\|ghv^j\|}{\|v^j\|}\right)_{1 \le j \le m},$$
(3.6)

and the second equality is due to (2.7) and

 $\chi_j^c(\sigma(gh,\eta) - \kappa(g) - n\sigma_\mu) = \chi_j^c(c(gh) - c(g) - n\sigma_\mu) = \chi_j^c(c(h) - n\sigma_\mu).$ 

Let

$$X^{i}(n,\eta) := X(n,h_{i},\eta).$$

$$(3.7)$$

In order to use Lemma 3.11, we need to linearise some function related to  $A_d^n(\mathbf{h}_{n+1},\eta)$  with  $\mathbf{h}_{n+1}$  fixed. We will multiply  $B_d$  by its denominator, and all the Galois conjugate to get a polynomial on  $\|X_i^i\|^2$ , which can be realized as a linear functional.

The function  $B_d$  can be seen as a rational function on

$$(\underline{x}) := (x^1, \cdots, x^{d+1}) = (x^i_j)_{1 \le i \le d+1, 1 \le j \le m}.$$

By definition,  $B_d$  has a special form. Each term in  $B_d$  can be expressed as a quotient of two monomials. Let  $D_d$  be the lowest common denominator of  $B_d$  such that  $D_d B_d$  is a polynomial on (<u>x</u>). In other words, suppose that

$$B_d = \sum_{\mathbf{n} \in \mathbb{Z}^{m(d+1)}} b_{\mathbf{n}} \prod_{1 \le j \le m, 1 \le i \le d+1} (x_j^i)^{n_{ij}},$$

where **n** is a multi index and  $b_{\mathbf{n}}$  is the coefficient. Let  $q_{ij} = \sup_{\mathbf{n} \in \mathbb{Z}^{m(d+1)}} \{-n_{ij}, 0\}$  for  $1 \leq j \leq m, 1 \leq i \leq d+1$ . Then  $D_d = \prod_{1 \leq j \leq m, 1 \leq i \leq d+1} (x_j^i)^{q_{ij}}$ .

**Definition 3.12.** Let F be a polynomial on  $(x^1, \dots, x^k)$  where  $x^1, \dots, x^k$  are vectors in  $\mathbb{R}^n$ . Then we call F a multi homogeneous polynomial of degree  $\mathbf{q} = (q_1, \dots, q_n) \in \mathbb{N}^n$  if for  $\xi$  in  $(\mathbb{R}^*)^n$  we have

$$F(\xi x^1, \cdots, \xi x^k) = \xi^{\mathbf{q}} F(x^1, \cdots, x^k),$$

where  $\xi^{\mathbf{q}} = \prod_{1 \leq j \leq n} \xi_j^{q_j}$ .

Let  $\Gamma$  be the finite group  $(\mathbb{Z}/2\mathbb{Z})^{d(d+1)}$  which acts on  $\mathbb{R}^{d(d+1)}$ . Let  $(\underline{y}) := (y^1, \cdots, y^{d+1}) = (y^i_j)_{1 \le i \le d+1, 1 \le j \le d} \in (\mathbb{R}^d)^{d+1}$ . For  $\rho \in \Gamma$ , we write  $\rho(\underline{y})$  for the action on the coefficient  $y^i_j$ , which is of dimension d(d+1). Due to the definition of  $\Gamma$ , the product  $\prod_{\rho \in \Gamma} A_d \rho(y^1, \ldots, y^{d+1})$  is invariant under the action  $\Gamma$ , hence it is a polynomial on  $(y^i_j)^2$ . Let

$$F_d(x^1, \dots, x^{d+1}) = \prod_{\rho \in \Gamma} D_d A_d \rho(E_d x^1, \dots, E_d x^{d+1}),$$
(3.8)

then

**Lemma 3.13.**  $F_d$  is a multi homogeneous polynomial on  $((x^1)^2, \dots, (x^{d+1})^2)$  with degree  $\mathbf{q} = (q_1, \dots, q_m) \in \mathbb{N}^m$ .

*Proof.* We only need to verify that  $F_d$  is a multi homogeneous polynomial. The fact that the determinant is a multilinear function implies that for  $\lambda$  and  $y^i$  in  $\mathbb{R}^d$ 

$$A_d(\lambda y^1, \cdots, \lambda y^{d+1}) = \det(\lambda) A_d(y^1, \cdots, y^{d+1}), \tag{3.9}$$

where  $det(\lambda) = \lambda_1 \cdots \lambda_d$ . The functions  $E_d$  and  $D_d$  are group morphisms due to definition. Hence we have

$$E_d(\xi x) = E_d(\xi)E_d(x) \text{ and } D_d(\xi x^1, \cdots, \xi x^{d+1}) = D_d(\xi, \cdots, \xi)D_d(x^1, \cdots, x^{d+1}).$$
 (3.10)

Therefore by (3.8), (3.9) and (3.10), for  $\xi$  and  $x^i$  in  $\mathbb{R}^m$ ,

$$F_{d}(\xi x^{1}, \cdots, \xi x^{d+1}) = \prod_{\rho \in \Gamma} D_{d} A_{d} \rho(E_{d}(\xi x^{1}), \cdots, E_{d}(\xi x^{d+1}))$$
  
$$= \prod_{\rho \in \Gamma} D_{d} A_{d} \rho(E_{d}(\xi) E_{d}(x^{1}), \cdots, E_{d}(\xi) E_{d}(x^{d+1}))$$
  
$$= \prod_{\rho \in \Gamma} D_{d} A_{d} \rho(E_{d}(x^{1}), \cdots, E_{d}(x^{d+1})) \det(E_{d}(\xi)) D_{d}(\xi, \cdots, \xi)$$
  
$$= \xi^{\mathbf{q}} F_{d}(x^{1}, \cdots, x^{d+1}),$$

where **q** is a vector in  $\mathbb{N}^m$  such that  $\xi^{\mathbf{q}} = (\det(E_d(\xi))D_d(\xi, \cdots, \xi))^{|\Gamma|}$ .

For  $\mathbf{h}_{d+1} \in G^{\times (d+1)}$  and  $\eta$  in  $\mathscr{P}$ , we write

$$F(\mathbf{h}_{d+1},\eta) = F_d(X(n,h_1,\eta),\ldots,X(n,h_{d+1},\eta)).$$

Fix  $\mathbf{h}_{d+1}$ . By (3.6), F is a function on  $v^j$  for  $1 \leq j \leq m$ . Recall that  $v^j$  are vectors in  $V_{j,\eta}$ . Let

$$F_0(v^1, \cdots, v^m) = F(\mathbf{h}_{d+1}, \eta) \prod_{1 < j < m} \|v^j\|^{2q_j}$$

Now, we want to explain how to realize  $F_0$  as a linear functional.

**Lemma 3.14.** Let F be a multi homogeneous polynomial of degree  $\mathbf{q} = (q_1, \dots, q_{d+1}) \in (\mathbb{N})^{d+1}$ . Then  $F_0(v^1, \dots, v^m) := F((X^1)^2, \dots, (X^{d+1})^2) ||v^j||^{2q_j}$  is a linear functional  $F_1$  on the space  $V_0 = \bigotimes_{1 \leq j \leq m} (Sym^2V_j)^{\otimes q_j}$ , where  $X^j$  is defined in (3.7).

*Proof.* Since F is a multi homogeneous polynomial, it is sufficient to prove that every monomial in F has the same property. By Definition 3.12, a monomial of F is of the form

$$\Pi_{1 \le j \le m} \Pi_{1 \le i \le d+1} (x_j^i)^{2n_{ij}},$$

with  $n_{ij} \in \mathbb{N}$  and  $\sum_{1 \leq i \leq d+1} n_{ij} = q_j$ . The term  $\Pi \|v^j\|^{2q_j}$  is used to compensate  $\|v^j\|$  in the denominator of  $X_j^i$  in (3.6). Now, by multiplying  $\|v^j\|$ , we can view  $X_j^i$  as  $\|gh_iv^j\|$  with some coefficient. By (3.6) and  $\|ghv^j\|^2 = \langle ghv^j, ghv^j \rangle$ , the function  $(X_j^i)^2$  is a linear functional on  $Sym^2V_j$ . Hence  $\Pi_{1\leq i\leq d+1}(X_j^i)^{2n_{ij}}$  is a linear functional on  $(Sym^2V_j)^{\otimes q_j}$ . This is because if we have two linear functionals  $f_1$  and  $f_2$  on  $W_1$  and  $W_2$ , then  $f_1f_2$  is the linear functional on  $W_1 \otimes W_2$  given by  $f_1f_2(w_1 \otimes w_2) = f_1(w_1)f_2(w_2)$ . Then by the same reason, the monomial  $\Pi_{i,j}(X_i^i)^{2n_{ij}}$  is a linear functional on  $V_0$ . In order to express the linearity of  $F_0$ , we rewrite

$$F_1(\otimes_j((v^j)^2)^{\otimes q_j}) := F_0(v^1, \cdots, v^m),$$

where  $v^{j}$  is in  $V_{j,\eta}$  and  $F_{1}$  is understood as a linear functional on  $V_{0}$ .

Proof of Lemma 3.10. Recall  $\beta = \max_{\alpha \in \Pi} e^{-\alpha \sigma_{\mu} n}$ . Let  $\delta = e^{-\epsilon_2 n}$ , where the constant  $\epsilon_2$  will be determined later depending on  $\epsilon$ . We suppose that n is large enough such that  $\delta \leq 1/2$ . Because for small n, WNC can be obtained by enlarging the constant C.

**Step 1:** We take into account of measures. We want to reduce the condition of WNC on  $A_d^n$  to F, which is essentially a linear functional.

For this purpose, we will bound the measure of small  $A_d^n$  by the measure of small F.

**Lemma 3.15.** Let  $f_1, f_2$  be two Borel measurable functions on a locally compact Hausdorff space X and m be a Borel probability measure on X. Then for c > 0

$$m\{h \in X || f_1(h) \le c |\} \le m\{h \in X || f_1(h) f_2(h) | \le c \sup_X |f_2|\}.$$

In order to control  $F/A_d^n(\mathbf{h}_{d+1},\eta)$ , we take  $\mathbf{h}_{d+1}$  which is  $\eta$  good, that means for every i in  $\{1, \dots, d+1\}$ , the group element  $h_i$  is  $(n, \epsilon_2, \eta, \zeta_g^m)$  good (Definition 2.57). By Lemma 2.59 and (3.6), for  $1 \leq i \leq d+1, 1 \leq j \leq m$ 

$$|X_j^i| \le \delta^{-1}.$$

Since  $F/A_d^n$  is a polynomial on  $X_j^i$ , for  $\mathbf{h}_{d+1}$  which is  $\eta$  good, we have

$$F/A_d^n = D_d \Pi_{\rho \in \Gamma, \rho \neq e} D_d A_d^n \rho \le \delta^{-O(1)}.$$
(3.11)

Using Lemma 3.15 with  $f_1 = A_d^n$  and  $f_2 = F/A_d^n$ , hence by (3.11) and Lemma 2.58, we have

$$M := \mu^{*(d+2)n} \{ (\mathbf{h}_{d+1}, \ell) || A_d^n (\mathbf{h}_{d+1}, \ell\eta) | \le e^{-\epsilon n} \}$$
  

$$\le \mu^{*(d+2)n} \{ \mathbf{h}_{d+1} \text{ is } \ell\eta \text{ good, } \ell \in G || A_d^n (\mathbf{h}_{d+1}, \ell\eta) | \le e^{-\epsilon n} \} + O_{\epsilon_2}(\delta^c)$$
  

$$\le \mu^{*(d+2)n} \{ \mathbf{h}_{d+1} \text{ is } \ell\eta \text{ good, } \ell \in G || F(\mathbf{h}_{d+1}, \ell\eta) | \le e^{-\epsilon n} \delta^{-O(1)} \} + O_{\epsilon_2}(\delta^c)$$
  

$$\le \mu^{*(d+2)n} \{ (\mathbf{h}_{d+1}, \ell) || F(\mathbf{h}_{d+1}, \ell\eta) | \le e^{-\epsilon n} \delta^{-O(1)} \} + O_{\epsilon_2}(\delta^c).$$
(3.12)

**Step 2:** Lemma 3.13 implies that F is a multi homogeneous polynomial on  $(x_j^i)^2$  of degree  $\mathbf{q} = (q_1, \ldots, q_{d+1})$ . Lemma 3.14 implies that

$$F(\mathbf{h}_{d+1},\eta) = F_1(\otimes_j ((v^j)^2)^{\otimes q_j}) / \Pi ||v_j||^{2q_j},$$

where  $F_1$  is a linear functional on  $V_0 = \bigotimes_j (Sym^2 V_j)^{\otimes q_j}$ . To be more precise,  $F_1$  will be restricted to a linear form on W, the unique irreducible representation of  $V_0$  with maximal weight. (This is specific for real split Lie groups)

It remains to show that for most  $\mathbf{h}_{d+1}$  in  $G^{\times(d+1)}$ , the norm of  $F_1$  is robustly large. It is sufficient to find one  $\eta$  such that  $|F(\mathbf{h}_{d+1},\eta)|$  is large. We will prove that  $|D_d A_d \rho|$ is large for each  $\rho$  in  $\Gamma$ , which implies that  $|F(\mathbf{h}_{d+1},\eta)|$  is large.

Using the d + 1-th column expansion of the matrix  $\begin{pmatrix} y^1 & \cdots & y^{d+1} \\ 1 & \cdots & 1 \end{pmatrix}$ , we have

$$A_{d}(y^{1}, \cdots, y^{d+1}) = -A_{d-1}(r_{d}y^{1}, \cdots, r_{d}y^{d})y_{d}^{d+1} + \text{ other terms,}$$
  
=  $\sum_{1 \le j \le d} (-1)^{j+d+1}A_{d-1}(r_{j}y^{1}, \cdots, r_{j}y^{d})y_{j}^{d+1} + \det(y^{1}, \cdots, y^{d}),$  (3.13)

where  $r_j : \mathbb{R}^d \to \mathbb{R}^{d-1}$  is the map forgetting the *j*-th coordinate. Replacing  $y^i$  by  $E_d x^i$ , due to  $r_d E_d x^i = E_{d-1} x^i$ , we obtain

$$A_d(E_d x^1, \cdots, E_d x^{d+1}) = -A_{d-1}(E_{d-1} x^1, \cdots, E_{d-1} x^d)(E_d x^{d+1})_d + \text{ other terms.}$$
(3.14)

Using SNC on dimension d-1, we are able to give a lower bound of  $A_{d-1}(E_{d-1}X^1, \dots, E_{d-1}X^d)$ with a loss of exponentially small probability of  $\mathbf{h}_{d+1}$ . But the problem is in other similar terms. Due to  $y_j^{d+1} = \prod_{1 \le i \le m} (x_i^{d+1})^{-\alpha_j(H_i)}$  and the structure of root system, the degree of  $x_d^{d+1}$  in  $y_j^{d+1} = (E_d x^{d+1})_j$  is

$$-\alpha_d(H_d) = -2 \text{ and } -\alpha_j(H_d) \ge 0 \text{ for } j < d.$$
(3.15)

Hence, we will make  $X_d^{d+1} \leq \beta$ , which makes the first term in (3.13) greater than  $\delta^{O(1)}\beta^{-2}$ , and the other terms are less than  $\delta^{-O(1)}$ .

Now, here is the precise proof. Take  $h_{d+1}$  good, that means  $h_{d+1}$  is  $(n, \epsilon_2, \zeta_g^m)$  good (Definition 2.60). We take

$$\eta = \ell_{h_{d+1}}^{-1} \eta_{\alpha_d} \tag{3.16}$$

as in Lemma 2.61. By Lemma 2.61

$$X_j^{d+1} \in [\delta, \delta^{-1}] \text{ for } j \neq d \text{ and } X_d^{d+1} \leq \beta \delta^{-1}.$$

$$(3.17)$$

Let  $\Gamma_{d-1} = (\mathbb{Z}/2\mathbb{Z})^{(d-1)d}$ , seen as a subgroup of  $\Gamma$ , which acts on  $\mathbb{R}^{(d-1)d}$ . Then we demand that  $\mathbf{h}_d$  satisfies

$$|A_{d-1}^{n}\rho(\mathbf{h}_{d},\eta)| \ge \delta \text{ for all } \rho \in \Gamma_{d-1} \text{ and } \mathbf{h}_{d} \text{ is } \eta \text{ good.}$$
(3.18)

Recall that  $\mathbf{h}_d$  is  $\eta$  good means that  $h_i$  is  $(n, \epsilon_2, \eta, \zeta_g^m)$  good for  $1 \le i \le d$ . By Lemma 2.59 and (3.6),

$$X_j^i(\eta) \in [\delta, \delta^{-1}], \text{ for } 1 \le i \le d, 1 \le j \le m.$$
 (3.19)

Recall that W is the unique irreducible subrepresentation of  $V_0$  with the highest weight.

**Lemma 3.16.** We claim that if  $h_{d+1}$  is good  $((n, \epsilon_2, \zeta_g^m) \text{ good})$ ,  $\eta$  is taken as in (3.16) and the assumption (3.18) is satisfied for  $\mathbf{h}_d$ , then the operator norm satisfies

$$\|F_1\|_W \| \ge \delta^{O(1)}$$

Proof of Lemma 3.16. As we have already explained, it is sufficient to prove that for  $\rho$  in  $\Gamma$ , we have

$$|D_d A^n_d \rho(\mathbf{h}_d, \eta)| \ge \delta^{O(1)}.$$

The proof is similar for  $\rho$  in  $\Gamma$ , we will only prove the case  $\rho = e$ .

