Kleinian Schottky groups, Patterson-Sullivan measures, and Fourier decay

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Abstract

Let $\Gamma$ be a Zariski dense Kleinian Schottky subgroup of $\text{PSL}_2(\mathbb{C})$. Let $\Lambda \subset \mathbb{C}$ be its limit set, endowed with a Patterson-Sullivan measure $\mu$ supported on $\Lambda$. We show that the Fourier transform $\hat{\mu}(\xi)$ enjoys polynomial decay as $|\xi|$ goes to infinity. This is a PSL$_2(\mathbb{C})$ version of the result of Bourgain-Dyatlov [8], and uses the decay of exponential sums based on Bourgain-Gamburd sum-product estimate on $\mathbb{C}$. These bounds on exponential sums require a delicate non-concentration hypothesis which is proved using some representation theory and regularity estimates for stationary measures of certain random walks on linear groups.

1 Introduction and main result

1.1 Fourier dimension

Let $\mu$ be a Borel probability measure on $\mathbb{R}^d$, then its Fourier transform $\hat{\mu}(\xi)$ is defined for any $\xi \in \mathbb{R}^d$ by

$$\hat{\mu}(\xi) := \int_{\mathbb{R}^d} e^{-i\langle \xi, x \rangle} d\mu(x).$$

Here $\langle \cdot, \cdot \rangle$ is the usual scalar product on $\mathbb{R}^d$ and $|\cdot|$ is the associated euclidean norm. Let $K$ be a non empty compact subset of $\mathbb{R}^d$, then following Frostman [19] its Hausdorff dimension can be expressed as

$$\dim_H(K) = \sup \left\{ s \in [0, d] : \int_{\mathbb{R}^d} |\hat{\mu}(\xi)|^2 |\xi|^{s-d} d\xi < \infty \text{ for some } \mu \in \mathcal{P}(K) \right\},$$

where $\mathcal{P}(K)$ is the set of Borel probability measures on $K$. On the other hand, the Fourier dimension is defined by

$$\dim_F(K) = \sup \left\{ s \in [0, d] : \sup_{\xi} |\hat{\mu}(\xi)|^2 |\xi|^s < \infty \text{ for some } \mu \in \mathcal{P}(K) \right\}.$$

We therefore have $\dim_F(K) \leq \dim_H(K)$, and sets for which equality occur are called Salem sets. Constructing Salem sets with genuine fractal dimension is a difficult problem, and all the known constructions either rely on the use of a random process [26, 5] or specific number theoretic properties [24, 28]. A related problem and still widely open, is to build deterministic sets with positive Fourier dimension, i.e. compact sets $K$ with fractal Hausdorff dimension for which one can find a Borel probability measure $\mu$ on $K$ whose Fourier transform has polynomial decay:

$$\hat{\mu}(\xi) = O(|\xi|^{-\epsilon}),$$

for some $\epsilon > 0$. This is of course not always possible: in dimension 1, the celebrated example of the triadic Cantor set is known to have zero Fourier dimension, by the work of Kahane and

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Salem [27]. In higher dimension, any fractal set $K$ which is contained in an affine subspace will obviously not enjoy Fourier decay as $|\xi| \to \infty$. The problem of Fourier decay of fractal measures is not only interesting for itself but also for its relationship with optics and diffraction through the Huygens-Fresnel principle: see for example [1, 39] in the physics literature. In the mathematics literature, in addition to the above mentioned works, Fourier decay is deeply connected to the problem of restriction estimates in harmonic analysis, we just mention [29] and references therein. For a comprehensive introduction to fractal sets and the calculation of Hausdorff dimension, we refer to the classic textbook of Falconer [18]. For an in depth study of the relationships between Hausdorff dimension and Fourier transform, we recommend the book of Mattila [35].

1.2 Main result

A recent result of Bourgain-Dyatlov [8] shows (in dimension 1) that all limit sets of non-elementary convex co-compact Fuchsian groups have positive Fourier decay, which is an explicit family of examples. Recent works of Sahelsten et al [25, 40] and Li [32] also prove Fourier decay in deterministic situations (Cantor sets related to the Gauss map and stationary measures on $SL_2(\mathbb{R})$). Before we state our main theorem, we need to recall some notations. From now on we take $d = 2$ and we identify $\mathbb{R}^2 \simeq \mathbb{C}$. Let $D_1, \ldots, D_r, \ldots, D_{2r}$ be $2r$ bounded open topological discs\(^1\) in $\mathbb{C}$, with $r \geq 2$, whose closures are pairwise disjoint:

$$\forall i \neq j, \quad \overline{D_j} \cap \overline{D_i} = \emptyset.$$  

Assume that we are given $\gamma_1, \ldots, \gamma_r$ in $PSL_2(\mathbb{C})$ such that

$$\forall i = 1, \ldots, r, \quad \gamma_i(\overline{C} \setminus \overline{D_{i+2r}}) = D_i.$$  

Then the free group $\Gamma := \langle \gamma_1, \ldots, \gamma_r, \gamma_1^{-1}, \ldots, \gamma_r^{-1} \rangle$ is called a Schottky group. If in addition the discs $D_i$ are genuine euclidean discs, then $\Gamma$ is called classical. The limit set $\Lambda_\Gamma$ is the complementary set of the discontinuity set $\Omega_\Gamma \subset \hat{\mathbb{C}}$ for the action of $\Gamma$ on $\hat{\mathbb{C}}$. When, using Poincaré extension, $\Gamma$ is viewed as a set of isometries of the hyperbolic 3-space $\mathbb{H}^3$, then $\Lambda_\Gamma \subset \partial \mathbb{H}^3$ coincides with the limit set of $\Gamma$ for its action on $\mathbb{H}^3$. The Hausdorff dimension of the limit set coincides with the critical exponent of Poincaré series, and is denoted throughout the paper by $\delta := \delta_\Gamma$. We point out that there is a universal upper bound strictly smaller than 2 on the dimension $\delta$ for classical Schottky groups due to Doyle [16]. On the other hand, non-classical Schottky groups are rather ubiquitous, and free subgroups of co-compact subgroups of $PSL_2(\mathbb{C})$, see L. Bowen [11], provide examples of non-classical Schottky groups with $\delta$ arbitrarily close to 2. We recall that limit sets of convex co-compact manifolds are naturally equipped with a family of finite measures called Patterson-Sullivan measures, see §2 for more details. Our main result is as follows.

**Theorem 1.1.** Assume that $\Gamma$ is a Zariski dense Schottky group in $PSL_2(\mathbb{C})$, and let $\mu$ be a Patterson-Sullivan measure on $\Lambda_\Gamma$. Fix any neighborhood $U$ of $\Lambda_\Gamma$. Let $g$ be in $C^1(U, \mathbb{C})$ and $\varphi$ be in $C^2(U, \mathbb{R})$ with

$$M := \inf_{z \in U} |\nabla_z \varphi| > 0$$

on $\Lambda_\Gamma$. Assume that $\|g\|_{C^1} + \|\varphi\|_{C^2} \leq M'$. Then there exist $C := C(M, M', \Gamma) > 0$ and $\epsilon > 0$, with $\epsilon$ depending only on $\mu$, such that for all $t \in \mathbb{R}$ with $|t| \geq 1$,

$$\left| \int_{\Lambda_\Gamma} e^{it \varphi(z)} g(z) d\mu(z) \right| \leq C|t|^{-\epsilon}. \quad (1)$$

Moreover, there exists $\alpha > 0$ such that if we have $M \geq |t|^{-\alpha}$ for all $|t|$ large, the same conclusion holds.

\(^1\)In general $\partial D_j$ is just Hölder regular and $D_j$ is not convex.
Remark 1.2.  1. In the case of \( \text{PSL}_2(\mathbb{R}) \), Theorem 1.1 is obtained by Bourgain-Dyatlov \cite[Theorem 1.2]{BD} and they show that the decay rate \( \epsilon \) depends only on the Hausdorff dimension \( \delta_\Gamma \). In our setting, the decay rate \( \epsilon \) depends on \( \delta_\Gamma \) and the regularity constant \( \kappa_8 \) given in Lemma 4.4. It is natural to expect that in higher dimensions extra quantities will appear in the characterization of the decay rate. This is because there is no uniform decay rate for any fixed \( \delta_\Gamma < 1 \). Indeed, a Zariski dense Schottky group with \( \delta_\Gamma < 1 \) can be arbitrarily close to a subgroup contained in \( \text{PSL}_2(\mathbb{R}) \), which has no such Fourier decay (see Corollary 1.4). It would be interesting to find a geometric interpretation of the regularity constant \( \kappa_8 \).

2. The decay rate \( \epsilon \) does not change if we pass to finite index subgroups, because the Patterson-Sullivan measure remains the same when passing to finite index subgroups, a result due to Roblin \cite[Lemma 2.1.4. Theorem 2.2.2]{Roblin}.

3. Let \( \Gamma \) be as in Theorem 1.1. We consider the Selberg zeta function \( Z_\Gamma(s) \) for the quotient \( \Gamma \backslash \mathbb{H}^3 \). Using the method originating in the work of Dolgopyat \cite{Dolgopyat}, Stoyanov showed that \( Z_\Gamma \) has finitely many zeros in \( \{ \Re s \geq \delta_\Gamma - \epsilon_\Gamma \} \) for some \( \epsilon_\Gamma > 0 \) depending on \( \Gamma \) \cite{Stoyanov}. Now with Theorem 1.1 available, following the exact same arguments as in \cite{BD}, we can obtain an \( \epsilon_0 > 0 \) depending only on the Fourier decay rate \( \epsilon \) given in (1) such that \( Z_\Gamma(s) \) has only finitely many zeros in \( \{ \Re(s) > \delta - \epsilon_0 \} \). In particular, this yields a uniform "essential" spectral gap for any finite cover of \( \Gamma \backslash \mathbb{H}^3 \) by the above remark. We point out that under the assumption that \( \delta_\Gamma \) is close enough to 1, this was already obtained by Dyatlov and Zahl in \cite{DyatlovZahl}.

4. In \cite{firstauthor}, the first author established Fourier decay for split semi-simple groups. To deal with the non-split group \( \text{PSL}_2(\mathbb{C}) \), we will borrow ideas from \cite{firstauthor}, while following a new scheme. We would like to point out that it seems possible to combine the methods in this paper with the ones in \cite{firstauthor} to obtain a Fourier decay for Furstenberg measures on \( \hat{\mathbb{C}} \) but we do not pursue this generalization here.

1.3 \( C^2 \)-stable positive Fourier dimension

Theorem 1.1 motivates the following definition.

**Definition 1.3.** A compact set \( K \subset \mathbb{C} \) is said to have \( C^2 \)-stable positive Fourier dimension if and only if for all \( C^2 \)-diffeomorphism \( \phi : U \to \phi(U) \subset \mathbb{C} \), defined on a neighborhood \( U \) of \( K \), \( \phi(K) \) has positive Fourier dimension.

Theorem 1.1 implies the following characterization of "stable Fourier decay" for limit sets of Schottky groups.

**Corollary 1.4.** Let \( \Gamma \) be a Schottky group as above, then \( \Lambda_\Gamma \) has \( C^2 \)-stable positive Fourier dimension if and only if \( \Gamma \) is Zariski dense in \( \text{PSL}_2(\mathbb{C}) \).