By (3.13) and (3.14)

$$D_{d}A_{d}(E_{d}x^{1},\cdots,E_{d}x^{d+1}) = -A_{d-1}(E_{d-1}x^{1},\cdots,E_{d-1}x^{d})D_{d}(E_{d}x^{d+1})_{d} + \sum_{1 \le j < d} (-1)^{j+d+1}A_{d-1}(r_{j}E_{d}x^{1},\cdots,r_{j}E_{d}x^{d})D_{d}(E_{d}x^{d+1})_{j} + D_{d}\det(E_{d}x^{1},\cdots,E_{d}x^{d})$$
(3.20)

where  $r_j : \mathbb{R}^d \to \mathbb{R}^{d-1}$  is the map forgetting the *j*-th coordinate. Since  $x_d^{d+1}$  only appears in  $E_d x^{d+1}$ , by (3.15), we know that the degree of  $x_d^{d+1}$  in  $D_d$  equals  $\alpha_d(H_d) = 2$ , which implies that

$$D_d \le \delta^{-O(1)} \beta^2$$

Hence by (3.17)-(3.19) and the property (3.15) that the degree of  $X_d^{d+1}$  in  $(E_d X^{d+1})_d$  is -2, the degree in  $(E_d X^{d+1})_j$  is non negative for j < d, we have

$$D_{d}(E_{d}X^{d+1})_{d} \geq \delta^{O(1)}, |A_{d-1}(E_{d-1}X^{1}, \cdots, E_{d-1}X^{d})| \geq \delta^{O(1)},$$
  

$$D_{d}(E_{d}X^{d+1})_{j} \leq \delta^{-O(1)}\beta^{2}, |A_{d-1}(r_{j}E_{d}X^{1}, \cdots, r_{j}E_{d}X^{d})| \leq \delta^{-O(1)} \text{ for } 1 \leq j < d \qquad (3.21)$$
  
and  $D_{d} \det(E_{d}X^{1}, \cdots, E_{d}X^{d}) \leq \delta^{-O(1)}\beta^{2}.$ 

By (3.20) and (3.21), we have

$$|D_d A_d^n| \ge \delta^{O(1)} - \delta^{-O(1)} \beta^2 \ge \delta^{O(1)}$$

The proof is complete.

**Step 3.** We return to the proof of Lemma 3.10. We write  $\ell v$  for the vector  $\bigotimes_j (\ell(v^j)^2)^{\bigotimes q_j}$  in  $V_0$ . Then  $\mathbb{R}lv$  is exactly the image of  $\ell\eta$  in  $\mathbb{P}W$ . Using the Fubini theorem and (3.12), we have

$$M \leq \int d\mu^{*n}(h_{d+1}) \int d\mu^{*(d-1)n}(\mathbf{h}_d) \mu^{*n} \left\{ \ell \Big| \frac{|F_1(\ell v)|}{\|F_1\|_W \| \|\ell\|} \leq e^{-\epsilon n} \delta^{-O(1)} \|F_1\|_W \|^{-1} \right\} + O_{\epsilon_2}(\delta^c).$$

Using SNC on dimension d-1, for all  $\rho \in \Gamma_{d-1}$ , we have  $\mu^{*(d-1)n}\{(\mathbf{h}_d) || A_{d-1}^n \rho(\mathbf{h}_d, \eta) | \leq \delta\} = O_{\epsilon_2}(\delta^c)$ . (This is a stronger form of SNC on dimension d-1. Due to  $\Gamma_{d-1} \in O(d-1)^{\times d}$ , it follows from Remark 3.8 that SNC implies this stronger form.) By Lemma 2.58, the set that  $h_{d+1}$  is not  $(n, \epsilon_2, \zeta_q^m)$  good and  $\mathbf{h}_d$  is not  $\eta$  good have exponentially small measure. Hence

$$M \leq \int_{good} d\mu^{*n}(h_{d+1}) \int_{\mathbf{h}_{d} \text{ satisifes (3.18)}} d\mu^{*(d-1)n}(\mathbf{h}_{d}) \mu^{*n} \left\{ \ell \Big| \frac{|F_{1}(\ell v)|}{\|F_{1}\|_{W} \| \|\ell\|} \leq e^{-\epsilon n} \delta^{-O(1)} \|F_{1}\|_{W} \|^{-1} \right\} + O_{\epsilon_{2}}(\delta^{c}).$$

$$(3.22)$$

Due to Lemma 3.16, when  $\epsilon_2$  is small enough with respect to  $\epsilon$ , we have  $(\delta = e^{-\epsilon_2 n}$  and  $||F_1|_W|| \ll \delta^{-O(1)})$ 

$$e^{-\epsilon n} \delta^{-O(1)} \|F_1\|_W \|^{-1} \le e^{-\epsilon n} \delta^{-O(1)} \le e^{-\epsilon n/2}.$$

Using Lemma 3.11 with V = W, due to  $\ell v$  in W we conclude that under the condition of Lemma 3.16,

$$\mu^{*n} \left\{ \ell \Big| \frac{F_1(\ell v)}{\|F_1\|_W \| \|\ell\|} \le e^{-\epsilon n} \delta^{-O(1)} \|F_1\|_W \|^{-1} \right\} \le_{\epsilon} e^{-c\epsilon n}.$$
(3.23)  
the proof is complete.

By (3.22) and (3.23), the proof is complete.

### 3.4 Combinatoric tool

**Proposition 3.17.** Fix  $\kappa_1 > 0$ . Let  $C_0 > 0$ . Then there exist  $\epsilon_3$  and  $k \in \mathbb{N}, \epsilon > 0$  depending only on  $\kappa_1$  such that the following holds for  $\tau$  large enough depending on  $C_0$ . Let  $\lambda_1, \ldots, \lambda_k$  be Borel measures on  $([-\tau^{\epsilon_4}, -\tau^{-\epsilon_4}] \cup [\tau^{-\epsilon_4}, \tau^{\epsilon_4}])^m \subset \mathbb{R}^m$  where  $\epsilon_4 = \min\{\epsilon_3, \epsilon_3 \kappa_0\}/10k$ , with total mass less than 1. Assume that for all  $\rho \in [\tau^{-2}, \tau^{-\epsilon_3}]$  and  $j = 1, \ldots, k$ 

$$\sup_{a \in \mathbb{R}, v \in \mathbb{S}^{m-1}} (\pi_v)_* \lambda_j (B_{\mathbb{R}}(a, \rho)) = \sup_{a, v} \lambda_j \{ x \mid \langle v, x \rangle \in B_{\mathbb{R}}(a, \rho) \} \le C_0 \rho^{\kappa_1}.$$
(3.24)

Then for all  $\varsigma \in \mathbb{R}^m$ ,  $\|\varsigma\| \in [\tau^{3/4}, \tau^{5/4}]$  we have

$$\left|\int \exp(i\langle\varsigma, x_1\cdots x_k\rangle) \mathrm{d}\lambda_1(x_1)\cdots \mathrm{d}\lambda_k(x_k)\right| \leq \tau^{-\epsilon_3}$$

This is proved in [LI18b], based on a discretized sum-product estimate by He-de Saxcé [HdS18]. When n = 1, this is due to Bourgain in [Bou10].

#### 3.5 Application to our measure

From Proposition 3.4, we fix  $\epsilon_2 < \frac{1}{10} \min_{\alpha \in \Pi} \{\alpha \sigma_\mu\}$  and we can find  $c_1$  such that PNC holds. Let  $(\epsilon_2/2, c')$  be the constants in Lemma 2.59. Take

$$\kappa_0 = \frac{1}{10} \min\{c_1, c'\}$$

Using Proposition 3.17 with  $\kappa_1 = \kappa_0$ , we get  $\epsilon_3, \epsilon_4$ .

For g, h in G and  $\eta$  in  $\mathscr{P}$ , recall that  $Y^n(h, \eta) = (e^{-\alpha(\sigma(gh,\eta) - \kappa(g) - n\sigma_{\mu})})_{\alpha \in \Pi} \in \mathbb{R}^m$ . Let  $\lambda_{g,\eta}$  be a pushforward measure on  $\mathbb{R}^m$  of  $\mu^{*n}$  restricted on a subset  $G_{n,g,\eta}$  of G, which is defined by

$$\lambda_{g,\eta}(E) = \mu^{*n} \{ h \in G_{n,g,\eta} | Y^n(h,\eta) \in E \},\$$

for any Borel subset E of  $\mathbb{R}^m$ , where

$$G_{n,g,\eta} = \{h \in G | h \text{ is } (n,\epsilon,\eta,\zeta_q^m) \text{ good}\}$$
(3.25)

and where  $\epsilon_{\mu} \geq \epsilon > 0$  will be determined later.

PNC is only at one scale, we need to verify all the scales needed in the sum-product estimate. The idea is to separate the random variable and try to use PNC in other scale, where we need the cocycle property to change scale.

**Proposition 3.18** (Change scale). With  $\epsilon$  small enough depending on  $\epsilon_4 \epsilon_2$ , there exists  $C_0$  independent of n such that the measure  $\lambda_{g,\eta}$  satisfies the conditions in Proposition 3.17 with constant  $\tau = e^{\epsilon_2 n}$  for all  $n \in \mathbb{N}$ .

*Proof.* We abbreviate  $\lambda_{g,\eta}$  to  $\lambda$ . By taking  $\epsilon$  small depending on  $\epsilon_4 \epsilon_2$ , Lemma 2.59 implies that the support of  $\lambda$  is contained in the cube  $[\tau^{-\epsilon_4}, \tau^{\epsilon_4}]^m$ .

Then we verify (3.24). Let  $\rho \in [\tau^{-2}, \tau^{-\epsilon_3}]$ . Let  $n_1 = [\frac{|\log \rho|}{2\epsilon_2}]$ . and  $n_2 = n - n_1$ . Then  $n_1$  lies in  $[\epsilon_3 n/2, n]$ . We separate  $h = h_1 h_2$  such that  $h_1, h_2$  have distributions  $\mu^{*n_1}, \mu^{*(n-n_1)}$ , respectively. We have

$$Y^{n}(h,\eta) = Y^{n_{1}}(h_{1},h_{2}\eta)Y_{0}^{n_{2}}(h_{2},\eta), \qquad (3.26)$$

We can not use the cocycle property directly to change the scale. The problem is in (3.26), where the term  $Y_0^{n_2}$  behaves bad if  $n_2 \gg n_1$ , that is to say that the probability of  $h_2$  such that  $Y_0^{n_2}(h_2,\eta)$  is smaller than  $\rho = e^{-2\epsilon_2 n_1}$  is large. In order to overcome this difficulty, we use the support of  $Y^n$ . We will prove that if  $Y_0^{n_2}$  is too small, then the support of  $Y^n$  will force  $Y^{n_1}$  to become large, which can be controlled by the large deviation principle.

Now we give the details of the proof. For (3.24), due to the fact that the support of  $\lambda$  is contained in  $[\tau^{-\epsilon_4}, \tau^{\epsilon_4}]^m$ , we have

$$(\pi_w)_*\lambda(B(a,\rho)) \le \sup_{h_2,v} \mu^{*n_1}\{h_1 | \langle v, Y^{n_1}(h_1,h_2\eta) \rangle \in R^{-1}B(a,\rho), Y^n(h_1h_2,\eta) \in [\tau^{-\epsilon_4},\tau^{\epsilon_4}]^m\},$$
(3.27)

where  $R = ||wY_0^{n_2}(h_2, \eta)||$  depends on  $h_2$ .

• If  $R \ge \rho^{1/2}$ , then  $\rho R^{-1} \le \rho^{1/2} = e^{-\epsilon_2 n_1}$ . It follows by PNC at scale  $n_1$  that

$$\mu^{*n_1}\{h_1|\langle v, Y^{n_1}(h_1, h_2\eta)\rangle \in B(a, e^{-\epsilon_2 n_1})\} \ll_{\epsilon_2} e^{-c_1\epsilon_2 n_1} \le \rho^{\kappa_0}.$$
(3.28)

• If  $R \leq \rho^{1/2}$ . There exists one coordinate  $\alpha$  such that  $|Y_0^{n_2}(h_2,\eta)_{\alpha}| \leq \rho^{1/2}$ , which implies that  $Y^{n_1}(h_1,h_2\eta)_{\alpha} = Y^n(h,\eta)_{\alpha}/Y_0^{n_2}(h_2,\eta)_{\alpha} \geq \tau^{-\epsilon_4}\rho^{-1/2}$ . Due to  $\epsilon_3 \geq 4\epsilon_4$  and  $n_1 \geq \epsilon_3 n/2$ , we have  $\epsilon_2 n_1 \geq 2\epsilon_4 \epsilon_2 n$ . Therefore  $\tau^{-\epsilon_4}\rho^{-1/2} = \tau^{-\epsilon_4}e^{\epsilon_2 n_1} \geq e^{\epsilon_2 n_1/2}$ . For such  $h_2$ , we have

$$\mu^{*n_1}\{h_1|Y^{n_1}(h_1h_2,\eta)\in[\tau^{-\epsilon_4},\tau^{\epsilon_4}]^m\}\leq \sum_{\alpha\in\Pi}\mu^{*n_1}\{h_1|Y^{n_1}(h_1,h_2\eta)_\alpha\geq e^{\epsilon_2n_1/2}\}.$$
 (3.29)

It follows from Lemma 2.59 that

$$\mu^{*n_1}\{h_1|Y^{n_1}(h_1,h_2\eta)_{\alpha} \ge e^{\epsilon_2 n_1/2}\} \le \mu^{*n_1}\{h_1|\|\sigma(gh_1,h_2\eta) - \kappa(g) - n_1\sigma_{\mu}\| \ge \epsilon_2 n_1/2\}$$
  
$$\ll_{\epsilon_2} e^{-c'\epsilon_2 n_1} \le \rho^{\kappa_0}.$$
(3.30)

By (3.27)-(3.30), for  $\rho \in [\tau^{-2}, \tau^{-\epsilon_3}]$  we have

$$(\pi_w)_*\lambda(B(a,\rho)) \ll_{\epsilon_2} \rho^{\kappa_0}$$

The proof is complete.

# 4 Proof of the main theorems

In this section, we will use the results of Section 2 and Section 3 to give the proofs of the main theorems. In Section 4.2, we will prove Theorem 1.7, the simply connected case. For non simply connected case, please see Theorem 5.4 in Appendix 5.1. Then in Section 4.3-4.4, we will work on semisimple case and we prove all the other theorems in the introduction from Theorem 5.4.

We will add many assumptions on the elements of G and  $\mathscr{P}$ . The assumptions seem complicate. In fact, they are not really important. They are taken to make the result work outside a set of exponentially small measure. These assumptions say that the elements are away from certain closed subvarieties of G or  $\mathscr{P}$ , which also explains that they are true almost everywhere.

# 4.1 (C, r) good function

For a  $C^1$  function  $\varphi$  on the flag variety  $\mathscr{P}$ . We first lift it to  $\mathscr{P}_0 = G/A_e N$ . Let  $\partial_\alpha \varphi = \partial_{Y_\alpha} \varphi$ be the directional derivative on  $\mathscr{P}_0$ . By Lemma 2.36 the action of the group M only changes the sign of the directional derivative  $\partial_\alpha \varphi$ , hence  $|\partial_\alpha \varphi|$  is actually a function on  $\mathscr{P}$ . Although  $\partial_\alpha \varphi$  is not well-defined on  $\mathscr{P}$ , we can fix a local trivialization of the line bundle  $P_\alpha$  and define the directional derivative. This point of view will be used in G3.

Recall that for  $\eta, \eta'$  in  $\mathscr{P}$  and simple root  $\alpha$ , we have defined  $d_{\alpha}(\eta, \eta') = d(V_{\alpha,\eta}, V_{\alpha,\eta'})$ .

**Definition 4.1.** Let r be a continuous function on  $\mathscr{P}$ . Let J be the open set in  $\mathscr{P}$ , which is the 1/C-neighbourhood of the support of r. Let  $\varphi$  be a  $C^2$  function on  $\mathscr{P}$ . For a simple root  $\alpha$ , let  $v_{\alpha} = \sup_{\eta \in \text{suppr}} |\partial_{\alpha}\varphi(\eta)|$ . We say that  $\varphi$  is (C, r) good if:

(G1) For  $\eta, \eta'$  in J such that  $d(\eta, \eta') \leq 1/C$ ,

$$|\varphi(\eta) - \varphi(\eta')| \le C \sum_{\alpha \in \Pi} d_{\alpha}(\eta, \eta') v_{\alpha}, \tag{4.1}$$

(G2) For every simple root  $\alpha$  and for every  $\eta$  in the support of r, we have

$$|\partial_{\alpha}\varphi(\eta)| \ge \frac{1}{C}v_{\alpha},\tag{4.2}$$

(G3) For  $\eta, \eta'$  in J with  $d(\eta, \eta') \leq 1/C$ ,

$$|\partial_{\alpha}\varphi(\eta) - \partial_{\alpha}\varphi(\eta')| \le Cd(\eta, \eta')v_{\alpha}.$$
(4.3)

(G4)

$$\sup_{\alpha \in \Pi} v_{\alpha} \in [1/C, C].$$
(4.4)

**Remark 4.2.** The distance  $d_{\alpha}$  does not depend on the representation. For two different representation  $(\rho, V), (\rho', V')$  such that  $\Theta(\rho) = \Theta(\rho') = \{\alpha\}$ , by Lemma 5.7, when C is small enough, two distances  $d_V, d_{V'}$  are equivalent.

In the above definition, the G3 assumption (4.3) is equivalent to the inequality on  $\mathscr{P}_0$ , that is

$$\left|\partial_{\alpha}\varphi(z) - \partial_{\alpha}\varphi(z')\right| \le Cd_0(z, z')v_{\alpha},\tag{4.5}$$

for z, z' in  $\pi^{-1}(J)$  with  $d(z, z') \le 1/C$ .

G1 assumption is new in higher dimension which means that we can bound the difference by its difference in each representation  $V_{\alpha}$ , and in the representation  $V_{\alpha}$  the directional derivative  $|\partial_{\alpha}\varphi|$  can bound the Lipschitz norm. G2 and G3 assumptions are natural generalizations of the case m = 1,  $SL_2(\mathbb{R})$ . G4 assumption is used to normalize the function.

The role of J is to simplify the verification of (C, r) goodness. With this definition, we only need to verify assumptions on a neighbourhood of the support of r.

## 4.2 From sum-product estimates to Fourier decay

In this subsection we will prove Theorem 1.7, an estimate of Fourier decay, by using the results established in Section 2 and Section 3.

Recall that we have fixed  $(\epsilon_2, c_1)$  for Proposition 3.4 in Section 3.5, the constant  $(\epsilon_2/2, c')$  in Lemma 2.58 and

$$\kappa_0 = \frac{1}{10} \min\{c_1, c'\}$$

Take  $k, \epsilon_3, \epsilon_4$  from Proposition 3.17 with this  $\kappa_0$ . Let  $\epsilon$  be a positive number to be determined later (the only constant which is not fixed yet). The constant  $\epsilon_0$  in the hypothesis of Theorem 1.7 is defined as

$$\epsilon_0 = \frac{\epsilon}{\max_{\alpha \in \Pi} \{(2k+1)\alpha \sigma_\mu + \epsilon_2\} + \epsilon}$$
(4.6)

which will be fixed once  $\epsilon$  is fixed.

Here, we define and give relations of different constants. Let v be the vector in  $\mathbb{R}^m$  whose components are  $v_{\alpha} = \sup_{\eta \in \text{suppr}} |\partial_{\alpha} \varphi(\eta)|$ , for  $\alpha \in \Pi$ . Then by G4 assumption (4.4), we have

$$\sup_{\alpha \in \Pi} v_{\alpha} \in [\xi^{-\epsilon_0}, \xi^{\epsilon_0}]. \tag{4.7}$$

Let n be the minimal integer such that

$$e^{\epsilon_2 n} \ge \xi \max_{\alpha \in \Pi} \{ v_\alpha e^{-(2k+1)\alpha \sigma_\mu n} \}.$$

$$(4.8)$$

The existence is guaranteed by the positivity of Lyapunov constant, that is  $\alpha \sigma_{\mu} > 0$  for  $\alpha \in \Pi$  (Lemma 2.51). Let the regularity scale  $\delta$  be given by

$$\delta = e^{-\epsilon n} < 1/2,$$

where we take  $\xi$  large enough such that n is large enough. Let the contraction scale  $\beta$  given by

$$\beta_{\alpha} = e^{-\alpha \sigma_{\mu} n}, \beta = \max_{\alpha \in \Pi} \{\beta_{\alpha}\}.$$

The point is that the contraction speed  $\beta$  decides the magnitude of a term and  $\delta$  is only an error term, much larger than  $\beta$ .

Let the frequency  $\tau$  be defined by  $\tau = e^{\epsilon_2 n}$ . By (4.8), we have

$$\tau \ge \xi \max_{\alpha \in \Pi} \{ v_{\alpha} \beta_{\alpha}^{2k+1} \} \ge C_{\epsilon_2} \tau, \tag{4.9}$$

where  $C_{\epsilon_2} = e^{-\epsilon_2} \min_{\alpha \in \Pi} \{ e^{-(2k+1)\alpha \sigma_{\mu}} \}$ . By (4.7), there exists  $\alpha_o$  in  $\Pi$  such that  $v_{\alpha_o} \geq \xi^{-\epsilon_0}$ . Then (4.9) and (4.6) imply that

$$\xi \leq \tau v_{\alpha_o}^{-1} \beta_{\alpha_o}^{-2k-1} \leq \xi^{\epsilon_0} \tau \beta_{\alpha_0}^{-(2k+1)} \leq \xi^{\epsilon_0} e^{n\epsilon \frac{1-\epsilon_0}{\epsilon_0}}.$$

Hence the regularity scale satisfies

$$\xi^{\epsilon_0} \le e^{\epsilon n} = \delta^{-1}. \tag{4.10}$$

Notation: We state some notation which will be used throughout Section 4.2.

- Let  $\mathbf{g} = (g_0, \dots, g_k)$  be an element in  $G^{\times (k+1)}$ .
- Let  $\mathbf{h} = (h_1, \dots, h_k)$  be an element in  $G^{\times k}$ .
- We write  $\mathbf{g} \leftrightarrow \mathbf{h} = g_0 h_1 \cdots h_k g_k \in G$  for the product of  $\mathbf{g}, \mathbf{h}$ .
- We write  $T\mathbf{g} \leftrightarrow \mathbf{h} = g_0 h_1 \cdots g_{k-1} h_k \in G$ .
- For  $l \in \mathbb{N}$ , let  $\mu_{l,n}$  be the product measure on  $G^{\times l}$  given by  $\mu_{l,n} = \mu^{*n} \otimes \cdots \otimes \mu^{*n}$ .
- Recall that for g, h in G and  $\eta$  in  $\mathscr{P}$ , we define  $Y_g^n(h, \eta)_{\alpha} = \exp(-\alpha(\sigma(gh, \eta) \kappa(g) n\sigma_{\mu}))$ and  $Y_q^n(h, \eta) = (Y_q^n(h, \eta)_{\alpha})_{\alpha \in \Pi} \in \mathbb{R}^m$ .
- For z in  $\mathscr{P}_0$ , let  $\tilde{Y}_g^n(h, z)_{\alpha} = \alpha^{\sharp}(m(\ell_g^{-1}, hz))Y_g^n(h, \eta)_{\alpha}$ , where  $\alpha^{\sharp}$  is the corresponding algebraic character of the simple root  $\alpha$  and we make a choice of  $\ell_g$  and  $\eta = \pi(z)$ .
- For g in G, z in  $\mathscr{P}_0$  and  $\eta = \pi(z)$ , let  $\tilde{\lambda}_{g,z}$  be the pushforward measure on  $\mathbb{R}^m$  of  $\mu^{*n}$  restricted to a subset  $G_{n,g,\eta}$  under the map  $\tilde{Y}_g^n(\cdot, z)$ . In other words, for a Borel set E,

$$\tilde{\lambda}_{g,z}(E) = \mu^{*n} \{ h \in G_{n,g,\eta} | \tilde{Y}_q^n(h,z) \in E \}.$$

Recall that the set  $G_{n,g,\eta}$  is defined by  $G_{n,g,\eta} = \{h \in G | h \text{ is } (n,\epsilon,\eta,\zeta_g^m) \text{ good} \}.$ 

• After fixing **g**, we will also fix a choice of  $k_{g_j}$ ,  $\ell_{g_j}$  for  $g_j$  and let  $z_{g_j} = k_{g_j} z_o$ ,  $m_j(h) = m(\ell_{g_{j-1}}^{-1}, hk_{g_j})$  and  $\lambda_j = \tilde{\lambda}_{g_{j-1}, z_{g_j}}$ , for  $j = 1, \ldots, k$ .

**Lemma 4.3.** The measure  $\tilde{\lambda}_{g,z}$  satisfies the same property (3.24) as  $\lambda_{g,\eta}$  with  $C_0$  replaced by  $2^m C_0$ , where  $\eta = \pi(z)$ .

*Proof.* Since the difference is only in the sign, we have

$$(\pi_v)_* \tilde{\lambda}_{g,z}(B_{\mathbb{R}}(a,\rho)) \le \sum_{f \in (\mathbb{Z}/2\mathbb{Z})^m} (\pi_{fv})_* \lambda_{g,\eta}(B_{\mathbb{R}}(a,\rho)),$$

where we identify  $(\mathbb{Z}/2\mathbb{Z})^m$  with  $\{-1,1\}^m \subset \mathbb{R}^m$ . The result follows from this inequality. **First step:** For  $\eta, \eta'$  in  $\mathscr{P}$ , let

$$f(\eta,\eta') = \int_G e^{i\xi(\varphi(g\eta) - \varphi(g\eta'))} r(g\eta) r(g\eta') \mathrm{d}\mu^{*(2k+1)n}(g).$$

$$(4.11)$$

Lemma 4.4. We have

$$\left| \int_{\mathscr{P}} e^{i\xi\varphi(\eta)} r(\eta) \mathrm{d}\nu(\eta) \right|^2 \le \int_{\mathscr{P}^2} f(\eta, \eta') \mathrm{d}\nu(\eta) \mathrm{d}\nu(\eta').$$
(4.12)

*Proof.* By the definition of  $\mu$ -stationary measure and the Cauchy-Schwarz inequality,

$$\begin{split} \left| \int_{\mathscr{P}} e^{i\xi\varphi(\eta)} r(\eta) \mathrm{d}\nu(\eta) \right|^2 \\ &= \left| \int_{\mathscr{P}\times G} e^{i\xi\varphi(g\eta)} r(g\eta) \mathrm{d}\mu^{*(2k+1)n}(g) \mathrm{d}\nu(\eta) \right|^2 \leq \int_G \left| \int_{\mathscr{P}} e^{i\xi\varphi(g\eta)} r(g\eta) \mathrm{d}\nu(\eta) \right|^2 \mathrm{d}\mu^{*(2k+1)n}(g) \\ &= \int_{\mathscr{P}^2} \int_G e^{i\xi(\varphi(g\eta) - \varphi(g\eta'))} r(g\eta) r(g\eta') \mathrm{d}\mu^{*(2k+1)n}(g) \mathrm{d}\nu(\eta) \mathrm{d}\nu(\eta'). \end{split}$$

The proof is complete.

Recall that for  $\eta$  in  $\mathscr{P}$ , we write  $V_{\alpha,\eta}$  for its image in  $\mathbb{P}V_{\alpha}$  and  $d_{\alpha}(\eta, \eta') = d(V_{\alpha,\eta}, V_{\alpha,\eta'})$ .

**Definition 4.5** (Good Position). Let  $\eta, \eta'$  be in  $\mathscr{P}$ . We say that they are in good position if

$$\forall \alpha \in \Pi, \ d_{\alpha}(\eta, \eta') \ge \delta.$$

We fix  $\eta, \eta'$  in good position, which means that  $\eta, \eta'$  are far in all  $\mathbb{P}V_{\alpha}$ . We rewrite the formula.

Lemma 4.6. We have

$$\left| \int_{\mathscr{P}} e^{i\xi\varphi(\eta)} r(\eta) \mathrm{d}\nu(\eta) \right|^2 \leq \int_{\eta,\eta' \ good} f(\eta,\eta') \mathrm{d}\nu(\eta) \mathrm{d}\nu(\eta') + O(\delta^c).$$
(4.13)

*Proof.* By the regularity of stationary measure (2.40), we have

$$\nu\{\eta' \in \mathscr{P} | d_{\alpha}(\eta, \eta') \le \delta\} = \nu\{\eta' \in \mathscr{P} | d(V_{\alpha, \eta}, V_{\alpha, \eta'}) \le \delta\} \le C\delta^c.$$

$$(4.14)$$

Therefore by (4.14) and Fubini's theorem,

$$\nu \otimes \nu\{(\eta, \eta') \in \mathscr{P}^2 | d_{\alpha}(\eta', \eta) < \delta\} = \int_{\eta \in \mathscr{P}} \nu\{\eta' \in \mathscr{P} | d_{\alpha}(\eta, \eta') \le \delta\} \mathrm{d}\nu(\eta) \ll \delta^c.$$

Summing over simple roots  $\alpha$ , we obtain the result by  $||r||_{\infty} \leq 1$ .

**Second step:** The purpose of this part is to give a Ping-Pong Lemma in measure sense. We will eliminate sets with negligible measure such that the Ping-Pong condition is almost preserved by iteration on the complement. We fix  $g_j$  for  $j = 0, \ldots, k - 1$  which satisfies

$$\|\kappa(g_j) - n\sigma_{\mu}\| \le \epsilon n/C_A. \tag{4.15}$$

Recall that  $C_A$  is a constant in Definition 2.57. We also demand that

$$h_{j+1}$$
 is  $(n,\epsilon,\eta_{g_{j+1}}^M,\zeta_{g_j}^m)$  good. (4.16)

Recall that the Cartan subspace  $\mathfrak{a}$  is equipped with the norm induced by the Killing form, and with this norm  $\mathfrak{a}$  is isomorphic to the euclidean space  $\mathbb{R}^m$ .

**Lemma 4.7.** Suppose that  $\mathbf{g}$ ,  $\mathbf{h}$  satisfy the above conditions (4.15) and (4.16). Then the action of  $T\mathbf{g} \leftrightarrow \mathbf{h}$  on  $b^M_{V_{\alpha},g_k}(\delta)$  is  $\beta^{2k}_{\alpha} \delta^{-O(1)}$  Lipschitz and

$$e^{-\alpha\sigma(g_0h_1,x_{g_1}^M)}\cdots e^{-\alpha\sigma(g_{k-1}h_k,x_{g_k}^M)} \le \beta_{\alpha}^{2k}\delta^{-O(1)},$$
(4.17)

for every  $\alpha$  in  $\Pi$ . For  $t \in b_{g_k}^M(\delta)$ , let  $t_j = g_j h_{j+1} \cdots h_k t$  for  $j = 0, \ldots, k$ , where we let  $t_k = t$ . Then

$$t_j \in b_{g_j}^M(\beta\delta^{-2}) \subset b_{g_j}^M(\delta), \tag{4.18}$$

$$\|\sigma(g_j h_{j+1}, t_{j+1}) - \sigma(g_j h_{j+1}, \eta_{g_{j+1}}^M)\| \ll \beta \delta^{-O(1)}.$$
(4.19)

**Remark 4.8.** The contraction constant  $\beta$  here is a little different from the gap  $\gamma(g_j)$ , but  $\gamma(g_j)/\beta$ is in the interval  $[\delta^{O(1)}, \delta^{-O(1)}]$  by Lemma 2.59. Hence they are of the same largeness and we will not distinguish them.

The intuition here is that by controlling  $\kappa(g), \eta_g^M, \zeta_g^m$ , all the other position or length will also be controlled, which is similar to hyperbolic dynamics.

*Proof.* For every  $\alpha$  in  $\Pi$ , using Lemma 2.11 2k times, we obtain the Lipschitz property. By Lemma 2.59, we have (4.17) from (4.16) for all  $\alpha$  in  $\Pi$  at the same time.

We use induction to prove the inclusion. For j = k, it is due to the hypothesis of Lemma 4.7.

Suppose that the property holds for j + 1. By definition,  $t_j = g_j h_{j+1} t_{j+1}$ . We abbreviate  $g_j, h_{j+1}, t_{j+1}, \eta_{g_{j+1}}^M$  to  $g, h, \eta, \eta'$ . The condition becomes

$$d(\eta, \eta') \leq \delta$$
,  $\|\kappa(g) - n\sigma_{\mu}\| \leq \epsilon n/C_A$  and h is  $(n, \epsilon, \eta', \zeta_a^m)$  good.

By Lemma 2.59, we have  $\gamma(h) \leq \beta \delta^{-1}$ . By Lemma 2.16, due to  $\eta \in B(\eta', \delta) \subset B_h^m(\delta)$ , we have  $h\eta \in b_h^M(\beta/\delta^2) \subset B_g^m(\delta)$ . Therefore  $gh\eta \in b_g^M(\beta/\delta^2)$ , which is the inclusion condition. Then we will prove (4.19) and we keep the notation  $g, h, \eta, \eta'$ .

$$\|\sigma(gh,\eta) - \sigma(gh,\eta')\| \ll \|\sigma(g,h\eta) - \sigma(g,h\eta')\| + \|\sigma(h,\eta) - \sigma(h,\eta')\|.$$

By the same argument, due to Lemma 2.16 and  $\eta, \eta' \in B(\eta', \beta/\delta^2) \subset B_h^m(\delta)$ , we have  $h\eta, h\eta' \in B(\eta', \beta/\delta^2) \subset B_h^m(\delta)$ .  $b_h^{\check{M}}(\beta/\delta^2) \subset B_g^m(\delta)$ . Therefore by the Lipschitz property of Lemma 2.16

$$\|\sigma(gh,\eta) - \sigma(gh,\eta')\| \ll (d(\eta,\eta') + d(h\eta,h\eta'))\delta \ll \beta/\delta^3.$$

The proof is complete.

**Lemma 4.9.** Suppose that  $\mathbf{g}, \mathbf{h}$  satisfy the conditions (4.15) and (4.16). Let s be in  $\{z \in$  $\mathscr{P}_0|d_0(z, z_{g_k}) \leq \delta\}$ . Let  $s_j = g_j h_{j+1} \cdots h_k s$  for  $j = 0, \dots, k$ , where we let  $s_k = s$ . We have

$$\mathbf{m}(s_0, k_{g_0}) = \Pi_{1 \le j \le k} \mathbf{m}(\ell_{g_{j-1}}^{-1}, h_j k_{g_j}) = \Pi_{1 \le j \le k} m_j(h_j).$$
(4.20)

*Proof.* We let  $\eta = \pi(s)$ , then  $\eta$  is in  $b_{g_k}^M(\delta)$ . By (4.18) with j = 1 and (4.16) with j = 0, Lemma 2.30 implies

$$m(s_0, k_{g_0}) = m(k_{g_0}, g_0 h_1 s_1) = m(\ell_{g_0}^{-1}, h_1 k_{g_1}) m(s_1, k_{g_1}).$$

Iterating this formula, we obtain the result.

Third step: Here we mimic the proof of [BD17], where they heavily use the properties of Schottky groups and symbolic dynamics. But in our case, the group is much more complicate from the point of view of dynamics. We use the large deviation principle to get a similar formula.

By very careful control of  $g_l$ , with a loss of an exponentially small measure, we are able to rewrite the formula in a form to use the sum-product estimates. The key point is that by controlling the Cartan projection and the position of  $\eta_g^M$  and  $\zeta_g^m$  of each  $g_l$ , we are able to get good control of their product  $\mathbf{g} \leftrightarrow \mathbf{h}$ .

We should notice that the element  $g_j$  will be fixed, and we will integrate first with respect to  $h_j$ . This gives the independence of the cocycle  $\sigma(g_{j-1}h_j, \eta_{g_j}^M)$ , that is for different j they are independent, which is an important point to apply sum-product estimates.