**Proof.** Assume first that \( \Gamma \) is Zariski dense, and denote by \( \phi : U \to \phi(U) \) an arbitrary \( C^2 \)-diffeomorphism with finite \( C^2 \) norm on \( U \), with \( U \supset \Lambda_\Gamma \) an open bounded set. Let \( \phi^* \mu \) be the push-forward of a Patterson-Sullivan measure \( \mu \), then

\[
\phi^* \mu(\xi) = \int_{\Lambda_\Gamma} e^{-it\phi(\xi)} d\mu(z).
\]

Set \( \xi = t\theta \) where \( t > 0 \) and \( |\theta| = 1 \), so that we have

\[
\phi^* \mu(\xi) = \int_{\Lambda_\Gamma} e^{-it\phi_\theta(z)} d\mu(z),
\]
with $\varphi_\theta(z) = \langle \theta, \phi(z) \rangle$. Notice that for all $z = x + iy \in \mathcal{U}$, we have

$$|\nabla_z \varphi_\theta|^2 = (\langle \theta, \partial_x \phi(z) \rangle)^2 + (\langle \theta, \partial_y \phi(z) \rangle)^2.$$ 

Because $\phi$ is a diffeomorphism, for all $z \in \mathcal{U}$ we get that $\partial_x \phi(z)$ and $\partial_y \phi(z)$ are linearly independent vectors, which obviously implies that $\nabla_z \varphi_\theta \neq 0$. Because $\|\varphi_\theta\|_{C^2}$ can be bounded uniformly in $\theta$, we can apply Theorem 1.1 to deduce that for all $|\xi| \geq 1$, we have

$$|\hat{\varphi}_\theta^* \mu(\xi)| \leq C|\xi|^{-\epsilon},$$

for some $\epsilon > 0$. Hence $\hat{\phi}(\Lambda_\Gamma)$ has positive Fourier dimension.

Conversely, assume that $\Gamma$ is not Zariski dense. Consider the Zariski closure $H$ of $\Gamma$ in $\text{PSL}_2(\mathbb{C})$. There is a general fact, see for example [6] and references therein, which says that a non-compact proper Lie subgroup of $H \subset \text{Isom}^+(\mathbb{H}^3)$ which has no fixed point for its action on $\partial \mathbb{H}^3$ has an invariant totally geodesic proper submanifold of $\mathbb{H}^3$. Since $\Gamma$ is taken non-elementary, $\Gamma$ must therefore leave invariant a circle for its action on $\partial \mathbb{H}^3 = \hat{\mathbb{C}}$. It is not difficult to see then that $\Lambda_\Gamma$ must be included in that invariant circle. As a consequence, the limit set of $\Gamma$ can be mapped inside the real line $\mathbb{R}$ by a Möbius map $\phi$. But clearly for any finite Borel measure $\nu$ supported on $\phi(\Lambda_\Gamma) \subset \mathbb{R}$, we have

$$\check{\nu}(\xi) = \int_{\phi(\Lambda_\Gamma)} e^{-i\langle \xi, z \rangle} d\nu(z) = \nu(\phi(\Lambda_\Gamma)) \neq 0,$$

whenever $\text{Re}(\xi) = 0$. Hence $\phi(\Lambda_\Gamma)$ has zero Fourier dimension. □

Figure 1: On the left a Zariski dense case, on the right a Fuchsian case, where $C^2$-stable positive Fourier dimension fails.

### 1.4 About the proof of Theorem 1.1

Let us now comment on the structure of the proof. After some preliminary facts and notations related to Schottky subgroups of $\text{PSL}_2(\mathbb{C})$ gathered in §2, we show in §3 how Theorem 1.1 follows from an estimate on decay of exponential sums based on Bourgain-Gamburd sum-product estimate on $\mathbb{C}$, under a non-concentration hypothesis and this generalizes the main ideas of [8]. Unfortunately, this non-concentration hypothesis cannot be verified by elementary methods as was done in the $\text{PSL}_2(\mathbb{R})$ case in [8]. We have in particular to check that this non-concentration property holds uniformly ”in every direction”, which requires the use of some more sophisticated arguments of representation theory and some regularity properties of Patterson-Sullivan measures borrowed from the work on random walks by Guivarc’h [21]. The last section is devoted to the proof of this non-concentration hypothesis which is the main difficulty of the paper. In the appendix, the first author proves that Patterson-Sullivan measures arise as
stationary measures of certain random walks on \( \text{SL}_2(\mathbb{C}) \) with \textit{finite exponential moment}, which allows us to use the key regularity property of Guivarc’h.

Verifying non-concentration hypothesis is the main challenge when trying to apply discretized sum-product estimates. For example, in the breakthrough work of Bourgain-Gamburd [9], it is precisely the non-concentration hypothesis that prevents them from obtaining a spectral gap outside of elements with algebraic entries. What’s more, in our situation, the Fourier decay is almost equivalent to the non-concentration hypothesis, because the Fourier decay will imply a spectral gap of the transfer operator, which in turn can be used to get the non-concentration hypothesis.

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\section{Preliminary estimates on Schottky groups}

In this section, we gather notations and important but elementary bounds that will be used in §3. We use similar notations as the ones introduced in the Bourgain-Dyatlov paper [8]. Recall that we are given a set of pairwise disjoint open topological discs \( D_1, \ldots, D_{2r} \) and we fix a set of generators \( \gamma_1, \ldots, \gamma_r \) in \( \text{PSL}_2(\mathbb{C}) \) such that

\[ \forall i = 1, \ldots, r, \; \gamma_i(\hat{\mathbb{C}} \setminus D_{i+r}) = D_i, \]

see the figure below where \( r = 2 \).

\begin{figure}[h]
\centering
\includegraphics[width=0.5\textwidth]{schottky_pairing.png}
\caption{A Schottky pairing.}
\end{figure}

For convenience, for \( j = r + 1, \ldots, 2r \), we set \( \gamma_j := \gamma_{j-r}^{-1} \). By the usual ping-pong argument, \( \gamma_1, \ldots, \gamma_{2r} \) generate a free group denoted by \( \Gamma \) which is convex co-compact. We will frequently use the notation

\[ D := \bigcup_{j \in \mathcal{A}} D_j, \]

where the alphabet \( \mathcal{A} \) is just the finite set

\[ \mathcal{A} := \{1, \ldots, r, r + 1, \ldots, 2r\}. \]

Let \( \Lambda_\Gamma \) be the limit set of \( \Gamma \), defined as the set of accumulation points (in \( \partial\mathbb{H}^3 = \hat{\mathbb{C}} \)) for the action of \( \Gamma \) on \( \mathbb{H}^3 \). The action of \( \Gamma \) on \( \hat{\mathbb{C}} \setminus \Lambda_\Gamma \) is proper discontinuous and \( \hat{\mathbb{C}} \setminus D \) is a fundamental domain for this action.
• For \( a \in \mathcal{A} \), we set \( \overline{a} := a + r \mod 2r \) such that \( \gamma_a = \gamma_{a}^{-1} \).

• For \( n \in \mathbb{N}_0 \), define \( \mathcal{W}_n \), the set of reduced words of length \( n \), by

\[ \mathcal{W}_n := \{ a_1 \ldots a_n \mid a_1, \ldots, a_n \in \mathcal{A}, \ a_{j+1} \neq \overline{a}_j \text{ for } j = 1, \ldots, n - 1 \}. \]

Denote by \( \mathcal{W} := \bigcup_n \mathcal{W}_n \) the set of all words. The length of a word \( a = a_1 \ldots a_n \) is denoted by \( |a| = n \).

Denote the empty word by \( \emptyset \) and put \( \mathcal{W}^\circ := \mathcal{W} \setminus \{ \emptyset \} \). For \( a = a_1 \ldots a_n \in \mathcal{W} \), put \( \overline{a} := \overline{a_n} \ldots \overline{a_1} \in \mathcal{W} \). If \( a \in \mathcal{W}^\circ \), put \( \overline{a} := a_1 \ldots a_{n-1} \in \mathcal{W} \).

• For \( a = a_1 \ldots a_n, b = b_1 \ldots b_m \in \mathcal{W} \), we write \( a \rightarrow b \) if either at least one of \( a, b \) is empty or \( a_n \neq \overline{b}_1 \). Under this condition the concatenation \( ab \) is a word.

• For \( a, b \in \mathcal{W} \), we write \( a \prec b \) if \( a \) is a prefix of \( b \), that is \( b = ac \) for some \( c \in \mathcal{W} \).

• For \( a = a_1 \ldots a_n, b = b_1 \ldots b_m \in \mathcal{W}^\circ \), we write \( a \rightsquigarrow b \) if \( a_n = b_1 \). Note that when \( a \rightsquigarrow b \), the concatenation \( ab \) is a word of length \( n + m - 1 \).

• For each \( a = a_1 \ldots a_n \in \mathcal{W} \), define the group element \( \gamma_a \in \Gamma \) by

\[ \gamma_a := \gamma_{a_1} \ldots \gamma_{a_n}. \]

Note that each element of \( \Gamma \) is equal to \( \gamma_a \) for a unique choice of \( a \) and \( \gamma_a = \gamma_{\overline{a}}^{-1}, \gamma_{ab} = \gamma_a \gamma_b \) when \( a \rightarrow b \).

• We then define the cylinder sets associated to reduced words. Given \( a = a_1 \ldots a_n \in \mathcal{W}^\circ \), we set

\[ D_a := \gamma_{\overline{a}}(D_{a_n}). \]

Remark that cylinder sets are topological discs, but may not be convex at all in the non-classical case.

• Given \( \gamma \in \Gamma \) we will often write

\[ \gamma \simeq \begin{pmatrix} a & b \\ c & d \end{pmatrix}, \]

meaning that we have:

\[ \forall z \in \mathbb{C}, \ \gamma(z) = \frac{az + b}{cz + d} \text{ with } ad - bc = 1. \]

Finally, we warn the reader about constants: throughout the rest of this paper \( C_\Gamma \) is a positive constant that depends only on \( \Gamma \) (more accurately on the choice of generators as above). This constant \( C_\Gamma \) may change from line to line, while still being denoted the same. Given \( x, y, C > 0 \), we denote by \( x \approx_C y \) the set of inequalities:

\[ C^{-1}y \leq x \leq Cy. \]

The following estimates mimic the ones that are found in [8]. However, since we do not work a priori with convex cylinder sets, all diameter estimates are replaced with measure estimates with respect to the Patterson-Sullivan measure \( \mu \), see below.
2.1 A Lipschitz property

For the rest of the paper, fix \( \epsilon_0 > 0 \) such that \( 2\epsilon_0 > \inf_{j \neq l} d(D_j, D_l) \).

**Lemma 2.1.** There exists \( C_\Gamma > 0 \) such that the following holds. For any \( a \in W^\circ \) and any \( z, w \in C \) with \( d(z, D_a) > \epsilon_0 \) and \( d(w, D_a) > \epsilon_0 \), we have

\[
|\gamma'_a z - \gamma'_a w| \leq C_\Gamma |z - w|(|\gamma'_a z||\gamma'_a w|)^{1/2}.
\] (2)

**Proof.** Suppose that \( \gamma_a = \begin{pmatrix} a & b \\ c & d \end{pmatrix} \). Then

\[
|\gamma'_a z - \gamma'_a w| = \left| \frac{1}{(cz + d)^2} - \frac{1}{(cw + d)^2} \right|
\]

\[
= \left| \frac{c(z - w)}{(cz + d)(cw + d)} \right| \left( \frac{1}{cz + d} + \frac{1}{cw + d} \right).
\]

Observe that \( \frac{c(z - w)}{cz + d} = \frac{z - w}{z + d/c} \) and \( \frac{c(z - w)}{cw + d} = \frac{z - w}{w + d/c} \). Moreover, we have \(-d/c = \gamma_a^{-1}(\infty) \in D_\Gamma \) and \( d(z, D_\Gamma), d(w, D_\Gamma) > \epsilon_0 \). These facts imply that

\[
\left| \frac{z - w}{z + d/c} \right| \leq C_\Gamma |z - w|, \quad \left| \frac{z - w}{w + d/c} \right| \leq C_\Gamma |z - w|
\]

and the inequality (2) follows. \( \square \)

We recall the following fundamental formula for Möbius transformations that will be used throughout the paper.

**Lemma 2.2.** For any \( \gamma \in \Gamma \setminus \{e\} \) and any \( x, y \in C \setminus \{\gamma^{-1}(\infty)\} \), we have

\[
|\gamma x - \gamma y| = |x - y|(|\gamma x|^1/2|\gamma y|^1/2).
\]

The proof is by straightforward computation.