We return to (4.13). We call **g** "good" with respect to  $\eta, \eta'$  if

**g** satisfies (4.15), 
$$g_k$$
 satisfies conditions in Lemma 2.45,  $\eta_{g_0}^M \in \text{supp}r$   
and  $\delta(\eta, \zeta_{g_k}^m), \delta(\eta', \zeta_{g_k}^m), \delta(V_{\alpha,\eta} \wedge V_{\alpha,\eta'}, y_{\wedge^2 \rho_{\alpha} g_k}^m) \ge 4\delta.$  (4.21)

**Lemma 4.10.** If  $\eta$  and  $\eta'$  are in good position and  $\mathbf{g}$  is "good", then  $g_k\eta, g_k\eta'$  are in  $b_{g_k}^M(\delta)$ , and for  $\alpha \in \Pi$  the  $d_\alpha$  distance between  $g_k\eta$  and  $g_k\eta'$  is almost  $\beta_\alpha$ , that is

$$d_{\alpha}(g_k\eta, g_k\eta') \in \beta_{\alpha}[\delta^{O(1)}, \delta^{-O(1)}].$$

*Proof.* The inclusion is due to Lemma 2.16. Since g is good (4.21), by (2.13) we have the lower bound and by the Lipschitz property in Lemma 2.11 we have the upper bound.  $\Box$ 

For  $\eta, \eta'$  in  $\mathscr{P}$ , we can rewrite the formula of  $f(\eta, \eta')$  as

$$f(\eta, \eta') = \int e^{i\xi(\varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta'))} r(\mathbf{g} \leftrightarrow \mathbf{h}\eta) r(\mathbf{g} \leftrightarrow \mathbf{h}\eta') d\mu_{k,n}(\mathbf{h}) d\mu_{k+1,n}(\mathbf{g}).$$
(4.22)

We call  $\mathbf{h}$  is g-regular if  $\mathbf{h}$  satisfies (4.16). Let

$$f_{\mathbf{g}}(\eta, \eta') = \int_{\mathbf{g}-regular} e^{i\xi(\varphi(\mathbf{g}\leftrightarrow\mathbf{h}\eta) - \varphi(\mathbf{g}\leftrightarrow\mathbf{h}\eta'))} \mathrm{d}\mu_{k,n}(\mathbf{h})$$

Lemma 4.11. For  $\eta, \eta'$  in  $\mathscr{P}$ 

$$|f(\eta, \eta')| \le \int_{\mathbf{g}^{"good"}} |f_{\mathbf{g}}(\eta, \eta')| \mathrm{d}\mu_{k+1, n}(\mathbf{g}) + O_{\epsilon}(\delta^{c}), \qquad (4.23)$$

if  $\epsilon$  is small enough with respect to  $\gamma$ , that is  $\epsilon \leq \min_{\alpha \in \Pi} \{ \alpha \sigma_{\mu} \gamma / (2 + 2\gamma) \}.$ 

*Proof.* Let

$$\tilde{f}_{\mathbf{g}}(\eta,\eta') = \int_{\mathbf{g}-regular} e^{i\xi(\varphi(\mathbf{g}\leftrightarrow\mathbf{h}\eta)-\varphi(\mathbf{g}\leftrightarrow\mathbf{h}\eta'))} r(\mathbf{g}\leftrightarrow\mathbf{h}\eta) r(\mathbf{g}\leftrightarrow\mathbf{h}\eta') \mathrm{d}\mu_{k,n}(\mathbf{h}).$$

We call **g** "semi-good" if **g** satisfies (4.21) except the assumption of  $\eta_{g_0}^M \in \text{supp}r$  in (4.21). By large deviation principle (Proposition 2.52, Proposition 2.54, Lemma 2.65), we conclude that

$$\mu_{k+1,n}\{g \text{ not "semi-good"}\} \le O_{\epsilon}(\delta^c). \tag{4.24}$$

Then by (4.22), Lemma 2.58 and (4.24),

$$|f(\eta,\eta')| \leq \int_{\mathbf{g}} |\tilde{f}_{\mathbf{g}}(\eta,\eta')| \mathrm{d}\mu_{k+1,n}(\mathbf{g}) + O_{\epsilon}(\delta^{c}) \leq \int_{\mathbf{g}^{"semi-good"}} |\tilde{f}_{\mathbf{g}}(\eta,\eta')| \mathrm{d}\mu_{k+1,n}(\mathbf{g}) + O_{\epsilon}(\delta^{c}).$$

$$(4.25)$$

By Lemma 4.10, (4.18) with j = 0 and  $c_{\gamma}(r) \leq \xi^{\epsilon_0} \leq \delta^{-1}$ ,

$$|r(\eta_{g_0}^M)^2 - r(\mathbf{g} \leftrightarrow \mathbf{h}\eta)r(\mathbf{g} \leftrightarrow \mathbf{h}\eta')| \le 2||r||_{\infty}c_{\gamma}(r)(\beta\delta^{-2})^{\gamma} \le 2\beta^{\gamma}\delta^{-1-2\gamma} \le 2\delta,$$

if  $\epsilon$  is small enough with respect to  $\gamma$ . Hence

$$\begin{aligned} |\tilde{f}_{\mathbf{g}}(\eta,\eta')| &\leq \left| \int_{\mathbf{g}-regular} e^{i\xi(\varphi(\mathbf{g}\leftrightarrow\mathbf{h}\eta)-\varphi(\mathbf{g}\leftrightarrow\mathbf{h}\eta'))} r(\eta_{g_0}^M)^2 \mathrm{d}\mu_{k,n}(\mathbf{h}) \right| + O(\delta^c) \\ &\leq r(\eta_{g_0}^M)^2 |f_{\mathbf{g}}(\eta,\eta')| + O(\delta^c). \end{aligned}$$
(4.26)

If  $r(\eta_{g_0}^M) \neq 0$ , then that **g** is "semi-good" implies **g** is "good". Combined with (4.25) and (4.26), by  $||r||_{\infty} \leq 1$ , we have

$$\begin{aligned} |f(\eta, \eta')| &\leq \int_{\mathbf{g}^{"}semi-good"} \left( r(\eta_{g_0}^M)^2 |f_{\mathbf{g}}(\eta, \eta')| + O(\delta^c) \right) \mathrm{d}\mu_{k+1,n}(\mathbf{g}) + O_{\epsilon}(\delta^c) \\ &\leq \int_{\mathbf{g}^{"}good"} |f_{\mathbf{g}}(\eta, \eta')| \mathrm{d}\mu_{k+1,n}(\mathbf{g}) + O_{\epsilon}(\delta^c). \end{aligned}$$

The proof is complete.

Recall that  $\beta$  is the magnitude which is really small,  $\delta$  is only an error term and  $\tau$  is the frequency for applying the sum-product estimate, which lies between  $\delta^{-1}$  and  $\beta^{-1}$ .

**Proposition 4.12.** Let  $I_{\tau} = [\tau^{3/4}, \tau^{5/4}]$ . The following formula is true for  $\eta, \eta'$  in good position and **g** "good",

$$|f_{\mathbf{g}}(\eta,\eta')| \le \sup_{\|\varsigma\|\in I_{\tau}} \left| \int e^{i\langle\varsigma,x_1\cdots x_k\rangle} \mathrm{d}\lambda_1(x_1)\cdots\lambda_k(x_k) \right| + O(\beta\delta^{-O(1)}\tau), \tag{4.27}$$

when  $\epsilon$  is small enough with respect to  $\epsilon_2$ .

**Remark 4.13.** This is the most complicate step, where the difficulty comes from higher rank. We need to use the technique of changing flags to find the direction of slowest contraction speed, where we can use Newton-Leibniz's formula. Since the action of the sign group M is non trivial on the slowest directions, we also carefully treat the sign.

*Proof.* The element  $\eta, \eta'$  and **g** are already fixed. Since  $g_k$  satisfies the conditions in Lemma 2.45, we obtain two chains  $(\eta = \eta_o, \eta_1, \ldots, \eta_{l_1})$  and  $(\eta' = \eta'_o, \eta'_1, \ldots, \eta'_{l_2})$  as in Lemma 2.45. Then we write

$$\varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta') = \sum_{0 \le j \le l_1 - 1} (\varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta_j) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta_{j+1})) - \sum_{0 \le j \le l_2 - 1} (\varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta'_j) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta'_{j+1})) + (\varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta_{l_1}) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta'_{l_2})),$$

$$(4.28)$$

The terms for different j and for  $\eta, \eta'$  are similar. We fix j and we simplify  $\alpha(\eta_i, \eta_{i+1})$  to  $\alpha$ .

We compute the term  $\varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta_j) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta_{j+1})$ . In order to treat the sign, we will work on  $\mathscr{P}_0 = G/A_e N$ . Recall that  $\pi : \mathscr{P}_0 \to \mathscr{P}$  is the projection and we use  $z = kz_o$  to denote the element  $kA_e N$  in  $\mathscr{P}_0$ .

By Lemma 2.47 and (4.21), we know that  $g_k\eta_j, g_k\eta_{j+1}$  are in  $b_{g_k}^M(\delta)$ , which satisfy the condition of Lemma 4.7. Let  $z_0, z_1$  be preimages of  $g_k\eta_j$  and  $g_k\eta_{j+1}$  in  $\mathscr{P}_0$  such that  $m(z_0, z_1) = e$ . Notice that  $z_0, z_1$  are in the same  $\alpha$ -circle. By Lemma 2.45 (2.30) and Lemma 4.10

$$d(g_k\eta_j, g_k\eta_{j+1}) = d_\alpha(g_k\eta, g_k\eta') + O(\beta e^{-\alpha\kappa(g_k)}\delta^{O(1)}) \in \beta_\alpha[\delta^{O(1)}, \delta^{-O(1)}].$$

Due to  $m(z_0, z_1) = e$ , the arc length distance also satisfies

$$d_A(z_0, z_1) = \arcsin d(g_k \eta_j, g_k \eta_{j+1}) \in \beta_\alpha[\delta^{O(1)}, \delta^{-O(1)}].$$
(4.29)

Now, we lift  $\varphi$  to  $\mathscr{P}_0$ , becoming a right *M*-invariant function. By abuse of notation, we also use  $\varphi$  to denote the lifted function. Let  $\gamma$  be an arc connecting  $z_0, z_1$  with unit speed in the  $\alpha$ -circle with length less than  $\pi/2$ . Without loss of generality, we suppose that  $\gamma$  is in the positive direction (If not, we add minus in the right hand side of (4.30)). By Newton-Leibniz's formula (2.28), we have

$$\varphi(T\mathbf{g} \leftrightarrow \mathbf{h}z_0) - \varphi(T\mathbf{g} \leftrightarrow \mathbf{h}z_1) = \int_0^u \partial_\alpha \varphi(T\mathbf{g} \leftrightarrow \mathbf{h}\gamma(t)) e^{-\alpha\sigma(T\mathbf{g} \leftrightarrow \mathbf{h},\gamma(t))} \mathrm{d}t, \qquad (4.30)$$

where  $u = d_A(z_0, z_1)$ . Fix a time t in [0, u], let  $s_j = g_j h_{j+1} \cdots h_k \gamma(t)$ . Then  $\pi(\gamma(t))$  is in  $b_{g_k}^M(\delta)$ , because  $g_k \eta_j$  and  $g_k \eta_{j+1}$  are in  $b_{g_k}^M(\delta)$  and by (4.29). By (4.18), the element  $\pi(s_0)$ , the image of  $s_0 = T \mathbf{g} \leftrightarrow \mathbf{h} \gamma(t)$  in  $\mathscr{P}$ , is in  $b_{g_0}^M(\beta \delta^{-O(1)})$ .

Recall that we have made a choice of the Cartan decomposition of every  $g_j$  for  $0 \le j \le k$ . In particular,  $k_{g_0}$  is given in the decomposition of  $g_0 = k_{g_0} a_{g_0} \ell_{g_0} \in KA^+K$ . Let  $m_0 = m(s_0, k_{g_0})$  and  $\underline{s_0} = s_0 m_0$ , then  $m(\underline{s_0}, k_{g_0}) = e$ . By Lemma 2.36,

$$\partial_{\alpha}\varphi_{s_0} = \partial_{\alpha}\varphi_{\underline{s_0}m_0} = \alpha^{\sharp}(m_0)\partial_{\alpha}\varphi_{\underline{s_0}}.$$
(4.31)

By Lemma 5.8 and  $\pi s_0, \pi z_{g_0} = \eta_{g_0}^M$  in  $b_{g_0}^M(\beta \delta^{-O(1)})$ , we have

$$d_0(\underline{s_0}, z_{g_0}) \le d(\pi s_0, \pi z_{g_0}) < \beta \delta^{-O(1)}.$$
(4.32)

Due to **g** good (4.21), we have  $\eta_{g_0}^M \in \text{suppr.}$  By G2 assumption (4.2), we have  $|\partial_{\alpha}\varphi(z_{g_0})| \geq \delta v_{\alpha}$ . By (4.32), the point  $\pi s_0$  is in J, the  $\delta$  neighbourhood of suppr. By G3 assumption (4.3),  $|\partial_{\alpha}\varphi(\underline{s_0}) - \partial_{\alpha}\varphi(z_{g_0})| \leq \delta^{-1}v_{\alpha}d_0(\underline{s_0}, z_{g_0})$ , which implies

$$\partial_{\alpha}\varphi(\underline{s_0})/\partial_{\alpha}\varphi(z_{g_0}) \in [1 - \beta\delta^{-O(1)}, 1 + \beta\delta^{-O(1)}].$$

By Lemma 4.7 (4.19), we have

$$(1 - \beta \delta^{-O(1)})e^{-O(\beta/\delta)} \le \frac{\partial_{\alpha}\varphi(\underline{s_0})e^{-\alpha\sigma(g_0h_1,s_1)}\cdots e^{-\alpha\sigma(g_{k-1}h_k,s_k)}}{\partial_{\alpha}\varphi(z_{g_0})e^{-\alpha\sigma(g_0h_1,x_{g_1}^M)}\cdots e^{-\alpha\sigma(g_{k-1}h_k,x_{g_k}^M)}} \le (1 + \beta \delta^{-O(1)})e^{O(\beta/\delta)}.$$
(4.33)

By (4.17),

$$B_{\alpha} := e^{-\alpha\sigma(g_0h_1, x_{g_1}^M)} \cdots e^{-\alpha\sigma(g_{k-1}h_k, x_{g_k}^M)} \le \beta_{\alpha}^{2k} \delta^{-O(1)}$$

Together with (4.29)-(4.33)

$$|\varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta_j) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta_{j+1}) - d_A(z_0, z_1)\alpha^{\sharp}(m_0)\partial_{\alpha}\varphi(z_{g_0})B_{\alpha}| \le \beta\beta_{\alpha}^{2k+1}\delta^{-O(1)}v_{\alpha}.$$
(4.34)

We deal with the error term which comes from the process of changing flags. The Lipschitz property in Lemma 4.7 and Lemma 2.45 (2.31) imply that

$$d_{\alpha}(\mathbf{g} \leftrightarrow \mathbf{h}\eta_{l_1}, \mathbf{g} \leftrightarrow \mathbf{h}\eta'_{l_2}) \leq \beta_{\alpha}^{2k} \delta^{-O(1)} d_{\alpha}(g_k \eta_{l_1}, g_k \eta'_{l_2}) \leq \beta_{\alpha}^{2k+1} \beta \delta^{-O(1)}$$

Due to (4.18) in Lemma 4.7 and Lemma 2.47, the two points  $\mathbf{g} \leftrightarrow \mathbf{h}\eta_{l_1}, \mathbf{g} \leftrightarrow \mathbf{h}\eta'_{l_2}$  are in J, the  $\delta$  neighbourhood of suppr. Due to G1 assumption (4.1)

$$|\varphi(\mathbf{g} \leftrightarrow \mathbf{h} \eta_{l_1}) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h} \eta'_{l_2})| \leq \delta^{-1} \sum_{\alpha} v_{\alpha} d_{\alpha}(\mathbf{g} \leftrightarrow \mathbf{h} \eta_{l_1}, \mathbf{g} \leftrightarrow \mathbf{h} \eta'_{l_2}).$$

Therefore

$$|\varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta_{l_1}) - \varphi(\mathbf{g} \leftrightarrow \mathbf{h}\eta'_{l_2})| \le \delta^{-O(1)}\beta \sum_{\alpha} v_{\alpha}\beta_{\alpha}^{2k+1}.$$
(4.35)

We collect information for different simple roots. Recall that for a fixed g in G and for  $h \in G$ ,  $z \in \mathscr{P}_0$ , we have defined  $\tilde{Y}_g^n(h, z)_{\alpha} = e^{-\alpha(\sigma(gh, z) - \kappa(g) - n\sigma_{\mu})}\alpha(\mathrm{m}(\ell_g, hk))$ . Let

$$\varsigma_{\alpha} := \frac{\xi d_A(z_0, z_1) \alpha^{\sharp}(m_0) \partial_{\alpha} \varphi(z_{g_0}) B_{\alpha}}{\prod_{l=1}^k \tilde{Y}_{g_{l-1}}^n(h_l, z_{g_l})_{\alpha}}$$

Let  $\varsigma = (\varsigma_{\alpha})_{\alpha \in \Pi} \in \mathbb{R}^{m}$ . Hence by (4.28), (4.34), (4.35) and (4.9)

$$|\xi(\varphi(\mathbf{g}\leftrightarrow\mathbf{h}x)-\varphi(\mathbf{g}\leftrightarrow\mathbf{h}x'))-\langle\varsigma,\Pi_{l=1}^{k}\tilde{Y}_{g_{l-1}}^{n}(h_{l},z_{g_{l}})\rangle| \leq \beta\delta^{-O(1)}\sum_{\alpha}\beta_{\alpha}^{2k+1}v_{\alpha}\xi \ll \beta\delta^{-O(1)}\tau.$$
(4.36)

We want to verify that  $\|\varsigma\| \in I_{\tau}$ . By (4.20), we have

$$\varsigma_{\alpha} = \xi d_A(z_0, z_1) \partial_{\alpha} \varphi(z_{g_0}) \beta_{\alpha}^k e^{-\alpha \kappa(g_0) - \dots - \alpha \kappa(g_{k-1})}.$$

By (4.15), (4.29), (4.21) and (4.2) we have  $|\varsigma_{\alpha}| \in \xi v_{\alpha} \beta_{\alpha}^{2k+1}[\delta^{O(1)}, \delta^{-O(1)}]$ . Therefore by (4.9),

$$\|\varsigma\| \in \sup_{\alpha} \xi v_{\alpha} \beta_{\alpha}^{2k+1}[\delta^{O(1)}, \delta^{-O(1)}] \in \tau[\delta^{O(1)}, \delta^{-O(1)}] \subset [\tau^{3/4}, \tau^{5/4}] = I_{\tau}.$$

By definition, the distribution of  $\tilde{Y}_{g_{l-1}}^n(h_l, z_{g_l})$ , where  $h_l$  satisfies (4.16) with distribution  $\mu^{*n}$ , is the measure  $\lambda_l$ . Finally, due to  $|e^{ix} - e^{iy}| \leq |x - y|$  for  $x, y \in \mathbb{R}$ , the inequality (4.36) implies (4.27).