2.2 Facts on Patterson-Sullivan measures

We refer the reader to [42, 43] for the introduction of Patterson-Sullivan measures. Let us recall some basic facts of Patterson-Sullivan theory which will be used in this paper. Let \( \mathbb{H}^3 \) be the upper half-space model of the hyperbolic space, given by

\[
\mathbb{H}^3 = \mathbb{C}_z \times \mathbb{R}_{y}^+,
\]

endowed with the hyperbolic metric \( g \) given by

\[
g = \frac{dzd\bar{z} + dy^2}{y^2}.
\]

We will fix a base point \( o := (0,1) \in \mathbb{H}^3 \). Let \( \Gamma \) be a convex co-compact group of isometries of \( \mathbb{H}^3 \), for example a Schottky group as defined previously. For all \( s > \delta_\Gamma \) and \( x \in \mathbb{H}^3 \), one sets

\[
\mu^s_x := \frac{1}{P_\Gamma(o, s)} \sum_{\gamma \in \Gamma} e^{-sd(x, \gamma(o))} \mathcal{D}_{\gamma o},
\]

where \( P_\Gamma(o, s) \) is the convergent Poincaré series

\[
P_\Gamma(o, s) := \sum_{\gamma \in \Gamma} e^{-sd(0, \gamma(o))},
\]

and \( \mathcal{D}_x \) is the dirac mass at \( x \). The distance \( d(x, y) \) above is with respect to the hyperbolic metric.

By taking weak limits of these measures as \( s \to \delta_\Gamma \), one obtains a family of measures supported on the limit set (called Patterson-Sullivan measures) \( \Lambda_\Gamma \) satisfying the following properties:
For all $\gamma \in \Gamma$, $\gamma^* \mu_x = \mu_{\gamma^{-1} x}$.

- For all $x, x'$, we have $\mu_{x'} = e^{-\delta B_\xi(x',x)} \mu_x$, where $B_\xi(x',x)$ is the Busemann cocycle defined by (here $\xi \in \partial H^2$)
  \[ B_\xi(x, y) = \lim_{z \to \xi} (d(x, z) - d(y, z)). \]

The Busemann cocycle is a smooth function that can be expressed in terms of Poisson kernels.

An important fact is that given an isometry $\gamma$, we have
\[ e^{-B_\xi(\gamma^{-1} o, o)} = |\gamma'(\xi)|_{S^2}, \]
where $|\gamma'(\xi)|_{S^2}$ is the derivative of $\gamma$ at $\xi$ for the spherical metric on $\mathbb{C} \cup \{\infty\}$. In particular, Patterson-Sullivan measures satisfy the equivariant formula
\[ \forall \gamma \in \Gamma, \; \gamma^* \mu_x = e^{-\delta B_\xi(\gamma^{-1} x, x)} \mu_x. \tag{3} \]

Because these measures $\mu_x$ are all absolutely continuous with respect to each other, we will focus on $\mu := \mu_o$ and refer to it as the "Patterson-Sullivan measure" on $\Lambda$. Remark that given the above definition, it is a probability measure. Under the action of $\Gamma$, we have therefore the following key formula: for all bounded Borel function $f$ on $\mathbb{C}$ and $\gamma$ in $\Gamma$
\[ \int_{\Lambda \gamma} f(z) d\mu(z) = \int_{\Lambda \gamma} f(\gamma(z))|\gamma'(z)|_{S^2}^2 d\mu(z). \tag{4} \]

The spherical metric on $\mathbb{C} \cup \{\infty\}$ can be written as $\frac{4|dz|^2}{(1+|z|^2)^2}$ for $z \in \mathbb{C} \cup \{\infty\}$. Hence,
\[ |\gamma'(z)|_{S^2} = \frac{1 + |z|^2}{1 + |\gamma|_2^2} |\gamma'(z)|. \]

For $a \in W$, we will use the notation
\[ w_a(z) := |\gamma_a'(z)|_{S^2}^2. \tag{5} \]

### 2.3 Distortion estimates for Möbius transformations

Let $\Gamma$ be a Schottky group as above. For $\gamma \simeq \begin{pmatrix} a & b \\ c & d \end{pmatrix}$, set $\|\gamma\|_E := \sqrt{a^2 + b^2 + c^2 + d^2}$ and $\|\gamma\|_S := |c|$.

**Lemma 2.3.** There exists $C_\Gamma > 0$ such that for all $\gamma$ in $\Gamma \setminus \{e\}$, we have $\|\gamma\|_S \approx C_\Gamma \|\gamma\|_E$.

**Proof.** We will use the fact that $\hat{\mathbb{C}} \setminus \mathbb{D}$ is a fundamental domain for the action of $\Gamma$ on $\hat{\mathbb{C}} \setminus \Lambda$. In particular, if we have $\gamma \in \Gamma \setminus \{e\}$ and $z \in \hat{\mathbb{C}} \setminus \overline{\mathbb{D}}$, then $\gamma(z) \in \mathbb{D}$. First start with $C_\Gamma$ to be
\[ C_\Gamma := \max_{j \in \mathcal{A}} \sup_{z \in \overline{D}_j} |z|. \]

The bound $\|\gamma\|_E \geq |c|$ is trivial. Now pick any $\gamma \simeq \begin{pmatrix} a & b \\ c & d \end{pmatrix} \in \Gamma \setminus \{e\}$. Since $\infty$ is not contained in $\overline{\mathbb{D}}$, we have therefore $\gamma(\infty), \gamma^{-1}(\infty) \in \mathbb{D}$. Hence $c \neq 0$ and
\[ C_\Gamma \geq |\gamma(\infty)| = |a/c|, \; C_\Gamma \geq |\gamma^{-1}(\infty)| = |d/c|. \]
These imply $|a|, |d| \leq C_\Gamma |c|$.
Now we can bound \(|b|\). Observe that one of the points \(\gamma(0), \gamma^{-1}(0)\) must be in \(\overline{D}\). Otherwise, we have two points \(z = \gamma(0)\) and \(w = \gamma^{-1}(0)\) outside of \(\overline{D}\), but \(\gamma^2 w = z\). This forces \(\gamma^2\) to be the identity. But \(\Gamma\) is a free group, therefore \(\gamma\) is also the identity. A contradiction. Hence

\[
either \quad C_{\Gamma} \geq |\gamma(0)| = |b/d| \quad or \quad C_{\Gamma} \geq |\gamma^{-1}(0)| = |b/a|.
\]

This yields

\[
either \quad |b| \leq C_{\Gamma}|d| \quad or \quad |b| \leq C_{\Gamma}|a|.
\]

Therefore, we have \(\|\gamma\|_E \leq C_{\Gamma}|c|\).

**Lemma 2.4.** There exists \(C_{\Gamma} > 0\) such that for all \(b \in A\), all \(x \in D_b\) and all word \(a\) with \(a \leadsto b\), we have

\[
C_{\Gamma}^{-1}\|\gamma'_a\|_S^2 \leq |\gamma'_a x| \leq C_{\Gamma}\|\gamma_a\|_S^2.
\]

**Proof.** Suppose that \(\gamma'_a \simeq \begin{pmatrix} a & b \\ c & d \end{pmatrix} \in \text{SL}_2(\mathbb{C})\). Then \(\gamma'_a x = \frac{1}{c^2(x+d/c)^2}\). As \(x \in D_b\), we have that \(|x + d/c| = |x - \gamma^{-1}_a(\infty)| \geq 1/C_{\Gamma}\). Meanwhile, we have \(x, \gamma^{-1}_a(\infty) \in D\). Hence \(|\gamma'_a x| \in [1/C_{\Gamma}, C_{\Gamma}]\|\gamma\|_S^2\).

**Lemma 2.5.** For any \(a \in \mathcal{W}_\infty\), we have

\[
C_{\Gamma}^{-1}\mu(D_a) \leq |\gamma'_a(x)|^d \leq C_{\Gamma}\mu(D_a) \quad for \quad any \quad x \in D_a.
\]

**Proof.** Due to Lemma 2.4, it suffices to show \(\mu(D_a) \approx C_{\Gamma}\|\gamma'_a\|_S^{-2\delta}\). We have

\[
\mu(D_a) = \int_{D_a} w_a'(x)d\mu(x).
\]

By Lemma 2.4, we have

\[
C_{\Gamma}^{-1}\|\gamma'_a\|_S^{-2\delta} \leq w_a'(x) \leq C_{\Gamma}\|\gamma'_a\|_S^{-2\delta} \quad on \quad D_a
\]

and Lemma 2.5 follows.

2.4 More distortion estimates

**Lemma 2.6.** We have the following contraction property: for any \(a \in \mathcal{W}_\infty\), \(b \in A\), \(a \rightarrow b\), we have

\[
\mu(D_{ab}) \leq (1 - C_{\Gamma}^{-1})\mu(D_a).
\]

**Proof.** We have

\[
\mu(D_a \setminus D_{ab}) = \int_{D_a \setminus D_{ab}} w_a'(x)d\mu(x) \geq C_{\Gamma}^{-1}\mu(D_a)\mu(D_a \setminus D_{ab}),
\]

where we use Lemma 2.5 to obtain the inequality on the right. As we have a uniform non-trivial lower bound for the measure of the sets of the form \(D_a \setminus D_{ab}\), the proof of Lemma 2.6 is complete.

**Lemma 2.7** (Parent-child ratio). For any \(a \in \mathcal{W}_\infty\), \(b \in A\), \(a \rightarrow b\), we have

\[
C_{\Gamma}^{-1}\mu(D_a) \leq \mu(D_{ab}) \leq \mu(D_a).
\]
Proof. We just need to show the lower bound. We have

\[ \mu(D_{ab}) = \int_{D_{ab}} w_a(x)d\mu(x) \geq C_Γ^{-1} \mu(D_a) \mu(D_{a,b}), \]

where we use Lemma 2.5 to obtain the inequality on the right. As we have a uniform non-trivial lower bound for the measure of the sets of the form \( \mu(D_{a,b}) \), (9) yields (8). \( \square \)

**Lemma 2.8** (Concatenation). For any \( a, b \in \mathcal{W}^0 \), \( a \prec b \), we have

\[ C_Γ^{-1} \mu(D_a) \mu(D_b) \leq \mu(D_{a,b}) \leq C_Γ \mu(D_a) \mu(D_b). \] (10)

*Proof.* This follows from Lemma 2.5 similarly to Lemma 2.7, using that \( D_{a,b} = \gamma_a^* D_b \). \( \square \)

**Lemma 2.9** (Reversal). For any \( a \in \mathcal{W}^0 \), we have

\[ C_Γ^{-1} \mu(D_a) \leq \mu(D_{\overline{a}}) \leq C_Γ \mu(D_a). \] (11)

*Proof.* Without loss of generality, we may assume that \( |a| \geq 3 \). We write \( a = a_1 \ldots a_n \) and denote \( b := a_2 \ldots a_{n-1} \), so that \( a = a_1 b a_n \). Since \( D_a = \gamma_{a_1}(D_{b a_n}) \) and \( D_a = \gamma_{\overline{a}_n}(D_{\overline{b} \overline{a}}) \), it suffices to show that

\[ C_Γ^{-1} \mu(D_{b a_n}) \leq \mu(D_{\overline{b} \overline{a}}) \leq C_Γ \mu(D_{b a_n}). \] (12)

By Lemma 2.4, we have

\[ \mu(D_{b a_n}) \approx C_Γ \|b\|_S^{2\delta} \mu(D_{a_n}) \quad \text{and} \quad \mu(D_{\overline{b} \overline{a}}) \approx C_Γ \|\overline{a}\|_S^{-2\delta} \mu(D_{a_1}). \]

It follows from the definition of \( \| \cdot \|_S \) that \( \|b\|_S = \|\overline{a}\|_S \). Hence Lemma 2.9 follows. \( \square \)