Fourth step: We are able to apply sum-product estimates.

Proof of Theorem 1.7. For l = 1, 2, ..., k, Proposition 3.18 and Lemma 4.3 tell us that with  $\epsilon$  small enough depending on  $\epsilon_4 \epsilon$ , there exists  $C_0$  such that the measures  $\lambda_l$  satisfy the assumptions in Proposition 3.17 with  $\tau$ .

Proposition 3.17 implies that for  $\tau$  large enough,

$$\left|\int \exp(i\langle\varsigma, x_1\cdots x_k\rangle) \mathrm{d}\lambda_1(x_1)\ldots \mathrm{d}\lambda_k(x_k)\right| \leq \tau^{-\epsilon_3}.$$

Then by (4.13), (4.23) and (4.27), we have

$$\left|\int e^{i\xi\varphi(\eta)}r(\eta)\mathrm{d}\nu(\eta)\right|^2 \le O_\epsilon(\delta^c) + O(\beta\delta^{-O(1)}\tau) + \tau^{-\epsilon_3}.$$

Due to  $\beta \delta^{-O(1)} \tau = \max_{\alpha \in \Pi} e^{(-\alpha \sigma_{\mu} + O(1)\epsilon + \epsilon_2)n}$ , take  $\epsilon$  small enough. The proof is complete.  $\Box$ 

**Remark 4.14.** Another difference with [BD17] is that we avoid using the renewal idea, which simplifies the proof of this part. The renewal idea is that instead of using  $\mu^{*n}$ , we use a renewal measure  $\mu_t$ , which is defined to be the distribution of  $g_1 \cdots g_n$  for the first time that its Cartan projection exceeds t, where  $g_1, g_2 \ldots$  are i.i.d. random variables with distribution  $\mu$ . This is because we generalize the sum-product estimate to a form that the measure can have a support which depends on the frequency, and we use the large deviation principle to prove that our measure has a support not too large with respect to the frequency.

### 4.3 From Fourier decay to spectral gap

From now on, we will only consider semisimple groups. In this section, we will prove Theorem 1.3 and Theorem 1.4 by using Theorem 5.4.

#### Derivative of the cocycle

This part is devoted to the derivative of the cocycle. The results of this part imply that for most g, h in G, the difference of the Iwasawa cocycle  $\sigma(g, \cdot) - \sigma(h, \cdot)$  satisfies the (C, r) good condition in Definition 4.1 (See Lemma 4.25). Since the  $\alpha$ -bundle is trivial on  $\mathcal{P}_0$ , we will work on  $\mathcal{P}_0$ . We need to lift the Iwasawa cocycle  $\sigma$  to  $\mathcal{P}_0$  and we use the same notation  $\sigma$ .

Let V be an irreducible representation of G with a good norm. Recall that  $\sigma_V(g, x) = \frac{\|\rho(g)v\|}{\|v\|}$  for g in G and v in V. We will abbreviate  $\rho g$  to g in the proof, because  $(\rho, V)$  is the only representation to be studied in this part. Let  $\alpha$  be a simple root. Let  $e_1$  be a unit vector of highest weight in V and let  $e_2 = Y_{\alpha}e_1$ .

**Lemma 4.15.** Let V be an irreducible representation of G with a good norm. For  $z = kz_o$  in  $\mathscr{P}_0$ , we have

$$\partial_{\alpha}\sigma_V(g,z) = \frac{\langle \rho g v, \rho g u \rangle}{\|\rho g v\|^2},$$

where  $v = ke_1$  and  $u = ke_2$ .

*Proof.* Without loss of generality, we suppose that  $z = z_o$ . Since  $Y_{\alpha}$  is a left K invariant vector field on  $\mathscr{P}_0$ , we have

$$\begin{aligned} \partial_{Y_{\alpha}}\sigma_{V}(g,e) &= \partial_{t}\sigma_{V}(g,\exp(tY_{\alpha})z_{o})|_{t=0} = \partial_{t}\left(\log\frac{\|g\exp(tY_{\alpha})e_{1}\|}{\|\exp(tY_{\alpha})e_{1}\|}\right)\Big|_{t=0} \\ &= \frac{\langle ge_{1},gY_{\alpha}e_{1}\rangle}{\|ge_{1}\|^{2}} - \frac{\langle e_{1},Y_{\alpha}e_{1}\rangle}{\|e_{1}\|^{2}}. \end{aligned}$$

Since the norm is good, eigenvectors of different weights are orthogonal, we have  $\langle e_1, Y_{\alpha} e_1 \rangle = 0$ . The result follows.

Form this lemma, we know that the derivative of the cocycle  $\sigma_V$  in the direction  $Y_{\alpha}$  is nonzero only if  $\chi - \alpha$  is a weight of V. We fix the distance  $d_0$  on  $\mathscr{P}_0$ , which is defined in Appendix 5.2.

**Lemma 4.16.** Let  $\delta < 1/2$ . Let  $\widetilde{B^m_{V,g}(\delta)}$  be the preimage of  $B^m_{V,g}(\delta) \subset \mathbb{P}V$  in  $\mathscr{P}_0$ . For  $z = kz_o \in \widetilde{B^m_{V,g}(\delta)}$ ,

$$|\partial_{\alpha}\sigma_V(g,z)| \le \delta^{-O(1)}.\tag{4.37}$$

We also have

$$Lip_{\mathscr{P}_0}(\partial_\alpha \sigma_V(g,\cdot)|_{\widetilde{B^m_{V,g}(\delta)}}) \le \delta^{-O(1)}.$$
(4.38)

*Proof.* By Lemma 4.15, the hypothesis that  $\mathbb{R}ke_1 \in B^m_{V,q}(\delta)$  and (2.12)

$$|\partial_{\alpha}\sigma_{V}(g,z)| = \left|\frac{\langle gke_{1}, gke_{2}\rangle}{\|gke_{1}\|^{2}}\right| \le \frac{\|Y_{\alpha}\|\|g\|^{2}\|e_{1}\|^{2}}{\|g\|^{2}\delta^{2}\|e_{1}\|^{2}}.$$

Since the operator norm of  $Y_{\alpha}$  is bounded, we have

$$|\partial_{\alpha}\sigma_V(g,z)| \le \delta^{-O(1)}$$

The estimate of Lipschitz norm is more complicate. Let  $v = ke_1, v' = k'e_1, u = ke_2, u' = k'e_2$ . We have

$$|\partial_{\alpha}\sigma_V(g,z) - \partial_{\alpha}\sigma_V(g,z')| = \frac{|\langle gv, gu \rangle ||gv'||^2 - \langle gv', gu' \rangle ||gv||^2|}{||gv||^2 ||gv'||^2}$$

By the same argument, due to  $v = ke_1 \in B^m_{V,g}(\delta)$ , we use (2.12) to give a lower bound of the denominator, that is

$$||gv||^2 ||gv'||^2 \ge \delta^4 ||g||^4 ||v||^2 ||v'||^2 = \delta^4 ||g||^4 ||e_1||^4.$$

Use the difference to give a upper bound of the numerator, that is

$$\begin{aligned} |\langle gv, gu \rangle ||gv'||^2 - \langle gv', gu' \rangle ||gv||^2 | \\ \ll ||g||^3 ||e_1||^3 (||gv - gv'|| + ||gu - gu'||) \ll ||g||^4 ||v||^3 (||v - v'|| + ||u - u'||). \end{aligned}$$

Therefore we have

$$|\partial_{\alpha}\sigma_{V}(g,z) - \partial_{\alpha}\sigma_{V}(g,z')| \ll \delta^{-O(1)}(\|ke_{1} - k'e_{1}\| + \|ke_{2} - k'e_{2}\|).$$

Then by Lemma 5.6, the proof is complete.

Let V be a finite dimensional vector space with euclidean norm. Recall that  $\wedge^2 Sym^2 V$  is the exterior square of the symmetric square of V. It is a linear space generated by vectors of the form  $v_1v_2 \wedge v_3v_4$  where  $v_i$  are in V, for i = 1, 2, 3, 4. For g, h in GL(V), let  $F_{g,h}$  be the linear functional on  $\wedge^2 Sym^2 V$ , whose action on the vector  $v_1v_2 \wedge w_1w_2$  is defined by

$$F_{g,h}(v_1v_2 \wedge w_1w_2) = \langle hv_1, hv_2 \rangle \langle gw_1, gw_2 \rangle - \langle gv_1, gv_2 \rangle \langle hw_1, hw_2 \rangle.$$

This formula is well defined because  $v_1, v_2$  and  $w_1, w_2$  are symmetric, respectively. We also have  $F_{g,h}(v_1v_2 \wedge w_1w_2) = -F_{g,h}(w_1w_2 \wedge v_1v_2)$ . Since the vectors of form  $v_1v_2 \wedge w_1w_2$  generate the space  $\wedge^2 Sym^2V$ , the linear form  $F_{g,h}$  is uniquely defined.

Suppose that V is a super proximal representation of G with highest weight  $\chi$  (Definition 2.5). Let  $\alpha$  be the unique simple root such that  $\chi - \alpha$  is a weight of V. The space  $\wedge^2 Sym^2V$  may be reducible. The two highest weights of  $Sym^2V$  are  $2\chi$ ,  $2\chi - \alpha$ , whose eigenspaces have dimension 1. Hence, the highest weight of  $\wedge^2 Sym^2V$  is  $4\chi - \alpha$ , and the eigenspace has dimension 1. Let W be the irreducible subrepresentation of  $\wedge^2 Sym^2V$  with the highest weight  $\chi_1 := 4\chi - \alpha$ . In the following lemma, we abbreviate  $\rho(g), \rho(h)$  to g, h.

**Lemma 4.17.** Let  $\delta < 1/2$ . Let V be a super proximal representation of G and let  $\alpha$  be the unique simple root such that  $\chi - \alpha$  is a weight of V. Recall that  $V_{\chi_1,\eta}$  is the image of  $\eta \in \mathscr{P}$  in  $\mathbb{P}W$ . If g, h in G and  $z = kz_o \in \mathscr{P}_0, \eta = \pi(z)$  satisfy

(1) 
$$\ell_h^{-1} V^{\chi}, \ell_h^{-1} V^{\chi-\alpha} \in B^m_{V,g}(\delta), \gamma_{1,2}(g) \le \delta^3,$$
  
(2)  $\delta(V_{\chi_1,\eta}, F_{g,h}|_W) > \delta$  and  $V_{\chi,\eta} \in B^m_{V,g}(\delta) \cap B^m_{V,h}(\delta),$ 

then

$$\left|\partial_{\alpha}(\sigma_V(g,z) - \sigma_V(h,z))\right| \ge \delta^{O(1)}.$$

**Remark 4.18.** This is similar to the non local integrability property as defined in [Dol98] [Nau05] and [Sto11]. Although the above two conditions are complicate, we will see later that in the measure sense, most pairs g, h satisfy these conditions.

The key idea here is to use other representation to linearise polynomial functions on V. As long as the function is linear, we will have a good control of it. Another point is that the image of  $\mathscr{P}$  stays in the same irreducible subrepresentation.

Proof of Lemma 4.17. By Lemma 4.15, let

$$L := \partial_{\alpha}(\sigma_V(g, z) - \sigma_V(h, z)) = \frac{F_{g,h}(v^2 \wedge vu)}{\|gv\|^2 \|hv\|^2},$$
(4.39)

where  $v = ke_1$  and  $u = kY_{\alpha}e_1$  as in Lemma 4.15.

**Lemma 4.19.** If g, h satisfy assumption (1), then the operator norm satisfy

$$||F_{g,h}|_W|| \ge \delta^{O(1)} ||g||^2 ||h||^2$$

*Proof.* Using the Cartan decomposition and good norm, we can suppose that h is diagonal and  $h = \text{diag}(a_1, a_2, \dots, a_n)$  with  $a_1 \ge a_2 \ge \dots \ge a_n$ . By Definition 2.5, we know that  $he_1 = a_1e_1$  and  $he_2 = a_2e_2$ . The assumption (1) becomes

$$\delta(\mathbb{R}e_1, y_q^m), \delta(\mathbb{R}e_2, y_q^m) > \delta, \gamma_{1,2}(g) \le \delta^3.$$

$$(4.40)$$

In (4.39), let  $z = z_o$ , then  $v = e_1, u = e_2$ , which make

$$\langle hv, hu \rangle = \langle a_1 e_1, a_2 e_2 \rangle = 0.$$

Therefore, due to

$$\langle v_1, v_2 \rangle \ge ||v_1|| ||v_2|| - ||v_1 \wedge v_2||,$$

for  $v_1, v_2$  in V, we have

$$F_{g,h}(e_1^2 \wedge e_1 e_2) = a_1^2 \langle g e_1, g e_2 \rangle \ge a_1^2 (\|g e_1\| \|g e_2\| - \|g e_1 \wedge g e_2\|).$$

Then (2.12) and (4.40) imply

$$F_{g,h}(e_1^2 \wedge e_1 e_2) \ge ||h||^2 ||g||^2 (\delta^2 - \gamma_{1,2}(g)).$$

The proof is complete.

By Definition 2.5, the representation  $\wedge^2 Sym^2 V$  is a proximal representation. Due to  $\mathbb{R}(v^2 \wedge vu) = \mathbb{R}k(e_1^2 \wedge e_1e_2) = kV^{\chi_1}$ , the line  $\mathbb{R}(v^2 \wedge vu)$  is contained in the K-orbit of the subspace of highest weight  $V^{\chi_1}$ . Since  $V^{\chi_1}$  is in W, we see that  $v^2 \wedge vu$  is also in W. By (4.39),

$$L = \frac{F_{g,h}(v^2 \wedge vu)}{\|F_{g,h}\|_W} \frac{\|g\|^2 \|h\|^2}{\|gv\|^2 \|hv\|^2} \frac{\|F_{g,h}\|_W}{\|g\|^2 \|h\|^2}.$$

When  $\eta$  satisfies assumption (2), the result follows by applying (2.12) to  $||gv||^2$ ,  $||hv||^2$  and by Lemma 4.19.

#### Proof of the spectral gap

Here we will prove the theorem of uniform spectral gap. The first part is classic, where we use some ideas of Dolgopyat [Dol98] to transform the problem to an effective estimate Proposition 4.24, see also [Nau05] and [Sto11]. The key observation is that this effective estimate (Proposition 4.24) can be obtained by the Fourier decay, regarding the difference of cocycle as a function on  $\mathscr{P}$ . The intuition here is from Lemma 4.17. When g, h are in general position and  $\eta$  not too close to  $\zeta_g^m, \zeta_h^m$ , the difference  $\varphi(\eta) = \sigma(g, \eta) - \sigma(h, \eta)$  will be (C, r) good (Definition 4.1). But in order to accomplish this, we need some sophisticate cutoff, which makes the proof complicate.

Recall that the Iwasawa cocycle takes values in the Cartan subspace  $\mathfrak{a}$ . From now on, we will the use another family of representations  $\{\tilde{V}_{\alpha}\}_{\alpha\in\Pi}$ , which is defined in (2.5) and whose highest weight  $\tilde{\chi}_{\alpha}$  is a multiple of  $\omega_{\alpha}$ . For simplifying the notation, we abbreviate  $\tilde{V}_{\alpha}$  to  $V_{\alpha}$ . Because

it is the only family of representation considered here. This family is also super proximal by Lemma 2.6.

We are in semisimple case and we know that  $\mathfrak{b}^* = \mathfrak{a}^*$ . We can write  $\vartheta$  in  $\mathfrak{a}^*$  as a linear combination of weights,  $\{\tilde{\chi}_{\alpha} | \alpha \in \Pi\}$ , that is

$$\vartheta = \sum_{\alpha \in \Pi} \vartheta_{\alpha} \tilde{\chi}_{\alpha}$$

Set  $|\vartheta| = \max_{\alpha \in \Pi} |\vartheta_{\alpha}|.$ 

We want to treat the spectral gap on the flag variety  $\mathscr{P}$  and the projective space  $\mathbb{P}V$  at the same time, where V is an irreducible representation of G with good norm. Let X be  $\mathscr{P}$  or  $\mathbb{P}V$ . Let  $\sigma: G \times X \to E$  be the cocycle, which is

- given by the semisimple part of the Iwasawa cocycle  $\sigma$  and  $E = \mathfrak{a}$  when  $X = \mathscr{P}$ ,
- given by  $\sigma_V$  (defined in (2.11)) and  $E = \mathbb{R}$  when  $X = \mathbb{P}V$ .

Let  $E_{\mathbb{C}} = E \otimes_{\mathbb{R}} \mathbb{C}$  and  $E_{\mathbb{C}}^*$  be the dual space of  $E_{\mathbb{C}}$ . For  $z \in E_{\mathbb{C}}^*$ , write  $z = \varpi + i\vartheta$ , where  $\vartheta, \varpi$  are elements in  $E^*$ . Recall that the transfer operator  $P_z$  is defined as: For  $|\varpi|$  small enough and for f in  $C^0(X)$ , x in X

$$P_z f(x) = \int_G e^{z\sigma(g,x)} f(gx) \mathrm{d}\mu(g).$$

Recall that for f in  $C^{\gamma}(X)$  let  $c_{\gamma}(f) = \sup_{x \neq x'} \frac{|f(x) - f(x')|}{d(x,x')^{\gamma}}$  and  $|f|_{\gamma} = |f|_{\infty} + c_{\gamma}(f)$ .

**Remark 4.20.** Here we should be careful that the distances on  $\mathbb{P}V$  and  $\mathscr{P}$  are defined in (2.9) and (2.17). They are not the Riemannian distances defined in the introduction. But on a compact Riemannian manifold, different Riemannian distances are equivalent. In particular, every Riemannian distance on  $\mathscr{P}$  is equivalent to the K-invariant Riemannian distance on  $\mathscr{P}$ . By Corollary 5.9, we know it is equivalent to the distances defined (2.17). The case of the projective space  $\mathbb{P}V$  is similar. Hence, the norm  $|\cdot|_{\gamma}$  induced by different distances are equivalent.

We state our main result of this section

**Theorem 4.21.** Let  $\mu$  be a Zariski dense Borel probability measure on G with a finite exponential moment. For  $\gamma > 0$  small enough, there exist  $\rho < 1, C > 0$  such that for all  $\vartheta$  and  $\varpi$  in  $E^*$  with  $|\vartheta|$  large enough,  $|\varpi|$  small enough and f in  $C^{\gamma}(X)$ , n in  $\mathbb{N}$  we have

$$|P_{\varpi+i\vartheta}^n f|_{\gamma} \le C |\vartheta|^{2\gamma} \rho^n |f|_{\gamma}$$

**Remark 4.22.** Compared with Theorem 1.3, we make an additional assumption that the norm on V is a good norm here. We explain here that for other norms the result also holds.

If we have another norm  $\|\cdot\|_1$  on V. Let  $\sigma_1$  be the new cocycle defined with respect to the norm  $\|\cdot\|_1$ . Let  $\psi(x) = \frac{\|v\|_1}{\|v\|}$  for  $x = \mathbb{R}v$  in  $\mathbb{P}V$ . Then

$$\sigma_1(g, x) = \sigma_V(g, x) + \log \psi(gx) - \log \psi(x),$$

which means the difference of two cocycles is a coboundary. This function  $\psi$  is Lipschitz, due to equivalence of norms on finite dimensional vector spaces. Let  $T_z f(x) = e^{z \log \psi(x)} f(x)$ . By Lipschitz property of  $\psi$ , we have

$$|T_z f|_{\gamma} \le C e^{C|a|} |z|^{\gamma} |f|_{\gamma},$$

where C depends on  $|\psi|_{Lip}$ . We know that

$$P_{z\sigma_1} = T_z^{-1} P_{z\sigma_V} T_z,$$

hence the same spectral gap property also holds for the norm  $\|\cdot\|_1$  with different constants.

Theorem 1.3 and Theorem 1.4 follow directly from Theorem 4.21. The assumption on  $\mu$  will be needed throughout this section.

We start with standard *a priori* estimates. When z = 0, we will write P for  $P_0$ .

**Proposition 4.23.** For every  $\gamma > 0$  small enough, there exist C > 0 and  $0 < \rho < 1$  such that for all f in  $C^{\gamma}(X)$ ,  $|\varpi|$  small enough and  $n \in \mathbb{N}$ 

$$|P_z^n f|_{\infty} \le C^{|\varpi|n} |f|_{\infty}, \tag{4.41}$$

$$|P^n f|_{\infty} \le \left| \int_X f \mathrm{d}\nu \right| + C\rho^n |f|_{\gamma}, \tag{4.42}$$

$$c_{\gamma}(P_z^n f) \le C(C^{|\varpi|n} |\vartheta|^{\gamma} |f|_{\infty} + \rho^n c_{\gamma}(f)).$$
(4.43)

The inequality (4.41) is a consequence of exponential moment and the Hölder inequality. For (4.42), please see [BL85, V, Thm.2.5] and [BQ16, Prop 11.10, Lem.13.5] for more details. This inequality (4.42) is a consequence of the fact that the action of G on X is contracting. The third inequality (4.43) is called the Lasota-Yorke inequality. The proof is classic.

We reduce Theorem 4.21 to Proposition 4.24. The reduction is standard, using Proposition 4.23. Please see [Dol98] for more details. For f in  $C^{\gamma}(X)$ , we define another norm  $|f|_{\gamma,\vartheta} = |f|_{\infty} + c_{\gamma}(f)/|\vartheta|^{\gamma}$  for  $\vartheta \neq 0$ .

**Proposition 4.24.** For every  $\gamma > 0$  small enough, for  $|\vartheta|$  large enough and  $|\varpi|$  small enough, there exist  $\epsilon_2$ ,  $C_2 > 0$  such that for f in  $C^{\gamma}(X)$  and  $|f|_{\gamma,\vartheta} \leq 1$ , we have

$$\int \left| P_{\varpi+i\vartheta}^{[C_2 \ln |\vartheta|]} f \right|^2 \mathrm{d}\nu \le e^{-\epsilon_2 \ln |\vartheta|}.$$
(4.44)

Now we will distinguish two cases. We claim that the case of  $\mathbb{P}V$  is a corollary of the case of  $\mathscr{P}$  up to a constant. Recall that the stationary measure on  $\mathbb{P}V$  is written as  $\nu_V$ . Let f be a function in  $C^{\gamma}(\mathbb{P}V)$  and  $|f|_{\gamma,\vartheta} \leq 1$ . The estimate only depends on the value of f on the support of the stationary measure  $\nu_V$ . By Lemma 2.48, the stationary measure on  $\mathbb{P}V$  is the pushforward measure of the stationary measure  $\nu$  on  $\mathscr{P}$ . Hence we can define the function  $\tilde{f}$  on  $\mathscr{P}$  by

$$f(\eta) = f(V_{\chi,\eta}),$$

where  $\chi$  is the highest weight of V. Then by  $\sigma_V(g, V_{\chi,\eta}) = \chi \sigma(g, \eta)$  (see (2.7)),

$$\int \left| P_{\varpi+i\vartheta}^{[C_2 \ln |\vartheta|]} f \right|^2 \mathrm{d}\nu_V = \int \left| P_{(\varpi+i\vartheta)\chi}^{[C_2 \ln |\vartheta|]} \tilde{f} \right|^2 \mathrm{d}\nu.$$

We will verify that  $\tilde{f}$  satisfies  $|\tilde{f}|_{\gamma,\vartheta} \ll 1$ . By (2.21), for two distinct points  $\eta, \eta'$  in  $\mathscr{P}$  we have

$$\frac{|\tilde{f}(\eta) - \tilde{f}(\eta')|}{d(\eta, \eta')^{\gamma}} = \frac{|\tilde{f}(\eta) - \tilde{f}(\eta')|}{d(V_{\chi,\eta}, V_{\chi,\eta'})^{\gamma}} \frac{d(V_{\chi,\eta}, V_{\chi,\eta'})^{\gamma}}{d(\eta, \eta')^{\gamma}} \ll \frac{|f(V_{\chi,\eta}) - f(V_{\chi,\eta'})|}{d(V_{\chi,\eta}, V_{\chi,\eta'})^{\gamma}} = |f|_{\gamma}.$$

Hence with some change of constant, we can deduce the case of  $\mathbb{P}V$  from the case of  $\mathscr{P}$ .

We only need to prove Proposition 4.24 for the case of  $\mathscr{P}$ .

From Fourier decay to Proposition 4.24. We need to reduce (4.44) to Fourier decay (Theorem 5.4). Let

$$n = [C_2 \log |\vartheta|] \text{ and } \delta = e^{-\epsilon n}$$
(4.45)

(with  $C_2 \ge \max_{\alpha \in \Pi} \{1/\alpha \sigma_{\mu}\} + 1$  and  $\epsilon > 0$  to be determined later), and let  $G_{n,\epsilon,\alpha}$  be the subset of  $G \times G$ , defined as the set of couples which satisfy Lemma 4.17 (1) with  $V = V_{\alpha}$ . Let

$$G_{n,\epsilon} = \{g \in G | \|\kappa(g) - n\sigma_{\mu}\| \le n\epsilon\}^2 \bigcap_{\alpha \in \Pi} G_{n,\epsilon,\alpha} \subset G \times G.$$

For  $|f|_{\gamma,\vartheta} \leq 1$ , let

$$A_{g,h} := \int_X e^{z\sigma(g,\eta) + \bar{z}\sigma(h,\eta)} f(g\eta) \bar{f}(h\eta) \mathrm{d}\nu(\eta).$$

Then

$$\int |P_z^n f|^2 d\nu = \int e^{z\sigma(g,\eta) + \bar{z}\sigma(h,\eta)} f(g\eta) \bar{f}(h\eta) d\nu(\eta) d\mu^{*n}(g) d\mu^{*n}(h)$$

$$= \int_{G_{n,\epsilon}} A_{g,h} d\mu^{*n}(g) d\mu^{*n}(h) + \int_{G_{n,\epsilon}^c} A_{g,h} d\mu^{*n}(g) d\mu^{*n}(h).$$
(4.46)

We first compute the term with (g, h) outside of  $G_{n,\epsilon}$ , where the behaviour is singular. By the Cauchy-Schwarz inequality,

$$\left| \int_{G_{n,\epsilon}^{c}} A_{g,h} \mathrm{d}\mu^{*n}(g) \mathrm{d}\mu^{*n}(h) \right|^{2} \leq \mu(G_{n,\epsilon}^{c}) \int |A_{g,h}|^{2} \mathrm{d}\mu^{*n}(g) \mathrm{d}\mu^{*n}(h).$$
(4.47)

By large deviation principle (Proposition 2.52, Proposition 2.53), the set  $G_{n,\epsilon}^c$  has exponentially small  $\mu^{*2n}$  measure, that is

$$\mu(G_{n,\epsilon}^c) \ll_{\epsilon} \delta^c. \tag{4.48}$$

By  $||f||_{\infty} \leq 1$  and (4.41), we have

$$\int |A_{g,h}|^2 \mathrm{d}\mu^{*n}(g) \mathrm{d}\mu^{*n}(h) \le |P_{2\varpi}^n \mathbb{1}|_{\infty}^2 \le C^{4n\varpi}.$$
(4.49)

When  $|\varpi|$  is small enough depending on  $\epsilon$ , by (4.47), (4.48) and (4.49)

$$\int_{G_{n,\epsilon}^c} A_{g,h} \mathrm{d}\mu^{*n}(g) \mathrm{d}\mu^{*n}(h) \ll_{\epsilon} \delta^{c/2} \le |\vartheta|^{-c\epsilon/(2C_2)}.$$
(4.50)

We compute the major term, that is (g,h) in  $G_{n,\epsilon}$ . We want to use Theorem 5.4 to control this part with  $\varphi = |\vartheta|^{-1} \vartheta(\sigma(g,\eta) - \sigma(h,\eta))$  and a suitable r. For applying Theorem 5.4, we need that  $\varphi$  is (C,r) good, which will be accomplished by multiplying smooth cutoffs. The most important is G2 assumption (4.2), which will be verified with the help of Lemma 4.17. Hence we want that r vanishes when  $\eta$  does not satisfy Lemma 4.17 (2).

Let  $X_{g,h,\alpha}$  be the subset of  $\mathscr{P}$ , defined as the set of elements which satisfy Lemma 4.17 (2) with  $V = V_{\alpha}$ . Let  $X_{g,h} = \bigcap_{\alpha \in \Pi} X_{g,h,\alpha}$ . Recall that  $\tau$  be a smooth function on  $\mathbb{R}$  such that  $\tau|_{[0,\infty)} = 1, \tau$  takes values in [0,1],  $\operatorname{supp} \tau \subset [-1,\infty)$  and  $|\tau'| \leq 2$ . For  $\delta > 0$ , set  $\tau_{\delta}(x) = \tau(x/\delta)$ for  $x \in \mathbb{R}$ . Let  $\sigma_{\alpha} := \sigma_{V_{\alpha}} = \tilde{\chi}_{\alpha}\sigma$  and

$$\varphi(\eta) = |\vartheta|^{-1} \vartheta(\sigma(g,\eta) - \sigma(h,\eta)) = |\vartheta|^{-1} \sum_{\alpha \in \Pi} \vartheta_{\alpha}(\sigma_{\alpha}(g,\eta) - \sigma_{\alpha}(h,\eta))$$
(4.51)

and

$$r(\eta) = f(g\eta)\bar{f}(h\eta)e^{\varpi(\sigma(g,\eta) + \sigma(h,\eta))} \prod_{\alpha \in \Pi} \tau_{\alpha}, \qquad (4.52)$$

where

$$\tau_{\alpha}(\eta) = \tau_{\delta}(4\delta_{\alpha}(\eta, \zeta_g^m) - 4\delta)\tau_{\delta}(4\delta_{\alpha}(\eta, \zeta_h^m) - 4\delta)\tau_{\delta}(4\delta(V_{4\tilde{\chi}_{\alpha}-\alpha,\eta}, F_{\rho_{\alpha}g,\rho_{\alpha}h}) - 4\delta),$$

where  $\delta_{\alpha}$  is defined to be

$$\delta_{\alpha}(\eta, \zeta_g^m) = \delta(V_{\alpha,\eta}, y_{\rho_{\alpha}(g)}^m).$$

The choice of  $\tau_{\alpha}$  is sophisticate. We only need to keep in mind that they come from Lemma 4.17. Then  $e^{i|\vartheta|\varphi}r(\eta)$  equals  $e^{z\sigma(g,\eta)+\bar{z}\sigma(h,\eta)}f(g\eta)f(h\eta)$  on  $X_{g,h}$ .

**Lemma 4.25.** Let  $\epsilon_0, \epsilon_1$  be given by Theorem 5.4. Let (g, h) be in  $G_{n,\epsilon}$ . With  $\epsilon$  small enough depending on  $\epsilon_0$  and  $|\varpi|$  small enough depending on  $\epsilon$  and  $\epsilon_1$ , for  $\varphi, r$  defined in (4.51) and (4.52) we have that  $\varphi$  is  $(|\vartheta|^{\epsilon_0}, r)$  good and  $c_{\gamma}(r) \leq |\vartheta|^{\epsilon_0}, |r|_{\infty} \leq |\vartheta|^{\epsilon_1/2}$ .

By Lemma 4.25, we can fix a value of  $\epsilon$  and the functions  $\varphi$  and  $r|\vartheta|^{-\epsilon_1/2}$  satisfy the condition in Theorem 5.4. (Theorem 5.4 still holds when r is a complex function) Hence Theorem 5.4 implies

$$\left| \int e^{i|\vartheta|\varphi(\eta)} r(\eta) \mathrm{d}\nu(\eta) \right| \le |\vartheta|^{-\epsilon_1/2}.$$
(4.53)

The difference between  $A_{q,h}$  and  $\int e^{i|\vartheta|\varphi(\eta)}r(\eta)d\nu(\eta)$  is bounded by

$$\nu(X_{g,h}^c) \le \sum_{\alpha \in \Pi} \nu(X_{g,h,\alpha}^c).$$
(4.54)

Using the regularity of stationary measure (2.40) with  $V = W_{\alpha}$ , the irreducible subrepresentation of  $\wedge^2 Sym^2 V_{\alpha}$  with the highest weight, we have

$$\nu\{\eta \in \mathscr{P} | \delta(V_{4\tilde{\chi}_{\alpha}-\alpha,\eta}, F_{\rho_{\alpha}g,\rho_{\alpha}h}) < \delta\} \ll_{\epsilon} e^{-c\epsilon n}.$$

$$(4.55)$$

Using the regularity of stationary measure (2.40) with  $V = V_{\alpha}$ , we obtain

$$\nu\{\eta \in \mathscr{P} | V_{\alpha,\eta} \in B_h^m(\delta) \cup B_g^m(\delta)\} \ll_{\epsilon} e^{-c\epsilon n}.$$
(4.56)

Hence by (4.54)-(4.56), we have

$$\nu(X_{g,h}^c) \ll_{\epsilon} e^{-c\epsilon n} = |\vartheta|^{-c\epsilon/C_2}.$$
(4.57)

For (g, h) in  $G_{n,\epsilon}$ , by (4.53) and (4.57)

$$A_{g,h} \ll |\vartheta|^{-\epsilon_1/2} + |\vartheta|^{-c\epsilon/C_2}$$

Combined with (4.46) and (4.50), the proof is complete by setting  $\epsilon_2 = \min\{\frac{\epsilon_1}{2}, \frac{c\epsilon}{4C_2}\}$ .

It remains to prove Lemma 4.25.

Proof of Lemma 4.25. We first verify that  $\varphi$  is  $(|\vartheta|^{\epsilon_0}, r)$  good. Since  $\epsilon$  will be taken small enough, we can suppose  $|\vartheta|^{-\epsilon_0} \leq \delta/4$ . Let J be the  $|\vartheta|^{-\epsilon_0}$  neighbourhood of suppr. Then for  $\eta \in J$ , we have  $\delta_{\alpha}(\eta, \zeta_q^m) \geq \delta/2$  for  $\alpha$  in  $\Pi$ .

The function  $\varphi$  is a sum of functions. Each function is the lift of a function on  $\mathbb{P}V_{\alpha}$  for some simple root  $\alpha$ . We write  $\varphi = \sum_{\alpha \in \Pi} \varphi_{\alpha}$  where  $\varphi_{\alpha}(\eta) = |\vartheta|^{-1} \vartheta_{\alpha}(\sigma_{\alpha}(g,\eta) - \sigma_{\alpha}(h,\eta))$ . By Lemma 2.37, that is  $\partial_{\alpha'}\varphi_{\alpha} = 0$  for  $\alpha' \neq \alpha$ , in order to verify  $(|\vartheta|^{\epsilon_0}, r)$  good condition, it is enough to verify G1-G3 assumptions (4.1)-(4.3) for  $\varphi_{\alpha}$  and the G4 assumption (4.4) for  $\varphi$ . Since G1-G3 are linear, we can forget the coefficients  $|\vartheta|^{-1}\vartheta_{\alpha}$  in  $\varphi_{\alpha}$ .