**Lemma 2.10** (Separation). For any \( a \in \mathcal{W}^0 \) and any \( b, c \in \mathcal{A} \) so that \( a \rightarrow b \) and \( a \rightarrow c \), we have

\[ \text{dist}_E(D_{ab}, D_{ac}) \geq C_Γ^{-1} \mu(D_a)^{1/\delta}. \] (13)

*Proof.* Denote \( a = a_1 \ldots a_n \). For any \( x \in D_{ab}, y \in D_{ac} \), set \( \hat{x} = \gamma_a^{-1}(x) \in D_{a_n b}, \hat{y} = \gamma_a^{-1}(y) \in D_{a_n c} \). Using Lemma 2.2 and 2.5, we obtain

\[ |x - y| \geq C_Γ^{-1} |\hat{x} - \hat{y}| \mu(D_a)^{1/\delta}. \] (14)

As we have a uniform non-trivial lower bound for the Euclidean distance between the second generation of discs, (13) follows. \( \square \)

### 2.5 Patterson-Sullivan measures II

**Lemma 2.11.** Let \( \Omega \) be any Euclidean disc of radius \( \sigma \) contained in \( D_a \) for some \( a \in \mathcal{A} \). Then

\[ \mu(\Omega) \leq C_Γ \sigma^δ. \] (15)

*Proof.* We may assume \( \#(\Omega \cap \Lambda_Γ) \geq 2 \). Let \( a \in \mathcal{W}^0 \) be the longest word such that \( \Omega \cap \Lambda_Γ \subseteq D_a \). Then there are two different \( b, c \in \mathcal{A} \) so that \( a \rightarrow b \), \( a \rightarrow c \) and \( \Omega \cap D_{ab} \neq \emptyset, \Omega \cap D_{ac} \neq \emptyset \). By Lemma 2.10, the distance between \( D_{ab} \) and \( D_{ac} \) is bounded from below by \( C_Γ^{-1} \mu(D_a)^{1/\delta} \). Hence (15) follows. \( \square \)

Armed with Section 2.4, the following two lemmas do follow directly from the arguments in [8].

**Lemma 2.12** (Lemma 2.14 in [8]). Assume that \( \tau \in (0, 1], b \in \mathcal{W}^0 \). Then

\[ \#\{a \in \mathcal{W}^0 \mid b \prec a, \mu(D_a) \geq \tau^δ\} \leq C_Γ \tau^{-δ} \mu(D_b). \] (16)

**Lemma 2.13** (Lemma 2.15 in [8]). Let \( \Omega \) be any Euclidean disc of radius \( \sigma \) contained in \( D_a \) for some \( a \in \mathcal{A} \). For all \( C_0 \geq 2 \), we have

\[ \#\{a \in \mathcal{W}^0 \mid \tau^δ \leq \mu(D_a) \leq C_0 \tau^δ, D_a \cap \Omega \neq \emptyset\} \leq C_Γ^{-1} \mu(D_a)^{1/\delta} + C_Γ \log C_0. \] (17)
2.6 Partitions and transfer operators

A partition $Z$ is a subset of words in $W^\circ$ which is such that

$$\Lambda(\Gamma) = \bigcup_{a \in Z} (\Lambda(\Gamma) \cap D_a).$$

By the definition of Schottky groups, an obvious family of partitions is given for all $n \geq 1$ by

$$Z = W_n.$$

However, this natural choice turns out to be not the most convenient for our purpose, simply because elements corresponding to words with same length may have very different distortion (derivative). Instead, similarly as in [8], we will consider $\tau > 0$ a parameter (destined to be taken small later on), and set

$$Z(\tau) := \{a \in W^\circ \mid \mu(D_a) \leq \tau^\delta < \mu(D_{a'})\}. \tag{18}$$

The fact that for all $\tau > 0$ small enough $Z(\tau)$ is a partition follows readily from Lemma 2.6 and its consequence: there exist uniform $C_\Gamma$ and $0 < \rho < 1$ such that for all $a \in W^\circ$,

$$\mu(D_a) \leq C_\Gamma \rho^{|a|}.$$

Notice that by definition of $Z(\tau)$ and using Lemma 2.7 we get that as $\tau \to 0$,

$$\tau^{-\delta} \leq \# Z(\tau) \leq C \tau^{-\delta}.$$

Moreover, we have the following estimate.

**Lemma 2.14.** For $\tau > 0, C > 1$, let

$$Z(C, \tau) = \{b \in W \mid C^{-1} \tau^\delta \leq \mu(D_b) \leq C \tau^\delta\}. \tag{19}$$

Then there exists $l \in \mathbb{N}$ independent of $\tau$ such that

$$Z(C, \tau) \subset Z(C \tau) \times \cup_{0 \leq n < 1} W_n.$$

**Proof.** For any $b \in Z(C, \tau)$, the construction of $Z(C, \tau)$, we can express $b$ as $b = b'b''$ with $b' \in Z(C \tau)$. Note that $|b''|$ is bounded by a constant depending on $C$ due to the uniform contracting property (Lemma 2.6). Then Lemma 2.14 follows. \hfill \Box

To each partition $Z(\tau)$ we will associate a transfer operator $L_{Z(\tau)}$ acting on functions which is such that for all $f$ bounded Borel on $\Lambda(\Gamma)$, we have

$$\int_{\Lambda(\Gamma)} f \, d\mu = \int_{\Lambda(\Gamma)} L_{Z(\tau)}(f) \, d\mu.$$ 

Formula (4) shows that for all $j \in A$,

$$L_{Z(\tau)}(f)(x) = \sum_{a \in Z(\tau), a \to j} w_{a'}(x) f(\gamma_{a'}(x)) \text{ if } x \in D_j.$$ 

This formula can be iterated to give

$$L_{Z(\tau)}^k(f)(x) = \sum_{a_1, \ldots, a_k} w_{a'_1 \ldots a'_k}(x) f(\gamma_{a'_1 \ldots a'_k}(x)). \tag{20}$$

The rough strategy of the proof of Theorem 1.1 is then to write

$$\int_{\Lambda(\Gamma)} e^{it\varphi(x)} g(x) \, d\mu = \int_{\Lambda(\Gamma)} L_{Z(\tau)}^k(e^{it\varphi} g) \, d\mu,$$

one hopes to catch cancellations in the exponential sums

$$\sum_{a_1, \ldots, a_k} w_{a'_1 \ldots a'_k}(x) e^{it\varphi(\gamma_{a'_1 \ldots a'_k}x)} g(\gamma_{a'_1 \ldots a'_k}x),$$

with $\tau \asymp |t|^{-\beta}$, and $\beta > 0$ suitably chosen. We will make this more precise in §3.
3 Sum-products and decay of oscillatory integrals

In this section, we prove Theorem 1.1. The key tool is an estimate on the decay of exponential sums (Proposition 3.1). In Lemma 2.1, we’ve established the Lipschitz property of the derivatives of the elements in $Γ$. Using this and the Hölder inequality, we follow the scheme in [8, Lemma 3.4, 3.5] to obtain a combinatorial description of the oscillatory integral in concern which allows us to control it via certain exponential sums. We finish the proof of Theorem 1.1 by applying Proposition 3.1 to the exponential sum in (30).

**Proposition 3.1.** Given $κ > 0$, there exist $ε > 0$ and $k ∈ \mathbb{N}$ such that the following holds for $η ∈ \mathbb{C}$ with $|η| > 1$. Let $C_0 > 0$ and let $λ_1, \ldots, λ_k$ be Borel measures supported on the annulus $\{z ∈ \mathbb{C} : 1/C_0 ≤ |z| ≤ C_0\}$ with total mass less than $C_0$. Assume that each $λ_j$ satisfies the projective non concentration property, that is,

$$∀σ ∈ [C_0|η|^{-1}, C_0^{-1}|η|^{-σ}], \quad \sup_{a ∈ \mathbb{R}, θ ∈ \mathbb{R}} |λ_j(z ∈ \mathbb{C} | |\text{Re}(e^{iθ}z) - a| ≤ σ ≤ C_0σ^κ. \quad (21)$$

Then there exists a constant $C_1$ depending only on $C_0, κ$ such that

$$\left| \int \exp(2πi\text{Re}(ηz_1 \cdots z_k))dλ_1(z_1) \cdots dλ_k(z_k) \right| ≤ C_1|η|^{-ε}. \quad (22)$$

As for the proof of the proposition, it has already been pointed out in [8] that it can be shown by following the proof of Lemma 8.43 in [7] and replacing the real version of the sum-product theorem [7, Theorem 1] by its complex version established in [10, Proposition 2]. We refer readers to [33, Appendix 4.1] for more details.

### 3.1 A combinatorial description of the oscillatory integral

We now begin the proof of Theorem 1.1. In this section $C$ is a constant depending only on the Schottky data and the constants $M, M'$ in Theorem 1.1. It may change from line to line. Let $k ∈ \mathbb{N}$ be the constant in Proposition 3.1, which depends only on $κ$, which is fixed once for all and given by Proposition 3.4. Let $t$ be the frequency parameter in (1). Without loss of generality we may assume that $|t| ≥ C$. Define the small number $τ > 0$ by

$$|t| = τ^{-2k-3/2}. \quad (23)$$

Let $Z(τ) \subset W^0$ be the partition defined in (18) and let $L_{Z(τ)}$ be the associated transfer operator, see §2.6.

We follow the notation introduced in [8]:

- for $a = a_1 \cdots a_n ∈ W^0$ and $z ∈ \mathbb{C}$, write $a ↼ z$ if $z ∈ D_{a_n};$
- for $γ = γ_a ∈ Γ$ with $a = a_1 \cdots a_n$, we write $γ → z$ or $a → z$ if $z \notin D_{γ_a};$
- we denote $A = (a_0, \ldots, a_k) ∈ Z(τ)^{k+1}$, $B = (b_1, \ldots, b_k) ∈ Z(τ)^k$;
- we write $A ⇔ B$ if and only if $a_{j-1} ↼ b_j ⇔ a_j$ for all $j = 1, \ldots, k$;
- if $A ⇔ B$, then we define the words $A * B := a'_0 b'_1 a'_1 b'_2 \cdots a'_{k-1} b'_k a'_k$ and $A\#B := a'_0 b'_1 a'_1 b'_2 \cdots a'_{k-1} b'_k$;
- denote by $b(A) ∈ A$ the last letter of $a_k$;
- for each $a ∈ W^0$, fix a point $x_a ∈ D_a;$

12
• for \( j \in \{1, \ldots, k\} \) and \( b \in Z(\tau) \) such that \( a_{j-1} \sim b \sim a_j \), define
\[
\zeta_{j,A}(b) := \tau^{-2} \gamma'_{a_{j-1}} br(x_a).
\]

Using the functions \( \varphi, g \) from the statement of Theorem 1.1, define

\[
f(x) := \exp(it\varphi(x))g(x), \quad x \in \Lambda_f.
\]

By (20), the integral in (1) can be written as follows:

\[
\int_{\Lambda_f} f d\mu = \int_{\Lambda_f} L^{2k+1}_{Z(\tau)} f d\mu = \sum_{A:B:A \leftrightarrow B} \int_{D_b(A)} f(\gamma_{A\#B}(x))w_{A\#B}(x)d\mu(x).
\]

The following lemma follows almost the same lines with [8, Lemma 3.4]. This idea is to use the Lipschitz property of \( w_{A\#B} \) (Lemma 2.1) to obtain an approximation for \( w_{A\#B}(x) \) and then use Schwartz’s inequality to get the following bound.

**Lemma 3.2.** We have

\[
\left| \int_{\Lambda_f} f d\mu \right|^2 \leq C_\tau^{-2k+1}\delta \sum_{A:B:A \leftrightarrow B} \left| \int_{D_b(A)} e^{it\varphi(\gamma_{A\#B}(x))} w_{A\#B}(x)d\mu(x) \right|^2 + C_\tau^2.
\]  

(24)

**Proof.** It follows from Lemma 2.5 that for each \( a = a_1 \ldots a_n \in Z(\tau) \), we have

\[
C^{-1}\tau^\delta \leq w_{a}(x) \leq C\tau^\delta \quad \text{for} \quad x \in D_{a_i}.
\]  

(25)

This yields, using chain rule,

\[
C^{-1}\tau^{2k\delta} \leq w_{A\#B}(\gamma_{a_k}(x)) \leq C\tau^{2k\delta},
\]

(26)

\[
C^{-1}\tau^{2k\delta} \leq w_{A\#B}(x_{a_k}) \leq C\tau^{2k\delta}.
\]

(27)

Meanwhile, using Lemma 2.1 and 2.2, we deduce that

\[
\exp(-C\tau) \leq \frac{|w_{A\#B}(\gamma_{a_k}(x))|}{w_{A\#B}(x_{a_k})} \leq \exp(C\tau).
\]  

(28)

Observe that \( |g(\gamma_{A\#B}(x)) - g(x_{a_k})| \leq C\tau \). Combining this with (25)-(28), we obtain

\[
\left| \int_{\Lambda_f} f d\mu - \sum_{A:B:A \leftrightarrow B} \int_{D_b(A)} e^{it\varphi(\gamma_{A\#B}(x))} g(x_{a_0})w_{A\#B}(x_{a_k})w_{a_k}(x)d\mu(x) \right| \leq C\tau.
\]

Using Schwarz’s inequality and (27), we get

\[
\left| \sum_{A:B:A \leftrightarrow B} \int_{D_b(A)} e^{it\varphi(\gamma_{A\#B}(x))} g(x_{a_0})w_{A\#B}(x_{a_k})w_{a_k}(x)d\mu(x) \right|^2 \leq C_\tau^{-2k+1}\delta \sum_{A:B:A \leftrightarrow B} \left| \int_{D_b(A)} e^{it\varphi(\gamma_{A\#B}(x))} w_{a_k}(x)d\mu(x) \right|^2,
\]

completing the proof of the lemma. \( \square \)
To handle the first term on the right-hand side of (24), we estimate using (25)

\[
\sum_{A,B:A \leftrightarrow B} \left| \int_{D_h(A)} e^{it\varphi(\gamma_{A\#B}(x))} w_{\alpha_k}(x) d\mu(x) \right|^2
= \sum_{A} \int_{D_h(A)}^2 w_{\alpha_k}(x) \sum_{B:A \leftrightarrow B} e^{it(\varphi(\gamma_{A\#B}(x)) - \varphi(\gamma_{A\#B}(y)))} d\mu(x) d\mu(y)
\leq C T^{2\delta} \sum_{A} \int_{D_h(A)}^2 \sum_{B:A \leftrightarrow B} e^{it(\varphi(\gamma_{A\#B}(x)) - \varphi(\gamma_{A\#B}(y)))} d\mu(x) d\mu(y).
\]

With Lemma 2.1 available, the proof of the following lemma is almost the same as in [8, Lemma 3.5].

Lemma 3.3. Denote

\[ J_r := \{ \eta \in \mathbb{C} \mid |\eta| \in [r^{-1/8}, C r^{-1/2}] \}, \]

where C is sufficiently large. Then

\[
\left| \int_{J_r} f d\mu \right|^2 \leq C T^{(2k+1)\delta} \sum_{\eta \in J_r} \left| \sum_{A:B \leftrightarrow B} e^{2\pi i \text{Re}(\eta \zeta_{1,A}(b_1) - \zeta_{1,A}(b_2))} \right| + C T^{4/4}.
\]

Proof. Fix A. Take \( x, y \in D_h(A) \) and put

\[ \tilde{x} := \gamma_{\alpha_k}(x), \quad \tilde{y} := \gamma_{\alpha_k}(y) \in D_{\alpha_k}, \quad v := \frac{\tilde{x} - \tilde{y}}{|\tilde{x} - \tilde{y}|}. \]

Assume that \( A \leftrightarrow B \). Note that \( \varphi \) is real valued function defined on a neighborhood of \( \Lambda_r \).

For \( z = x + iy \), we use the notation

\[ \varphi'(z) = \partial_x \varphi(z) - i \partial_y \varphi(z), \]

so that we can write

\[ D_{\varphi}(w) = \partial_x \varphi(z) w_1 + \partial_y \varphi(z) w_2 = \text{Re}(\varphi'(z) w), \]

for \( w = w_1 + iw_2 \). Since \( \gamma_{A\#B}(x) = \gamma_{A\#B}(\tilde{x}), \gamma_{A\#B}(y) = \gamma_{A\#B}(\tilde{y}) \), we have

\[
\varphi(\gamma_{A\#B}(x)) - \varphi(\gamma_{A\#B}(y)) = \int_0^{[\tilde{x} - \tilde{y}]} (\varphi \circ \gamma_{A\#B} \circ p)'(s) ds
= \int_0^{[\tilde{x} - \tilde{y}]} \text{Re} \left( \varphi'(\gamma_{A\#B} \circ p(s)) \gamma_{A\#B}'(p(s)) v \right) ds,
\]

where \( p : [0, |\tilde{x} - \tilde{y}|] \to \mathbb{C} \) is the path defined by \( s \mapsto \tilde{y} + vs \).

Observe that, by Lemma 2.1 and 2.2, we have

\[
|\alpha_k - p(s)| \leq C_1 r \quad \text{for any} \; s \in [0, |\tilde{x} - \tilde{y}|], \]

\[
|\gamma_{\alpha_k}\gamma_{\alpha_{k+1}} \cdots \gamma_{b_k}(p(s)) - \alpha_k| \leq C_1 r \quad \text{for any} \; 0 \leq j \leq k - 1 \; \text{and} \; s \in [0, |\tilde{x} - \tilde{y}|].
\]

These yield for any \( s \in [0, |\tilde{x} - \tilde{y}|] \)

\[
|\varphi'(x_{\alpha_k}) - \varphi'(\gamma_{A\#B}(p(s)))| \leq C r.
\]

Hence we get

\[
\left| (\varphi \circ \gamma_{A\#B} \circ p)'(s) - \tau^{2k} \text{Re} \left( \varphi'(x_{\alpha_k}) \zeta_{1,A}(b_1) \cdots \zeta_{k,A}(b_k)v \right) \right| \leq C r^{2k+1}.
\]
It follows that
\[
\left| \varphi(\gamma_{A\cdot B}(y)) - \varphi(\gamma_{A\cdot B}(x)) - \tau^{2k} \text{Re} \left( \varphi'(x_{a_0}) \zeta_{1,A}(b_1) \cdots \zeta_{k,A}(b_k)(\bar{y} - \bar{x}) \right) \right| \leq C\tau^{2k+2},
\]
where we recall that \( \zeta_{j,A}(b) := \tau^{-2} \gamma'_{a_j-1}b(x_{a_j}) \). Denote
\[
\eta := \frac{\tau^{-3/2}}{2\pi} \varphi'(x_{a_0}) \cdot (\bar{x} - \bar{y}) \cdot \text{sign}(t).
\]
Then by Lemma 2.2 and 2.5,
\[
C^{-1} M \tau^{-1/2} |x - y| \leq |\eta| \leq C \tau^{-1/2} |x - y|.
\]
Notice that we have used here that on a neighborhood of \( \Lambda_\Gamma \),
\[
M := \inf_{z \in \mathcal{U}} |\varphi'(z)| = \inf_{z \in \mathcal{U}} |\nabla_z \varphi| > 0.
\]
Recall that \( |t| = \tau^{-2k-3/2} \). By Lemma 3.2, (29) and (31), we have
\[
\left| \int_{\Lambda_\Gamma} f d\mu \right|^2 \leq C \tau^{2k+1/2} \sum_{A} \int_{\mathcal{B}(A)} \left| \sum_{B \cdot A \leftrightarrow b} e^{2\pi i \text{Re} (\eta \zeta_{1,A}(b_1) \cdots \zeta_{k,A}(b_k))} \right| d\mu(x) d\mu(y) + C\tau^{1/2}.
\]
By Lemma 2.11, for a fixed \( C_0 \)
\[
\mu \times \mu\{(x,y) \in \Lambda^2_\Gamma : |x - y| \leq C_0 \tau^{1/4}\} \leq C \tau^{5/4}.
\]
We therefore take the double integral over the set of \( x,y \) such that \( |x - y| \geq C_0 \tau^{1/4} \), which assuming that
\[
M \geq \tau^{1/8}
\]
implies for a large \( C_0 \) that \( \eta \in J_\tau \). This finishes the proof. \( \square \)

### 3.2 End of the proof of Theorem 1.1

We will apply Proposition 3.1 to suitably defined discrete measures \( \lambda_j \)'s (see below) to estimate the sum in Lemma 3.3 and hence finish the proof of Theorem 1.1. The following technical proposition verifies that these \( \lambda_j \)'s satisfy the required projective non-concentration property in (21).

For any \( a, b \in Z(\tau) \) and \( x \in \mathbb{C} \) with \( a \sim b \sim x \), write
\[
\zeta_{a,x}(b) := \tau^{-2} \gamma'_{a/b}(x).
\]

**Proposition 3.4.** Assuming that \( \Gamma \) is Zariski dense, there exist \( C > 0, \kappa > 0 \) with \( \kappa \) depending only on the Patterson-Sullivan measure \( \mu \) such that for any \( a \in Z(\tau), x \in \mathbb{C} \) and \( \sigma \in (\tau^{1/2}, 1) \), we have
\[
\sup_{a \in \mathbb{R}, \theta \in \mathbb{R}} \tau^j \# \{ b \in Z(\tau) : a \sim b \sim x, |\text{Re}(e^{i\theta} \zeta_{a,x}(b)) - a| \leq \sigma \} \leq C \sigma^\kappa.
\]

Let us show how this proposition implies the main result. For each \( A \in Z(\tau)^{k+1} \) and for \( 1 \leq j \leq k \), we define the following measure on \( \mathbb{C} \):
\[
\lambda_j(E) := \tau^j \# \{ b \in Z(\tau) : \zeta_j,A(b) \in E \}
\]
for any Borel set \( E \subset \mathbb{C} \).

Notice that by Lemma 2.5, the chain rule, and the very definition of \( Z(\tau) \), we know that the rescaled derivatives \( \zeta_{j,A}(b) \) satisfy uniformly
\[
C_0^{-1} \leq |\zeta_{j,A}(b)| \leq C_0,
\]
for some \( C_0 > 0 \), and \( C_0 \) can definitely be taken large enough so that the total mass of each \( \lambda_j \) is less than \( C_0 \). Now recall that the constant \( k \) in Proposition 3.1, is determined by \( \kappa \) from Proposition 3.4. Moreover we have:
• $|t| = \tau^{-2k-3/2}$ and $|t|$ is taken large.
• $|\eta| \in [\tau^{-1/8}, \tau^{-1/2}]$.

Therefore for each $\sigma \in [C_0|\eta|^{-1}, C^{-1}_0|\eta|^{-\epsilon}]$, we get that $\sigma \in [C_0C^{-1}_0\tau^{1/2}, C^{-1}_0\tau^{-\epsilon/8}]$. Taking again $C_0 > 0$ large enough so that $C_0C^{-1}_0 > 1$, we can make sure that $\sigma \in (\tau^{1/2}, 1)$ in order to apply Proposition 3.4. Hypothesis (21) from Proposition 3.1 is now satisfied, we can combine it with Lemma 3.3 to obtain

$$\left| \int_{\Lambda_{\Gamma}} e^{it\varphi} g d\mu \right|^2 \leq C|\eta|^{-\epsilon} + C\tau^{\delta/4} = O(\tau^{\epsilon/8}) + O(\tau^{\delta/4}) = O(|t|^{-\epsilon}),$$

with $\tilde{\epsilon} = \min\left\{ \frac{\delta}{4(2k+3/2)} : \frac{\epsilon}{8(2k+3/2)} \right\}$. The proof is done.