Now, we verify G1-G3 assumptions and we fix a simple root  $\alpha$  and consider  $\varphi = \varphi_{\alpha} = \sigma_{\alpha}(g, \cdot) - \sigma_{\alpha}(h, \cdot)$ . Recall that  $v_{\alpha} = \sup_{\eta \in \text{suppr}} |\partial_{\alpha}\varphi(\eta)|$ . Since J satisfies the hypothesis of Lemma 4.16 with  $V = V_{\alpha}$ , we have

$$v_{\alpha}, Lip_{\mathscr{P}_0}(\partial_{\alpha}\varphi|_{\pi^{-1}J}) < \delta^{-O(1)}.$$

$$(4.58)$$

Since  $(g,h) \in G_{n,\epsilon}$  satisfies Lemma 4.17(1) and the support of r satisfies Lemma 4.17(2), for  $\eta$  in the support of r, by Lemma 4.17,

$$|\partial_{\alpha}\varphi(\eta)| > \delta^{O(1)} \ge \delta^{O(1)} v_{\alpha}$$

which is G2 assumption (4.2). This also implies

$$v_{\alpha} > \delta^{O(1)},\tag{4.59}$$

G4 assumption (4.4). By (4.58), we have G3 assumption (4.3). Let  $\varphi_1$  be a function on  $\mathbb{P}V_{\alpha}$  such that  $\varphi_1(V_{\alpha,\eta}) = \varphi(\eta)$ . Since J satisfies hypothesis of Lemma 2.11, this Lemma implies

$$\frac{|\varphi(\eta) - \varphi(\eta')|}{d_{\alpha}(\eta, \eta')} = \frac{|\varphi_1(V_{\alpha, \eta}) - \varphi_1(V_{\alpha, \eta'})|}{d(V_{\alpha, \eta}, V_{\alpha, \eta'})} \le |Lip_{\mathbb{P}V_{\alpha}}\varphi_1| < \delta^{-O(1)} \le \delta^{-O(1)}v_{\alpha},$$

which is G1 assumption (4.1).

For general  $\varphi$ , it remains to verify G4 assumption (4.4). There exists a simple root  $\alpha$  such that  $|\vartheta_{\alpha}| = |\vartheta|$ . Since  $\varphi_{\alpha}$  satisfies G4 assumption and  $|\partial_{\alpha}\varphi| = |\partial_{\alpha}\varphi_{\alpha}|$  by Lemma 2.37, the function  $\varphi$  also satisfies G4 assumption.

Finally, we verify the term  $c_{\gamma}(r)$  and  $|r|_{\infty}$ .

**Lemma 4.26.** For  $0 < \gamma \leq 1$ , let  $f, \tau$  be two  $\gamma$ -Hölder functions on a compact metric space X. Then

 $c_{\gamma}(\tau f) \leq c_{\gamma}(\tau) \|f|_{\mathrm{supp}\tau}\|_{\infty} + |\tau|_{\infty} c_{\gamma}(f|_{\mathrm{supp}\tau}).$ 

The proof of Lemma 4.26 is elementary. Recall that

$$r(\eta) = f(g\eta)\bar{f}(h\eta)e^{\varpi(\sigma(g,\eta) + \sigma(h,\eta))} \prod_{\alpha \in \Pi} \tau_{\alpha}.$$

For the infinity norm, due to  $(g,h) \in G_{n,\epsilon}$ , we have

$$|r| \le e^{|\varpi|(\|\kappa(g)\| + \|\kappa(h\|)} \le e^{|\varpi|(2\|\sigma_{\mu}\| + \epsilon)n} \le |\vartheta|^{|\varpi|C_2(2\|\sigma_{\mu}\| + \epsilon)}.$$

Take  $|\varpi|$  small enough, then  $|r|_{\infty} \leq |\vartheta|^{\epsilon_1/2}$ .

For the term  $c_{\gamma}(r)$ , we only need to verify that each term in the formula of r has a bounded  $c_{\gamma}$  value. Due to Lemma 4.26, we only need to verify the  $c_{\gamma}$  value on  $X_{q,h}$ .

• Since the action of g on  $X_{g,h}$  is contracting, by Lemma 2.16, we have

$$c_{\gamma}(f(g \cdot)|_{X_{q,h}}) \le c_{\gamma}(f)(Lip \ g|_{X_{q,h}})^{\gamma} \le (|\vartheta|\beta\delta^{-2})^{\gamma}.$$

Due to (4.45), we have  $\log \beta = -n \min_{\alpha \in \Pi} \alpha \sigma_{\mu} < -n/C_2 \leq -\log |\vartheta|$ . Therefore  $c_{\gamma}(f(g \cdot)|_{X_{g,h}}) \leq \delta^{-O(1)}$ .

• Due to

$$|e^a - e^b| \le \max\{e^a, e^b\}|a - b|^{\gamma}$$

for all a, b in  $\mathbb{R}$  and  $0 \le \gamma \le 1$ , by Lemma 2.16,

$$c_{\gamma}(e^{\varpi\sigma(g,\cdot)}|_{X_{g,h}}) \leq e^{|\varpi| \|\kappa(g)\|} (Lip \varpi\sigma(g,\cdot)|_{X_{g,h}})^{\gamma} \leq e^{|\varpi|(\|\sigma_{\mu}\|+\epsilon)n+\epsilon\gamma n} |\varpi|^{\gamma}.$$

Hence when  $|\varpi|$  is small enough depending on  $\sigma_{\mu}$ , we obtain  $c_{\gamma}(e^{\varpi\sigma(g,\cdot)}|_{X_{g,h}}) \leq \delta^{-O(1)}$ .

• In  $c_{\gamma}(\tau_{\alpha})$ , the only term we need to be careful about is  $\tau_{\delta}(4\delta(V_{4\tilde{\chi}_{\alpha}-\alpha,\eta}, F_{\rho_{\alpha}g,\rho_{\alpha}h}) - 4\delta)$ . By Lemma 2.17, we have  $d(V_{4\tilde{\chi}_{\alpha}-\alpha,\eta}, V_{4\tilde{\chi}_{\alpha}-\alpha,\eta'}) \ll d(\eta,\eta')$ . Hence the  $c_{\gamma}$  value of this term is also bounded by  $\delta^{-O(1)}$ .

The proof is complete.

### 4.4 Exponential error term

In this section, we will prove Theorem 1.1 that the speed of convergence in the renewal theorem is exponential using our result on the spectral gap. (Theorem 4.21) Recall  $X = \mathbb{P}V$ , where V is an irreducible representation of G with a norm and the highest weight  $\chi$  is in  $\mathfrak{b}^*$ . We have defined a renewal operator R as follows: For a positive bounded Borel function f on  $\mathbb{R}$ , a point x in X and a real number t, we set

$$Rf(x,t) = \sum_{n=0}^{+\infty} \int_G f(\sigma_V(g,x) - t) \mathrm{d}\mu^{*n}(g)$$

and

$$R_P f(x,t) = \sum_{n=0}^{+\infty} \int_G f(\log \|\rho(g)\| - t) \mathrm{d}\mu^{*n}(g).$$

Recall  $P_z$  is the transfer operator defined by  $P_z f(x) = \int_G e^{z\sigma_V(g,x)} f(gx) d\mu(g)$ . Using the analytical Fredholm theorem, we summarize the property of  $P_z$ .

**Proposition 4.27.** With the same assumption as in Theorem 1.1, for any  $\gamma > 0$  small enough, there exists  $\eta > 0$  such that when  $|\Re z| < \eta$ , the transfer operator  $P_z$  is a bounded operator on  $C^{\gamma}(X)$  and depends analytically on z. Moreover there exists an analytic operator U(z) on a neighbourhood of  $|\Re z| < \eta$  such that the following holds for  $|\Re z| < \eta$ 

$$(I - P_z)^{-1} = \frac{1}{\sigma_{V,\mu z}} N_0 + U(z),$$

where  $N_0$  is the operator defined by  $N_0 f = \int_X f d\nu_V$ . There exists C > 0 such that for  $|\Re z| \leq \eta$ 

 $||U(z)||_{C^{\gamma} \to C^{\gamma}} \le C(1+|\Im z|)^{2\gamma}.$ (4.60)

This is a generalization of [LI18a, Prop. 4.1] and [Boy16, Theorem 4.1], and the proof is exactly the same. The main difference is that the spectral radius of  $P_z$  is bounded below 1 in a strip of imaginary line (except at 0), due to Theorem 4.21. From this we have the analytic continuation of U(z) to the strip and the bound of the operator norm of U(z).

Now, we give the precise statement and the proof of Theorem 1.1.

**Proposition 4.28.** With the same assumption as in Theorem 1.1, there exists  $\epsilon > 0$  such that for  $f \in C_c^{\infty}(\mathbb{R})$ , we have

$$Rf(x,t) = \frac{1}{\sigma_{V,\mu}} \int_{-t}^{\infty} f(u) \mathrm{d}u + e^{-\epsilon|t|} O(e^{\epsilon|\mathrm{supp}f|} (|f''|_{L^1} + |f|_{L^1})),$$

where |supp f| is the supremum of the absolute value of x in supp f.

*Proof.* By the same computation as in [L118a, Lemma 4.5] and [Boy16, Prop. 4.14], we have

$$Rf(x,t) = \frac{1}{\sigma_{V,\mu}} \int_{-t}^{\infty} f(u) du + \lim_{s \to 0^+} \frac{1}{2\pi} \int e^{-it\xi} \hat{f}(\xi) U(s+i\xi) \mathbb{1}(x) d\xi,$$

where  $\hat{f}$  is the Fourier transform of f given by  $\hat{f}(\xi) = \int e^{i\xi u} f(u) du$ . Hence, we only need to control the error term.

By Proposition 4.27, we know that U(z) is analytical on  $\{z \in \mathbb{C} | |\Re z| \leq \eta\}$  and uniformly bounded by  $(1 + |\Im z|)^{2\gamma}$ . Since f is a compactly supported smooth function, the Fourier transform  $\hat{f}$  is an analytic function on  $\mathbb{C}$ . By  $|\hat{f}(i\epsilon + \xi)| \leq e^{\epsilon |\operatorname{supp} f|} \frac{1}{|\xi|^2} |f''|_{L^1}$ , and  $|\hat{f}(i\epsilon + \xi)| \leq e^{\epsilon |\operatorname{supp} f|} |f|_{L^1}$  for  $\epsilon, \xi$  in  $\mathbb{R}$ , we have

$$|\hat{f}(i\epsilon + \xi)| \le e^{\epsilon |\operatorname{supp} f|} \frac{2}{1 + |\xi|^2} (|f''|_{L^1} + |f|_{L^1}).$$
(4.61)

By (4.60), (4.61) and the dominant convergence theorem, we have

$$\lim_{s \to 0^+} \frac{1}{2\pi} \int e^{-it\xi} \hat{f}(\xi) U(s+i\xi) \mathbb{1}(x) d\xi = \frac{1}{2\pi} \int e^{-it\xi} \hat{f}(\xi) U(i\xi) \mathbb{1}(x) d\xi.$$
(4.62)

**Lemma 4.29.** [RS75, Thm.IX14] If T is in  $S'(\mathbb{R})$ , tempered distributions, the distribution T has analytic continuation to  $|\Im\xi| < a$  and  $\sup_{|b| < a} \int |T(ib+y)| dy < \infty$ , then  $\check{T}$  is a continuous function. For all b < a, let  $C_b = \max \int |T(\pm ib+y)| dy$ . We have

$$|\check{T}(t)| \le C_b e^{-b|t|}$$

Using Lemma 4.29 with  $T(\xi) = \hat{f}(\xi)U(i\xi)\mathbb{1}(x)$ , we have

$$\left| \int \hat{f}(\xi) U(i\xi) \mathbb{1}(x) e^{-it\xi} \mathrm{d}\xi \right| = |\check{T}(t)| \le e^{-\epsilon|t|} \max |T(\pm i\epsilon + \xi)|_{L^1(\xi)}.$$
(4.63)

By (4.61), we have

$$\max |T(\pm i\epsilon + \xi)|_{L^{1}(\xi)} \leq e^{\epsilon |\operatorname{supp} f|} \int \frac{2}{1 + |\xi|^{2}} (|f''|_{L^{1}} + |f|_{L^{1}}) |U(\mp \epsilon + i\xi) \mathbb{1}(x)| d\xi \ll_{\gamma} e^{\epsilon |\operatorname{supp} f|} (|f''|_{L^{1}} + |f|_{L^{1}}).$$

$$(4.64)$$

Combining (4.62), (4.63) and (4.64), we have the result.

**Proposition 4.30.** With the same assumption as in Theorem 1.1, there exists  $\epsilon > 0$  such that for  $f \in C_c^{\infty}(\mathbb{R})$ , we have

$$R_P f(x,t) = \frac{1}{\sigma_{V,\mu}} \int_{-t}^{\infty} f(u) \mathrm{d}u + e^{-\epsilon|t|} O(e^{\epsilon|\mathrm{supp}f|} (|f''|_{L^1} + |f|_{L^1})).$$

*Proof.* The ideal of the proof is the same as [L118a, Lemma 4.11 or Proposition 4.28], where we only need to replace the error term by the error term in the above Proposition 4.28.

We summarize the main idea here. By the large deviation principle, the main contribution of the renewal sum is given by n in a small interval containing  $t/\sigma_{V,\mu}$ . Since the norm is good, we have the interpretation of the norm by the Cartan projection (2.8). Then we use [BQ16, Lemma 17.8] to replace the norm by the cocycle  $\sigma_V$  for each n in the small interval. The proof is complete by applying Proposition 4.28.

# 5 Appendix

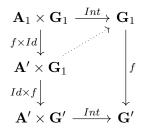
## 5.1 Non simply connected case

We explain here how to get Theorem 5.4 for connected algebraic semisimple Lie groups defined and split over  $\mathbb{R}$  from Theorem 1.7 for connected  $\mathbb{R}$ -split reductive  $\mathbb{R}$ -groups whose semisimple part is simply connected, which is proved in Section 4.

See [Mar91] and [Bor90, §22] for more facts about algebraic groups and central isogeny.

**Lemma 5.1.** Let  $\mathbf{G}'$  be a connected algebraic semisimple Lie groups defined over  $\mathbb{R}$ . Then there exist a connected reductive  $\mathbb{R}$ -group  $\mathbf{G}$  with simply connected derived group  $\mathscr{D}\mathbf{G}$  and an algebraic group morphism  $\psi : \mathbf{G} \to \mathbf{G}'$  which is surjective between real points. Moreover, the restriction of  $\psi$  to  $\mathscr{D}\mathbf{G}$  gives a central isogeny from  $\mathscr{D}\mathbf{G}$  to  $\mathbf{G}'$  and the connected centre of  $\mathbf{G}$  is  $\mathbb{R}$ -split.

Proof. Let  $\mathbf{A}'$  be a maximal  $\mathbb{R}$ -split torus of  $\mathbf{G}'$ . Let  $\mathbf{G}_1$  be a cover of  $\mathbf{G}'$  which is simply connected and let f be the isogeny map from  $\mathbf{G}_1$  to  $\mathbf{G}'$ . Let  $\mathbf{A}_1$  be the preimage of  $\mathbf{A}'$  in  $\mathbf{G}_1$ , which is a maximal  $\mathbb{R}$ -split torus of  $\mathbf{G}_1$  [Bor90, Theorem 22.6 (ii)]. Let  $\mathbf{N} = \ker f \cap \mathbf{A}_1$ , then  $\mathbf{A}'$ is isomorphic to  $\mathbf{A}_1/\mathbf{N}$  as torus. Consider the conjugate action of  $\mathbf{A}_1$  on  $\mathbf{G}_1$ , that is for  $s \in \mathbf{A}_1$ and  $g \in \mathbf{G}_1$  we define  $Int_s(g) = s^{-1}gs$ . Since the kernel of f is in the centre of  $\mathbf{G}_1$ , the conjugate action of  $\mathbf{N}$  on  $\mathbf{G}_1$  is trivial. By [Bor90, Corollary 6.10], the quotient group  $\mathbf{A}' \simeq \mathbf{A}_1/\mathbf{N}$  acts  $\mathbb{R}$ -morphically on  $\mathbf{G}_1$ .



Hence, we can define the semidirect product  $\mathbf{G} = \mathbf{A}' \ltimes \mathbf{G}_1$ , given by the action of  $\mathbf{A}'$  on  $\mathbf{G}_1$ . The derived group  $[\mathbf{G}, \mathbf{G}]$  equals  $\mathbf{G}_1$ , which is simply connected. The group  $\mathbf{G}$  is defined over  $\mathbb{R}$ , because  $\mathbf{A}', \mathbf{G}_1$  and  $\psi$  are also. The restriction of the action of  $\mathbf{A}'$  on  $\mathbf{A}_1$  is trivial and  $\mathbf{A}' \times \mathbf{A}_1$  is a maximal  $\mathbb{R}$ -split torus of  $\mathbf{G}$ . Hence the group  $\mathbf{G}$  is a connected reductive  $\mathbb{R}$ -group.

We only need to find the surjective morphism  $\psi$ . Let  $\mathbf{A}' \ltimes \mathbf{G}'$  be the semidirect product given by the conjugation action of  $\mathbf{A}'$  on  $\mathbf{G}'$ . As  $\mathbf{A}'$  is a subgroup of  $\mathbf{G}'$ , this semidirect product is isomorphic to a product. We have a group morphism

$$\mathbf{G} = \mathbf{A}' \ltimes \mathbf{G}_1 \to \mathbf{A}' \ltimes \mathbf{G}' \to \mathbf{G}'$$
$$\psi : (s,g) \mapsto (s,f(g)) \mapsto sf(g).$$

It is well-known that the real part of a semisimple simply connected group  $G_1$  is connected in analytic topology. (See for example [Ste68]) Let  $(G')^o$  be the analytic connected component of the identity element in G'. Then the image of real points of **G** under  $\psi$  is  $A'(G')^o$ , which is equal to G' by a theorem of Matsumoto [Mat64] ([BT65, Théorème 14.4]).

**Example 5.2.** When  $\mathbf{G}' = \mathbf{P}\mathbf{G}\mathbf{L}_2$ , the above construction gives  $\mathbf{G} = \mathbf{G}\mathbf{L}_2 = \mathbf{G}\mathbf{L}_1 \ltimes \mathbf{S}\mathbf{L}_2$  and the map  $\psi$  is the quotient map from  $\mathbf{G}\mathbf{L}_2$  to  $\mathbf{P}\mathbf{G}\mathbf{L}_2$ .