We point out that this ”non-concentration” result is really where the Zariski density hypothesis will be used and where our techniques deviate completely from the elementary arguments used in [8]. Section §4 is fully devoted to the proof of this Proposition 3.4.

4 Proving the non-concentration property

We prove Proposition 3.4 in this section. Here is an overview of the strategy. Roughly speaking, we want to count the elements in $Z(\tau)$ whose derivatives lie in a neighborhood of an affine line. We use the H"older’s inequality and reduce the problem to counting triples of elements whose derivatives are close to an affine line. A key observation is that the area or determinant condition enables us to obtain the desired supremum statement and hence Proposition 4.1 will lead to Proposition 3.4. In § 4.2, we discuss real polynomials (defined by (43)) which are related to the determinant in Proposition 4.1. We establish an estimate regarding the measure of small values of a real polynomial (Lemma 4.3). Real proximal representations of $\text{SL}_2(\mathbb{C})$ and Guivarc’h regularity property will naturally come into the picture. The last two subsections are about using Lemma 4.3 to obtain Proposition 4.1.

4.1 Proof of Proposition 3.4

In the rest of the paper, given two real functions $f$ and $g$, we write $f \ll g$ if there exists a constant $C_1$ only depending on $C_{\Gamma}$ such that $f \leq C_1g$. We write $f \approx g$ if $f \ll g$ and $g \ll f$. Recall that $\mu$ is the Patterson-Sullivan measure of $\Gamma$ on the extended complex plane $\hat{\mathbb{C}}$.

Proposition 3.4 follows from the following proposition.

**Proposition 4.1.** There exist $\epsilon = \epsilon(\mu) > 0$, $N > 0$ and $C > 0$ such that for any $a \in \mathcal{W}$, $\tau, \tau_1 > 0$, $1/N > \sigma > \tau, \tau_1$ and $z_0 \in \mathbb{C}$

$$\# \{ (b, c, d) \in Z(\tau)^3, e \in Z(\tau_1) \mid a \sim c \sim e \sim z_0, d \mid \det(\gamma_a b, \gamma_a c, \gamma_a d, \gamma_a e, \gamma_a z_0) \mid \leq \|\gamma_a\|^{-4\tau^2\sigma} \} \leq C\tau^{-3\delta\tau_1^{-\delta}\sigma^\epsilon},$$

(33)

where $\det(\gamma_1, \gamma_2, \gamma_3, z)$ is defined to be

$$\det(\gamma_1, \gamma_2, \gamma_3, z) = \det \begin{pmatrix} \text{Re} \gamma_1'z & \text{Im} \gamma_1'z & 1 \\ \text{Re} \gamma_2'z & \text{Im} \gamma_2'z & 1 \\ \text{Re} \gamma_3'z & \text{Im} \gamma_3'z & 1 \end{pmatrix}.$$  

(34)
One of the advantages to consider determinant is that it yields an estimate regardless of the choice of affine lines as stated in Proposition 3.4. In [8, Lemma 3.6], a weaker version of Proposition 3.4 was proved.

Note that the absolute value of (34) equals $\|u_1 \wedge u_2 + u_2 \wedge u_3 + u_3 \wedge u_1\|$ when taking $u_i = \gamma'_i z$, viewed as an element in $\mathbb{R}^2$. We need an elementary lemma in linear algebra.

**Lemma 4.2** (Corollary 3.6 in [34]). Let $X_1, X_2, X_3$ be i.i.d. random vectors in $\mathbb{R}^2$ bounded by $C > 0$. Let $l$ be an affine line in $\mathbb{R}^2$. Then for any $c > 0$, we have

$$
P\{d(X_1, l) < c\}^3 \leq \mathbb{P}\{\|X_1 \wedge X_2 + X_2 \wedge X_3 + X_3 \wedge X_1\| < 8Cc\}.$$  \hspace{1cm} (35)

A similar statement can be found in [17]. It says that if three points (in a bounded region) are near an affine line in $\mathbb{R}^2$, then the area of the triangle formed by these three points must be small, in a quantitative way. A key fact is that this bound for the area is independent of the affine line.

We show how Proposition 4.1 yields Proposition 3.4. The proof is similar to the proof of Lemma 3.9 in [34]. We use Hölder inequality: expressing $b$ as a product of $b_1$ and $b_2$ we can apply Hölder inequality, which will allow us to use Lemma 4.2.

**From Proposition 4.1 to Proposition 3.4.** Recall that $x \in D_b$ and $a \in Z(\tau)$, we want to compute

$$\#\{b \in Z(\tau) \mid a \sim b \sim b \sim x, |\text{Re}(e^{i\theta} \gamma'_a b(x)) - a| \leq \tau^2 \sigma\}. \hspace{1cm} (36)$$

Let $\tau_1 = \tau^{1/2}$. We divide $b$ into $b = b'_1 b'_2$ such that $b_1 \sim b_2$ and $b_1 \in Z(\tau_1)$. Then $b_2$ might not be in $Z(\tau_1)$, but we have a control by Lemma 2.8

$$C^{-1}_1 \mu(D_{b_1}) \mu(D_{b_2}) \leq \mu(D_b) \leq C_1 \mu(D_{b_1}) \mu(D_{b_2}). \hspace{1cm} (37)$$

Hence by (37) and (19),

$$Z(\tau) \subset Z(\tau_1) \times Z(C_1, \tau_1).$$

Then (36) is less than

$$\#\{b_1 \in Z(\tau_1), b_2 \in Z(C_1, \tau_1) \mid a \sim b_1 \sim b_2 \sim x, |\text{Re}(e^{i\theta} \gamma'_a b'_1 \gamma'_b(x) b'_2(x)) - a| \leq \tau^2 \sigma\}. \hspace{1cm} (38)$$

Fix $b_2$, the length of $\gamma'_b(x)$ is approximately $\tau_1$, due to Lemma 2.5. We consider $b_1$. Let $X_1$ be the random variable in $\mathbb{R}^2$ given by $\gamma'_a b'_1 \gamma'_b(x)$ (viewed a vector in $\mathbb{R}^2$) with $b_1$ in $Z(\tau_1)$ and $a \sim b_1 \sim b_2$. Let

$$N_{b_2} = \#\{b_1 \in Z(\tau_1) \mid a \sim b_1 \sim b_2\}.$$  

Every choice of $b_1$ has probability $1/N_{b_2}$. Then (38) equals

$$\sum_{b_2 \in Z(C_1, \tau_1), b_2 \sim x} N_{b_2} \mathbb{P}_{b_2} \left\{ \left| \text{Re} \left( \frac{\gamma'_b(x)}{\gamma'_b(x)} \right) e^{i\theta} X_1 - a \right| \leq C_1 \tau^{-3/2} \sigma \right\}. \hspace{1cm} (39)$$

The definition of $X_1$ depends on $b_2$ and we use $\mathbb{P}_{b_2}$ to emphasize. Equation (39) means the distance of $X_1$ to a line $l$ is less than $C_1 \tau^{-3/2} \sigma$. Be careful that $X_1$ and the line $l$ depend on $b_2$. The length of $X_1$ is approximately $\tau^{3/2}$. We write $e^{i\theta} := \frac{\gamma'_b(x)}{\gamma'_b(x)} e^{i\theta}$. By (35), we have

$$\mathbb{P}_{b_2} \{ |\text{Re}(e^{i\theta} X_1) - a| \leq C_1 \tau^{-3/2} \sigma \} \leq \mathbb{P}_{b_2} \{ \|X_1 \wedge X_2 + X_2 \wedge X_3 + X_3 \wedge X_1\| < C_2 \tau^{3/2} \sigma \}.$$
By the Hölder inequality, we have
\[
\frac{1}{\#Z(C_T, τ_1, x)} \sum_{b_2} N_{b_2}^3 P_{b_2} \{ |Re(e^{iθ} X_1) - a| ≤ C_T r^{3/2} σ \}
\leq \left( \frac{1}{\#Z(C_T, τ_1, x)} \sum_{b_2} N_{b_2}^3 P_{b_2} \{ ||X_1 ∧ X_2 + X_2 ∧ X_3 + X_3 ∧ X_1|| < C_T^2 r^{3} σ \} \right)^{1/3},
\]
where \( Z(C_T, τ_1, x) = \{ b ∈ Z(C_T, τ_1), b ↷ x \} \) and \( \sum_{b_2} = \sum_{b_2 ∈ Z(C_T, τ_1, x)} \). Note that
\[
\sum_{b_2} N_{b_2}^3 P_{b_2} \{ ||X_1 ∧ X_2 + X_2 ∧ X_3 + X_3 ∧ X_1|| < C_T^2 r^{3} σ \} = \# \{ (c, d, e) ∈ Z(τ_1)^3, b_2 ∈ Z(C_T, τ_1) | \]
\[
c \sim d \sim b_2 ↷ x, \ | det(γ_a e', γ_a d', γ_a e', γ_b x)| ≤ C_T^2 r^{3} σ \}. \tag{41}
\]
Using Lemma 2.14, we obtain
\[
\# \{ (c, d, e) ∈ Z(τ_1)^3, b_2 ∈ Z(C_T, τ_1) : \]
\[
c \sim d \sim b_2 ↷ x, \ | det(γ_a e', γ_a d', γ_a e', γ_b x)| ≤ C_T^2 r^{3} σ \}
\leq \#(∪_{0 ≤ n ≤ N} W_n) \cdot \# \{ (c, d, e) ∈ Z(τ_1)^3, b_2 ∈ Z(C_T τ_1) | \]
\[
c \sim d \sim b_2 ↷ x', \ | det(γ_a e', γ_a d', γ_a e', γ_b x')| ≤ C_T^2 r^{3} σ \}. \tag{42}
\]
Observe that \( r^3 ≈ ||γ_a||^{-4} (τ_1)^2 \). Now we apply (33) to estimate (42). Then combining this with (36), (38), (39), (40) and (41), we prove Proposition 3.4. \( \square \)

### 4.2 Estimate on the measure of small values of a real polynomial

We introduce the following notion. Let \( P \) be a polynomial in \( z, \bar{z} \) with complex coefficients (not necessarily homogeneous). We call \( P \) a real polynomial if
\[
P(z, \bar{z}) = \overline{P(z, \bar{z})}. \tag{43}
\]
It is worthwhile to point out that the numerator of the determinant considered in Proposition 4.1 is a real polynomial. Recall that \( μ \) is the Patterson-Sullivan measure of \( Γ \) on the extended complex plane \( \mathbb{C} \). We establish the following estimate.

**Lemma 4.3.** Fix \( n > 0 \). There exist \( C_n, κ_n > 0 \) with \( κ_n \) depending only on the regularity of the push-forward measure \( e_n^* μ \) on \( PV_n^* \) (defined in (47)) such that the following holds. Let \( P \) be a real polynomial in \( z \) and \( \bar{z} \) of highest degree \( n \). Then for \( 0 ≤ r < 1 \)
\[
μ\{ z ∈ \mathbb{C} | |P(z)| ≤ rh(P) \} ≤ C_n r^{κ_n}, \tag{44}
\]
where \( h(P) \) is the maximum of the absolute values of the coefficients of \( P \). Moreover, for \( 0 < τ < r < 1 \) and \( z ∈ \mathbb{C} \), we have
\[
\# \left\{ d ∈ Z(τ) | d ↷ z, \frac{|P(γ_a z)|}{h(P)} ≤ r \right\} ≤ C_n r^{κ_n} #Z(τ). \tag{45}
\]
4.2.1 Real proximal representations and the use of Guivarc'h regularity

For the rest of this subsection, we regard the Schottky group $\Gamma$ in concern as a group in $\text{SL}_2(\mathbb{C})$. To prove Lemma 4.3, we consider real proximal irreducible representations of $\text{SL}_2(\mathbb{C})$.