Let  $G' = \mathbf{G}'(\mathbb{R})$  be the group of real points of a connected algebraic semisimple Lie groups defined and split over  $\mathbb{R}$ . Recall that  $\mu$  is a Zariski dense Borel probability measure on G' with a finite exponential moment. If  $\mathbf{G}'$  is simply connected, then Theorem 1.7 holds for  $\mathbf{G}'$ . If not, let  $G = \mathbf{G}(\mathbb{R})$  be as in Lemma 5.1. Recall that  $\psi$  is a group morphism from G to G',

**Lemma 5.3.** There exists a Zariski dense Borel probability measure  $\tilde{\mu}$  on G with a finite exponential moment such that

$$\psi_* \tilde{\mu} = \mu. \tag{5.1}$$

The proof of Lemma 5.3 will be given at the end this section. We will explain why the results also hold for G' and  $\mu$ . We state the non simply connected version of Theorem 1.7 here

**Theorem 5.4** (Fourier decay). Let  $\mathbf{G}'$  be a connected algebraic semisimple Lie group defined and split over  $\mathbb{R}$  and let  $\mathbf{G}' = \mathbf{G}'(\mathbb{R})$  be its group of real points. Let  $\mu$  be Zariski dense Borel probability measure on  $\mathbf{G}'$  with finite exponential moment. Let  $\nu$  be the  $\mu$ -stationary measure on the flag variety  $\mathcal{P}$ .

For every  $\gamma > 0$ , there exist  $\epsilon_0 > 0, \epsilon_1 > 0$  depending on  $\mu$  such that the following holds. For any pair of real functions  $\varphi \in C^2(\mathscr{P}), r \in C^{\gamma}(\mathscr{P})$  and  $\xi > 0$  such that  $\varphi$  is  $(\xi^{\epsilon_0}, r)$  good,  $\|r\|_{\infty} \leq 1$  and  $c_{\gamma}(r) \leq \xi^{\epsilon_0}$ , then

$$\left| \int e^{i\xi\varphi(\eta)} r(\eta) \mathrm{d}\nu(\eta) \right| \le \xi^{-\epsilon_1} \quad \text{for all } \xi \text{ large enough.}$$
(5.2)

*Proof.* By Lemma 5.1, since  $\mathbf{G}'$  is  $\mathbb{R}$ -split,  $\mathbf{G}$  is also  $\mathbb{R}$ -split. Hence Theorem 1.7 holds for  $G, \tilde{\mu}$ . By Lemma 5.3, we only need to prove the flag varieties of G and G' are isomorphic, then the result follows.

By [Bor90, Prop.20.5], we know that  $(\mathbf{G}_2/\mathbf{P}_2)(\mathbb{R}) = \mathbf{G}_2(\mathbb{R})/\mathbf{P}_2(\mathbb{R})$  for any connected reductive  $\mathbb{R}$ -group  $\mathbf{G}_2$  and its parabolic  $\mathbb{R}$ -subgroup  $\mathbf{P}_2$ . Hence it is sufficient to prove for a minimal parabolic  $\mathbb{R}$ -subgroup  $\mathbf{P}$  of  $\mathbf{G}$  and  $\mathbf{P}'$  its image in  $\mathbf{G}'$  that

$$\mathbf{G}/\mathbf{P} \simeq \mathbf{G}'/\mathbf{P}'.\tag{5.3}$$

As if (5.3) holds, then  $\mathbf{P}'$  is also a parabolic subgroup by definition and it is minimal because  $\mathbf{P}$  is. Due to [Bor90, Thm.11.16], the normalizer of a parabolic subgroup is itself. Then the centre of  $\mathbf{G}$  is contained in the parabolic group  $\mathbf{P}$ . It suffices to prove that ker  $\psi$  is in the centre, then

$$\mathbf{G}/\mathbf{P} \simeq (\mathbf{G}/\ker\psi)/(\mathbf{P}/\ker\psi) \simeq \mathbf{G}'/\mathbf{P}'.$$

By [Bor90, Prop.14.2], we know that  $\mathbf{G} = \mathbf{C} \cdot \mathscr{D}\mathbf{G}$ , where  $\mathbf{C}$  is the connected centre and  $\mathbf{C} \cap \mathscr{D}\mathbf{G}$  is finite. Since  $\mathbf{G}'$  is semisimple, the connected centre  $\mathbf{C}$  is in ker  $\psi$ . As the restriction of  $\psi$  on  $\mathscr{D}\mathbf{G}$  to  $\mathbf{G}'$  is a central isogeny, hence ker  $\psi \cap \mathscr{D}\mathbf{G}$  is in the centre of  $\mathscr{D}\mathbf{G}$ , which is also in the centre of  $\mathbf{G}$ . Therefore the kernel of  $\psi$  is in the centre of  $\mathbf{G}$ . The proof is complete.  $\Box$ 

It remains to prove Lemma 5.3.

Proof of Lemma 5.3. We will first construct a measure  $\tilde{\mu}_1$  which has a finite exponential moment. In the construction of Lemma 5.1, there exists a finite subgroup F of A' such that  $\psi$  from  $F \ltimes G_1$  to G' is already surjective. Let  $F_1$  be the kernel of this covering, which is finite. Then there exists a unique Borel probability measure  $\tilde{\mu}_1$  on  $F \ltimes G_1 < G$  which is  $F_1$ -left invariant and the pushforward measure is  $\mu$ .

The moment condition is also satisfied. Because  $\psi$  induce an isomorphism between  $\mathfrak{a}_1$  to  $\mathfrak{a}'$ (Recall the notation in Section 2) and this isomorphism identifies the Cartan projections  $\kappa(g)$ and  $\kappa(\psi(g))$ . Let mg be an element in  $F \ltimes G_1$  with  $m \in F$  and  $g \in G_1$ , then by the sub additivity of the Cartan projection ([BQ16, Corollary 8.20]),

$$\|\kappa(mg)\| \le \|\kappa(m)\| + \|\kappa(g)\| = \|\kappa(m)\| + \|\kappa(\psi(g))\| \le \|\kappa(m)\| + \|\kappa(\psi(m))\| + \|\kappa(\psi(mg))\|.$$

Hence

$$\int_{G} e^{\epsilon \|\kappa(g)\|} \mathrm{d}\tilde{\mu}_1(g) \ll \int_{G} e^{\epsilon \|\kappa(\psi(g))\|} \mathrm{d}\tilde{\mu}_1(g) = \int_{G'} e^{\epsilon \|\kappa(g')\|} \mathrm{d}\mu(g').$$

In order to get a Zariski dense measure  $\tilde{\mu}$ , we replace the above measure  $\tilde{\mu}_1$  by  $\tilde{\mu} = \frac{1}{2}(\tilde{\mu}_1 + c_*\tilde{\mu}_1)$ , where c is an element in the connected centre C such that the group  $C_1 = \langle c \rangle$  generated by c is Zariski dense in C. Due to  $d\psi|_c = 0$ , the connected centre C is in the kernel of  $\psi$ . Hence  $\psi_*(c_*\tilde{\mu}_1) = \psi_*\tilde{\mu}_1$ . This measure  $\tilde{\mu}$  satisfies (5.1) and has a finite exponential moment. We will prove that it is also Zariski dense.

Let **H** be the Zariski closure of  $\Gamma_{\tilde{\mu}}$ , the group generated by the support of  $\tilde{\mu}$ . Let  $\mathfrak{h}$  be the Lie algebra of H. Since the group **G** is a connected  $\mathbb{R}$ -group, it is sufficient to prove that  $\mathfrak{h} = \mathfrak{g}$ . Recall that  $\mathfrak{g} = \mathfrak{c} \oplus \mathscr{D}\mathfrak{g}$ . Due to c in H, the Zariski closure of the  $C_1$  is also in H. Hence  $\mathfrak{h} \supset \mathfrak{c}$ . For the semisimple part, consider the adjoint action of  $\Gamma_{\tilde{\mu}}$  on  $\mathscr{D}\mathfrak{g}$ . Because the group  $\Gamma_{\mu}$ is Zariski dense in  $\mathbf{G}'$ , the adjoint action of  $\Gamma_{\mu}$  on  $\mathfrak{g}'$  is irreducible. The map  $d\psi|_{\mathscr{D}\mathfrak{g}} : \mathscr{D}\mathfrak{g} \to \mathfrak{g}'$ is an isomorphism of Lie algebras. By

$$\mathrm{d}\psi(\mathrm{Ad}_q X) = \mathrm{Ad}_{\psi(q)} \mathrm{d}\psi X \text{ for } X \in \mathscr{D}\mathfrak{g},$$

we obtain that the action of  $\Gamma_{\tilde{\mu}}$  on  $\mathscr{D}\mathfrak{g}$  is irreducible. Since  $\mathfrak{h} \cap \mathscr{D}\mathfrak{g}$  is nonzero and  $\Gamma_{\tilde{\mu}}$ -invariant, we know that  $\mathfrak{h} \cap \mathscr{D}\mathfrak{g} = \mathscr{D}\mathfrak{g}$ . Therefore  $\mathfrak{h} = \mathfrak{g}$ . The proof is complete.

## 5.2 Equivalence of distances

**Definition 5.5.** Let (X, d) be a metric space. Let d' be another metric on X. We say that d, d' are equivalent metrics if there exist c, C > 0 such that for all  $x_1, x_2$  in X

$$cd(x_1, x_2) \le d'(x_1, x_2) \le Cd(x_1, x_2)$$

Recall that  $\mathscr{P}_0$  is the homogeneous space  $G/A_eN$ , on which the compact group K acts simply transitively. Recall that  $\{V_\alpha\}_{\alpha\in\Pi}$  is the family of representation fixed in Lemma 2.3. We will define three distances on  $\mathscr{P}_0$ . Due to the fact that  $\mathscr{P}_0$  is homeomorphic to K, a distance on  $\mathscr{P}_0$  is also a distance on K and we will continue our argument on K. Let k, k' be two points in K. If they are not in the same connected component, we define their distance as 1. From now on, we always suppose that k, k' are in the connected component  $K^o$ .

- $d_0(k, k') = \sup_{\alpha \in \Pi} ||kv_\alpha k'v_\alpha|| / \sqrt{2}$ , where  $v_\alpha$  is a unit vector in  $V_\alpha$  with highest weight. This is also the distance induced by the embedding of  $\mathscr{P}_0$  to  $\prod_{\alpha \in \Pi} \mathbb{S}V_\alpha$ .
- $d_1(k, k') = ||k k'||$ , where  $|| \cdot ||$  is a K invariant norm on the space of  $(m+1) \times (m+1)$  square matrices  $M_{m+1}(\mathbb{R}) \supset K$ .
- $d_2(k, k')$  is the distance induced by the bi-invariant Riemannian metric on K.

We can easily verify that they are distances.

**Lemma 5.6.** The three distances  $d_0, d_1$  and  $d_2$  on  $\mathscr{P}_0$  are equivalent.

**Lemma 5.7.** Let V be an irreducible representation with good norm and with highest weight  $\chi$ , which satisfies  $\chi(H_{\alpha}) > 0$  for only one simple root  $\alpha$ . Then there exists  $t_0 > 0$  such that the following holds. Let Z be a unit vector in  $\mathfrak{k}$ , given by  $Z = \sum_{\alpha \in \mathbb{R}^+} c_{\alpha} K_{\alpha}$ . Let

$$Z_{\alpha} = \sum_{\beta \ge \alpha, \beta \in R^+} c_{\beta} K_{\beta}$$

Then for  $0 < t < t_0$ ,  $k = \exp(tZ)$  and a unit vector v with highest weight, we have

$$d(k\mathbb{R}v,\mathbb{R}v) \asymp ||kv-v|| \asymp t ||Z_{\alpha}||.$$

*Proof.* For a positive root  $\beta$ , let

$$A_{\beta} := \mathrm{d}\rho(K_{\beta})v = \mathrm{d}\rho(Y_{\beta})v.$$

Consider the representation of  $\mathfrak{s}_{\beta} = \{Y_{\beta}, X_{\beta}, H_{\beta}\} \simeq \mathfrak{sl}_2$ . Due to the classification of the representations of  $\mathfrak{sl}_2$ , the vector  $A_{\beta}$  is non zero if and only if  $\chi(H_{\beta}) > 0$ .

Fix an inner product  $(\cdot, \cdot)$  on  $\mathfrak{a}^*$  which is invariant under the Weyl group, then we can identify  $H_\beta$  with  $2\frac{\beta}{(\beta,\beta)}$ , that is

$$\chi(H_{\beta}) = (\chi, 2\frac{\beta}{(\beta, \beta)}).$$

By hypothesis,  $(\chi, \alpha) > 0$  for only one simple root  $\alpha$ , this implies that  $\chi(H_{\beta}) = 2(\chi, \beta)/(\beta, \beta) > 0$ if and only if  $\beta \ge \alpha$  and  $\beta$  is a positive root. Therefore only the vectors  $\{A_{\beta}\}_{\beta \ge \alpha, \beta \in \mathbb{R}^+}$  are non zero. They are also orthogonal since they are of different weights. When t is small enough, by Lipschitz property we conclude

$$d(k\mathbb{R}v,\mathbb{R}v) \asymp ||kv-v|| = ||\exp(tZ)v-v|| \asymp t ||d\rho(Z)v|| = t ||\sum_{\beta \ge \alpha, \beta \in R^+} c_\beta A_\beta|| \asymp t ||Z_\alpha||.$$

The proof is complete.

Proof of Lemma 5.6. First we observe that the three distances are left K invariant. It is sufficient to prove the equivalence for k' equal to the identity e.

Fix  $\epsilon$  small depending on K. Let  $B_{\epsilon}$  be the neighbourhood of e given by  $\{k \in K | d_1(k, e) < \epsilon\}$ . Then  $B_{\epsilon}^c$  is a compact subset of K. Consider the function  $f_{i,j}(k) = \frac{d_i(k,e)}{d_j(k,e)}$  for  $k \in B_{\epsilon}^c$  and  $i, j \in \{0, 1, 2\}$ . Then  $f_{i,j}$  is a positive continuous function  $B_{\epsilon}^c$ . The compactness of  $B_{\epsilon}^c$  implies that it has positive minimum on  $B_{\epsilon}^c$ . Hence there exists  $c_{i,j} > 0$  such that for k outside of  $B_{\epsilon}$ 

$$d_i(k,e) \ge c_{i,j}d_j(k,e).$$

Finally, we only need to consider a small neighbourhood of the identity. We take  $\epsilon$  small such that the exponential map at e is bi-Lipschitz. Suppose that  $k = \exp(tZ)$  with Z a unit vector in  $\mathfrak{k}$  and t > 0. Then

$$d_1(k, e) = ||e - \exp(tZ)|| \approx t = d_2(k, e).$$

Due to  $d_0(k, e) = \max_{\alpha \in \Pi} ||kv_\alpha - v_\alpha|| / \sqrt{2} \ll ||k - e|| = d_1(k, e)$ , it remains to prove that  $d_0$  is not small.

We can decompose Z as in Lemma 5.7. There exists  $\alpha \in \Pi$  such that  $||Z_{\alpha}|| \gg 1$ . By Lemma 5.7, we have

$$||kv_{\alpha} - v_{\alpha}|| \asymp t ||Z_{\alpha}|| \gg t.$$

Then we have  $d_0(k, e) \gg d_2(k, e)$ . The proof is complete.

Recall the definition of the sign function m of Section 2.5.

**Lemma 5.8.** Let  $z = kz_o, z' = k'z_o$  be two points in  $\mathcal{P}_0$ , then

$$\sqrt{2}d_0(z, z') \ge d(\pi(z), \pi(z')).$$

We have

$$\mathbf{m}(z, z') = e \iff d_0(z, z') < 1.$$

If m(z, z') = e, then

$$d(\pi(z), \pi(z')) \ge d_0(z, z').$$

*Proof.* Suppose that the angle between  $kv_{\alpha}$  and  $k'v_{\alpha}$  is  $\vartheta \in [0, \pi)$ , then  $||kv_{\alpha} - k'v_{\alpha}|| = 2\sin\frac{\vartheta}{2}$ and  $d(V_{\alpha,k\eta_o}, V_{\alpha,k'\eta_o}) = ||kv_{\alpha} \wedge k'v_{\alpha}|| = \sin\vartheta = 2\sin\frac{\vartheta}{2}\cos\frac{\vartheta}{2} \le 2\sin\frac{\vartheta}{2}$ , which implies the first inequality.

The assumption  $d_0(z, z') < 1$  is equivalent to that for every simple root  $\alpha$ , the angle  $\vartheta$  is less than  $\pi/2$ , which is equivalent to m(z, z') = e due to Lemma 2.24.

If m(z, z') = e, then for every simple root  $\alpha$ , the angle  $\vartheta$  is less than  $\pi/2$ . Hence  $\sin \vartheta = 2 \sin \frac{\vartheta}{2} \cos \frac{\vartheta}{2} \ge \sqrt{2} \sin \frac{\vartheta}{2}$ , which implies the result.

**Corollary 5.9.** The K-invariant Riemannian distance on  $\mathscr{P}$  is equivalent to the distance defined in (2.17).

Proof. By  $\mathscr{P} = \mathscr{P}_0/M$  and since the group M is a subgroup of K which preserves the distance, let  $d_2$  also be the quotient Riemannian distance on  $\mathscr{P}$ . By the same argument of the proof as in Lemma 5.6, it is sufficient to prove on a small neighbourhood of  $\eta_0$ . For any two points  $\eta$ ,  $\eta'$  in this small neighbourhood, we can find z, z' in  $\mathscr{P}_0$  such that  $\pi(z) = \eta$ ,  $\pi(z') = \eta'$  and  $d_2(z, z') = d_2(\eta, \eta')$ . Due to  $d_2(z, z')$  small, we see that  $d_0(z, z')$  is less than 1. Hence by Lemma 5.8, we have m(z, z') = e and then

$$d(\eta, \eta') \asymp d_0(z, z').$$

By Lemma 5.6, we have  $d_0(z, z') \approx d_2(z, z') = d_2(\eta, \eta')$ . The proof is complete.

# Acknowledgement

This is part of the author's Ph.D. thesis, written under the supervision of Jean-François Quint at the University of Bordeaux. The author gratefully acknowledges the many helpful suggestions and stimulating conversations of Jean-François Quint during the preparation of the paper.

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Jialun LI Université de Bordeaux jialun.li@math.u-bordeaux.fr