Let $n$ be any nonnegative integer. Set $W_n$ to be the complex vector space of polynomials in $u_1, u_2, \overline{u}_1, \overline{u}_2$ that are homogeneous of degree $n$ in $(u_1, u_2)$ and homogeneous of degree $n$ in $(\overline{u}_1, \overline{u}_2)$. Define a representation $\tilde{\Phi}_n$ of $\text{SL}_2(\mathbb{C})$ on $W_n$ by

$$
\tilde{\Phi}_n \begin{pmatrix} a & b \\ c & d \end{pmatrix} P \begin{pmatrix} u_1 \\ u_2 \end{pmatrix} = P \begin{pmatrix} a & b \\ c & d \end{pmatrix}^{-1} \begin{pmatrix} u_1 \\ u_2 \end{pmatrix}.
$$

(46)

This is a complex irreducible representation.

Let $V_n$ be the “real” part of $W_n$. More precisely, we let $V_n$ be the real vector space consisting of polynomials $P(u_1, u_2) \in W_n$ satisfying

$$P(u_1, u_2) = \overline{P(u_1, u_2)}.$$

Note that elements in $V_n$ can be intuitively thought as homogeneous real polynomials. We have $V_n \otimes \mathbb{R} \mathbb{C} \cong W_n$. Now define a representation $\Phi_n$ of $\text{SL}_2(\mathbb{C})$ on $V_n$ as in (46). The induced representation by $\Phi_n$ on $V_n \otimes \mathbb{R} \mathbb{C}$ is isomorphic to $\tilde{\Phi}_n$. So $\Phi_n$ is a real irreducible representation of $\text{SL}_2(\mathbb{C})$. As $\text{SL}_2(\mathbb{C})$ is Zariski connected, $\Phi_n$ is strongly irreducible. Moreover, it is proximal. This is because $\Phi_n \begin{pmatrix} e^t & 0 \\ 0 & e^{-t} \end{pmatrix}$ with $t > 0$ is a proximal matrix: $\mathbb{R} u_2^2 \overline{u}_2^2$ is the eigenspace that corresponds to the eigenvalue with the greatest absolute value.

Let $V^*_n$ be the dual space of $V_n$. Denote the dual representation of $\text{SL}_2(\mathbb{C})$ on $V^*_n$ by $\Phi_n^*$. It is also strongly irreducible and proximal. Consider the map

$$\tilde{e}_n : \mathbb{C}^2 \to V^*_n, \ (z_1, z_2) \mapsto \tilde{e}_n(z_1, z_2),$$

where $\tilde{e}_n(z_1, z_2)$ is given by $\tilde{e}_n(z_1, z_2) (P(u_1, u_2)) = P(z_1, z_2)$ for any $P(u_1, u_2) \in V_n$. It induces the following map

$$e_n : \mathbb{F}_\mathbb{C}^1 \to \mathbb{F} V^*_n$$

(47)

which is $\text{SL}_2(\mathbb{C})$-equivariant.

We fix an euclidean norm on $V^*_n$. In a finite dimensional vector space, different norms are equivalent. In particular, this norm is equivalent to the maximal norm. Note that the restriction of the maximal norm to the image of $\tilde{e}_n$ is equivalent to $|z_1|^{2n} + |z_2|^{2n}$. We take the operator norm on $V_n$. For $x = \Re v \in \mathbb{P} V_n$ and $y = \Re h \in \mathbb{P} V^*_n$, define

$$\Delta(x, y) = \frac{|h(v)|}{\|h\|\|v\|}.$$ 

Recall that since the group $\Gamma$ is Zariski dense, we can use Guivarc’h’s regularity property of $(e_n)_\mu$ (see [3, Theorem 14.1]).

Lemma 4.4. There exist $C, \kappa_n > 0$ such that for every $F \in \mathbb{P} V_n$, we have

$$(e_n)_\mu \{x \in \mathbb{P} V^*_n : \Delta(F, x) \leq r\} \leq C r^{\kappa_n}.$$

To be able to apply Theorem 14.1 from [3], we need to check that the Patterson-Sullivan measure $(e_n)_\mu$ is a Furstenberg measure (i.e. a stationary measure) arising from a random walk on $\text{GL}(V^*_n)$ with finite exponential moment. We also need to make sure that the Schottky group $\Gamma$ acts on $V^*_n$ proximally and strongly irreducibly.

For the Patterson-Sullivan measure related to a convex co-compact group, it is exactly shown in the appendix that it is a Furstenberg measure on $\hat{\mathbb{C}}$ with finite exponential moment.
arising from a random walk on $\Gamma$. Since the map (47) is $\text{SL}_2(\mathbb{C})$-equivariant, the measure $(e_n)_*\mu$ is a Furstenberg measure on $\mathbb{P}V^*_n$ arising from a random walk on $\Phi_n^*(\Gamma) < \text{GL}(V^*_n)$. For the proximal condition, we know that the subgroup $\Phi_n^*(\Gamma)$ is proximal iff its Zariski closure $\Phi_n^*(\text{SL}_2(\mathbb{C}))$ is proximal [23]. For the strongly irreducible condition, it follows from the fact that the subgroup $\Phi_n^*(\Gamma)$ is strongly irreducible iff its Zariski closure is so [23]. Furthermore, an irreducible representation of a connected group is always strongly irreducible.

4.2.2 Proof of Lemma 4.3

We show how to deduce Lemma 4.3 from Lemma 4.4.

Proof of Lemma 4.3. The idea is to express the polynomial $P$ as a linear functional on $V^*_n$. Write

$$P(z, \bar{z}) = \sum_{0 \leq j+l \leq n} c_{j,l} z^j \bar{z}^l$$

with $c_{j,l}$'s complex numbers.

Set $f(u_1, u_2) := \sum_{0 \leq j+l \leq n} c_{j,l} u_1^j u_2^{n-j} \bar{u}_1^l \bar{u}_2^{n-l}$. As $P$ is a real polynomial, we have $f \in V_n$. Take a homogeneous coordinate of $z \in \mathbb{C}$, that is, $z = \frac{z_1}{z_2}$. Then we have

$$P(z) = f(\bar{e}_n(z_1, z_2))/|z_2|^{2n}.$$

Since the support of Patterson-Sullivan measure is bounded by $C$, we have

$$|P(z)| = \frac{|f(\bar{e}_n(z_1, z_2))|}{|z_2|^{2n}} \leq C \frac{|f(\bar{e}_n(z_1, z_2))|}{\|\bar{e}_n(z_1, z_2)\|}.$$

Note that because of finite dimension, $\|f\|$ is equivalent to $h(P)$. Therefore

$$\frac{|P(z)|}{h(P)} \leq C \frac{|f(\bar{e}_n(z_1, z_2))|}{\|f\|\|\bar{e}_n(z_1, z_2)\|} = \Delta(\mathbb{P}f, e_n[z_1 : z_2]).$$

We apply Lemma 4.4 to $(e_n)_*\mu$ and hence the first statement is proved.

Now we prove the second inequality (45). The idea is to replace the counting by the Patterson-Sullivan measure. If $d$ in $Z(\tau)$ satisfies the condition in (45), then $\gamma_d^r z$ in $D_d$. Using Lemma 2.2, we have for $w$ in $D_d$

$$|P(w)| \leq |P(\gamma_d^r z)| + |P(\gamma_d^r z) - P(w)| \leq h(P)r + h(P)C_1|w - \gamma_d^r z|$$

$$\leq h(P)r + h(P)C_1\tau \leq C_1 h(P)r.$$

Note that $\mu(D_d) \geq C_1^{-1}\tau^\delta$, we can replace the counting measure by $C_1\tau^{-\delta}\mu_{D_d}$ and replace the condition $|P(\gamma_d^r z)| \leq h(P)r$ by $|P(w)| \leq C_1 h(P)r$. Hence the left hand side of (45) is less than

$$C_1 \tau^{-\delta}\mu\{w\|P(w)\| \leq C_1 h(P)r\}.$$ 

We then use (44) to finish the proof.

4.3 Proof of Proposition 4.1: non concentration of the real part

Let $C$ be a constant which depend only on $C_1$ and constants $C_n$ in Lemma 4.3 and it may vary from line to line.

We prove the non concentration of the real part using Lemma 4.3. The main point is to verify that $h(P)$ is large.
Proposition 4.5. There exists $\epsilon = \epsilon(\mu) > 0$ such that for any $0 < \tau, \tau_1 \leq \sigma \leq 1/N_0$, $z_0 \in \mathbb{C}$ and $a \in \mathcal{W}$, we have

$$\#\{ (b, c) \in Z(\tau)^2, d \in Z(\tau_1) \mid a \sim b, c \sim d \sim z_0, \quad |\text{Re}(\gamma'(a^* b^*) - \gamma'(a^* c^*))| \leq \|\gamma a\|_{\mathcal{S}}^{-2} \tau \sigma \leq C \tau^{-2} \tau_1^{-\delta} \sigma^\epsilon. \quad (48)$$

The idea is to prove that for “most” $\gamma_1, \gamma_2$ in $Z(C_\Gamma, \tau_2)$, where $\tau_2 > 0$ will be taken as $\|\gamma a\|_{\mathcal{S}}^{-2} \tau$, we have

$$\#\{ d \in Z(\tau_1) \mid d \sim z_0, |\text{Re}(\gamma_1(a^* d^*) - \gamma_2(a^* d^*))| \leq \tau_2 \sigma \} \leq \sigma^\epsilon \#Z(\tau_1).$$

We make some computation

$$\text{Re}(\gamma'_1 z - \gamma'_2 z) = \text{Re}\left(\frac{1}{(c_1 z + d_1)^2} - \frac{1}{(c_2 z + d_2)^2}\right) = \text{Re}\left(\frac{(c_2 z + d_2)^2 - (c_1 z + d_1)^2}{|c_1 z + d_1|^2 |c_2 z + d_2|^2}\right). \quad (49)$$

Let $P(z)$ be the real part of the numerator. We need to estimate $h(P)$. The following lemma is about the pairs $(b, c)$’s which yield the polynomials $P$ that might have small $h(P)$.

Lemma 4.6. Let $\sigma \geq \tau$. For each $b \in Z(\tau)$ such that $a \sim b$, we have

$$\#\{ c \in Z(\tau) \mid a \sim c, |\gamma_{ab}^{-1}(\infty) - \gamma_{ac}^{-1}(\infty)| \leq \sigma \} \leq C \tau^{-\delta} \sigma^\epsilon. \quad (50)$$

Proof. Denote $e = \overline{c}$. We have $\gamma_{ab}^{-1}(\infty) = \gamma_{ae}^{-1}(\infty) \in D_a$. Also, $C_\Gamma^{-1} \tau^\delta < \mu(D_a) < C_\Gamma^2 \tau^\delta$ by Lemma 2.7 and 2.9. Therefore, the left-hand side of (50) is bounded by

$$2r \cdot \#\{ e \in \mathcal{W} \mid C_\Gamma^{-1} \tau^\delta < \mu(D_e) < C_\Gamma^2 \tau^\delta, D_e \cap \Omega \neq \emptyset\},$$

where $\Omega$ is the disk of radius $\sigma$, centered at $\gamma_{ab}^{-1}(\infty)$. Now (50) follows from Lemma 2.13. \qed

Lemma 4.7. Let $\tau_1, \tau_2 > 0$ and $\sigma > \tau_1, \tau_2$. Let $\gamma_1, \gamma_2 \in Z(C_\Gamma, \tau_2)$ and $A_1 = \gamma_1^{-1}(\infty) - \gamma_2^{-1}(\infty)$. If $|A_1| \geq \sigma^{1/12}$, then we have

$$\#\{ d \in Z(\tau_1) \mid d \sim z_0, |\text{Re}(\gamma_1(a^* d^*) - \gamma_2(a^* d^*))| \leq \tau_2 \sigma \} \leq C \sigma^\epsilon \#Z(\tau_1), \quad (51)$$

where $\epsilon = \kappa_6/2$ in Lemma 4.3.

Proof. For $i = 1, 2$, we have $\|\gamma_i\|_{\mathcal{S}} \approx \tau_i^{-1/2}$ by Lemma 2.4 and Lemma 2.5. Observe that for $z \in D_1$, the union of disks $D_{1i}$, we have $|c_1 z + d_1| \leq C_\Gamma \|\gamma_i\|_{\mathcal{S}} \leq C_\Gamma^2 \tau_2^{-1/2}$. This implies the denominator of (49) is less than $C_\Gamma^{16} \tau_2^{-1}$ for $z \in D_1$. Now $\gamma_d z_0$ is in $D$ for $d \in Z(\tau_1), d \sim z_0$. Hence the formula (49) of $\text{Re}(\gamma_1 z - \gamma_2 z)$, it is sufficient to prove that

$$\#\{ d \in Z(\tau_1) \mid d \sim z, |\text{Re}(\gamma_d z_0)| \leq C_\Gamma^{16} \tau_2^{-3} \sigma \} \leq C \sigma^\epsilon \#Z(\tau_1). \quad (52)$$

It suffices to prove that $h(P)$ is greater than $c_1 \tau_2^{-3} \sigma^{1/2}$, where $c_1 > 0$ is a constant only depends on $C_\Gamma$, because then we can apply (45) to $P$ with $r = \sigma^{1/2} C_\Gamma^{16}/c_1$.

In order to prove $h(P)$ large, we will prove that for some choice of $z$ with bounded norm, the value $P(z)$ is large. WLOG, suppose that $|c_2| \geq |c_1|$. Take

$$z = A - \frac{d_1}{c_1},$$
where $A$ will be determined later. Then
\[
P(z) = \text{Re} \left( (c_0^2(A_1 + A)^2 - c_1^2 A) c_1^2 c_2^2 (A + A_1)^2 A^2 \right).
\]
We will take $|A|$ to get rid of the minus and $|A|$ not too small to have a lower bound. Take $|A| = |A_1|/10$, then the angle of the above formula almost only depends on $A$. With a suitable choice of the angle of $A$, the value of $P(z)$ is almost the absolute value. As $|c_1| = \|\gamma_i\|_S \approx \tau_2^{-1/2}$, we obtain
\[
P(z) > |c_0^2 A_1|^2 \|c_1^2 c_2^2 A_1\| > \tau_2^{-3} |A_1|^6 > \tau_2^{-3}\sigma^{1/2}.
\]
Since the norm of $z$ is bounded by $C_\Gamma$, we see that $h(P) > \tau_2^{-3} \sigma^{1/2}$. The proof is complete. \qed

4.4 Proof of Proposition 4.1: uniform non-concentration

We complete the proof of Proposition 4.1 in this subsection. The idea is the same with Proposition 4.5: we find conditions on $\gamma_1, \gamma_2, \gamma_3$ such that the real polynomial $P(z)$ showing up in the determinant has reasonable large height $h(P)$.

The following lemma is similar to Lemma 4.7, but much more involved.

**Lemma 4.8.** Let $\tau_1, \tau_2 > 0, \sigma > \tau_1, \tau_2$ and $z_0 \in \mathbb{C}$. Let $\gamma_i \in Z(C_\Gamma, \tau_2)$, $i = 1, 2, 3$. If $A_1 = \gamma_3^{-1}\infty - \gamma_1^{-1}\infty$, $A_2 = \gamma_3^{-1}\infty - \gamma_2^{-1}\infty$ satisfy
\[
|A_1|, |A_2| \geq \sigma^{1/128}
\]
and for $z_3 = \gamma_3^{-1}\infty$,
\[
|\text{Re}(\gamma'_1 z_3) - \text{Re}(\gamma'_2 z_3)| \geq \tau_2 \sigma^{1/8},
\]
then
\[
\# \{ e \in Z(\tau_1) \mid \gamma_i \to e \leadsto z_0, |\det(\gamma_1, \gamma_2, \gamma_3, \gamma' e, z_0)| \leq \tau_2^2 \sigma \} \leq C \sigma^\epsilon Z(\tau_1),
\]
where $\epsilon = \kappa_8/4$ in Lemma 4.3.

We first show how the above lemma will lead to Proposition 4.1.

**Proof of Proposition 4.1.** Let $\gamma_1 = \gamma_{a' b'}, \gamma_2 = \gamma_{a' c'}, \gamma_3 = \gamma_{a' d'}$. By Lemma 2.8, they are in $Z(C_\Gamma, \tau_2)$ with $\tau_2 = \|\gamma_a\|^{-2}\tau$.

We have a dichotomy. If $\gamma_1, \gamma_2, \gamma_3$ satisfy the conditions in Lemma 4.8, then the number of $e$ is small, less than $\#Z(\tau_1)\sigma^\epsilon$.

If not, the condition on $A_1, A_2$ can be dealt with Lemma 4.6. That is the number of $b, c, d$ not satisfying (53) is small by Lemma 4.6.

The main difficulty is to verify (54), but this can be dealt with Proposition 4.5. Because for $\gamma_3 = \gamma_{a' d'}$, then
\[
z_3 = \gamma_{a' d'}^{-1}\infty = \gamma_{a' d'} \infty = \gamma_{d'}(\gamma_a\infty),
\]
where $d'[a'] = d_{n-1} \cdots d_1 a_{m-1} \cdots a_1 = d_{n-1} \cdots d_2 a_m \cdots a_1 = (d')'a$ and we have $d' \leadsto a$. Let $f = d'$. The element $f$ is not always in $Z(\tau)$. By Lemma 2.9, it is in $Z(C_\Gamma, \tau)$. The number of $b, c, d$ not satisfying (54) is less than
\[
\# \{ (b, c) \in Z(\tau), f \in Z(C_\Gamma, \tau) \mid a \leadsto b, c, f \leadsto z, |\text{Re}(\gamma_{a' b'}(f'z)) - \text{Re}(\gamma_{a' c'}(f'z))| \leq \tau_2 \sigma^{1/8} \},
\]
where $z = \gamma_a\infty$ and $\tau_2 = \|\gamma_a\|^{-2}\tau$. Then by Lemma 2.14 and Proposition 4.5, the proof is complete. \qed

It remains to prove Lemma 4.8.
Proof of Lemma 4.8. By the same argument as in the proof of Lemma 4.7, we have an upper bound of the denominator of $\det(\gamma_1, \gamma_2, \gamma_3, z)$ with $z \in D$, which is less than $C_1^{24} \tau_2^{-6}$. Therefore it is enough to prove that
\[
\# \{ e \in Z(\tau_1) \mid e \sim z_0, |P(\gamma_0, z_0)| \leq C_1^{24} \tau_2^{-4} \sigma \} \leq C \# Z(\tau_1),
\]
where the polynomial $P(z)$ is the numerator of $\det(\gamma_1, \gamma_2, \gamma_3, z)$, given by
\[
P(z) = \det \begin{pmatrix}
\text{Re}(c_1z + d_1)^2 & \text{Im}(c_1z + d_1)^2 & |c_1z + d_1|^4 \\
\text{Re}(c_2z + d_2)^2 & \text{Im}(c_2z + d_2)^2 & |c_2z + d_2|^4 \\
\text{Re}(c_3z + d_3)^2 & \text{Im}(c_3z + d_3)^2 & |c_3z + d_3|^4
\end{pmatrix}.
\]

It suffices to prove that $h(P)$ is greater than $c_1 \tau_2^{-4} \sigma^{3/4}$, where $c_1 > 0$ is a constant depending on $C_1$, because then we can apply Lemma 4.3 (45) to $P$ with $r = \sigma^{1/4} C_1^{24} / c_1$.

In order to prove $h(P)$ large, we will prove that for some choice of $z$ with bounded norm, the value $P(z)$ is large. We take $z = A - \frac{d_3}{c_3} = A + z_3$, where $A$ will be determined later. Then
\[
P(z) = -\det \begin{pmatrix}
\text{Re}(c_1(A_1 + A))^2 & \text{Im}(c_1(A_1 + A))^2 & |c_1(A_1 + A)|^4 \\
\text{Re}(c_2(A_2 + A))^2 & \text{Im}(c_2(A_2 + A))^2 & |c_2(A_2 + A)|^4 \\
\text{Re}(c_3A)^2 & \text{Im}(c_3A)^2 & |c_3A|^4
\end{pmatrix}.
\]

(56)

We first fix the angle of $A$ such that $(c_3A)^2$ is an imaginary number, that is
\[
\text{Re}(c_3A)^2 = 0.
\]

(57)

We let $|A| = \sigma^{1/4}$ and we claim that $|P(z)| \gg \tau_2^{-4} \sigma^{3/4}$.

Now, we expand the determinant (56) with respect to the last line, using (57), which gives
\[
P(z) = P_1 + P_2,
\]

(58)

with
\[
P_1 := \text{Im}(c_3A)^2 \text{Re} \left( \gamma_1'z - \gamma_2'z \right) |\gamma_1'|^{-2} |\gamma_2'|^{-2}
\]

and
\[
P_2 := |c_3A|^4 \left( \text{Re}(c_1(A_1 + A))^2 \text{Im}(c_2(A_2 + A))^2 - \text{Re}(c_2(A_2 + A))^2 \text{Im}(c_1(A_1 + A))^2 \right).
\]

Due to $\gamma_i \in Z(C_1, \tau_2)$, by Lemma 2.3 we know that $|c_1|, |c_2| \ll \tau_2^{-1/2}$. By $|A_1 + A|, |A_2 + A| \ll 1$, we obtain
\[
|P_2| \ll |c_3A|^4 \tau_2^{-2}.
\]

(59)

Let
\[
B(z) = \text{Re}(\gamma_1'z - \gamma_2'z) |\gamma_1'|^{-2} |\gamma_2'|^{-2},
\]

then $P_1 = \text{Im}(c_3A)^2 B(z)$. The coefficients of $B$ are monomials of degree 6 on $c_1, d_1, c_2, d_2$. By Lemma 2.3, we obtain
\[
h(B) \leq \sup\{ \|\gamma_1\|_S, \|\gamma_2\|_S \}^6 \ll C \tau_2^{-3}.
\]

Due to $|A| = \sigma^{1/4}$, we know that
\[
|B(z) - B(z_3)| \ll |z - z_3|h(B) \ll |z - z_3|\tau_2^{-3} = |A|\tau_2^{-3} = \sigma^{1/4} \tau_2^{-3}.
\]

(60)

Thanks to $|A_1|, |A_2| \geq \sigma^{1/128}$ (53), we obtain $|\gamma_jz_3|^{-1} \geq |c_j|^2 |A_j|^2 \gg \tau_2^{-2} \sigma^{2/128}$ for $j = 1, 2$. Combining with (60) and (54), we obtain
\[
|B(z)| \geq |B(z_3)| - |B(z) - B(z_3)| \geq \sigma^2 \tau_2^{1/8} \tau_2^{-3} \sigma^{8/128} - \sigma^{-\alpha} \sigma^{1/4} \tau_2^{-3} \gg \tau_2^{-3} \sigma^{1/4}.
\]

(61)
measure on \( \partial \Gamma \), \( \tau \) and the Patterson-Sullivan measures.

This is what we need, then the bound \(|z| \leq C\Gamma\) implies \( h(P) \gg \tau_2^{-4} \sigma^{3/4} \).

Appendix A. Exponential moment and Stationarity of Patterson-Sullivan measures.

Jialun LI

In this appendix, we give a construction of a random walk on a convex cocompact subgroup of \( SO_0(1,n) \), which has exponential moment and such that the associated Patterson-Sullivan measure of the convex cocompact group can be realized as a stationary measure.

Let \( \Gamma \) be a convex cocompact subgroup of \( G = SO_0(1,n) \) \((n \geq 2)\) and \( \mu \) be an associated Patterson-Sullivan measure on the boundary at infinity \( \partial \mathbb{H}^n \). Let \( \nu \) be a Borel probability measure on \( G \). We call \( \mu \) a \( \nu \)-stationary measure or a Furstenberg measure if

\[
\mu = \nu \ast \mu := \int_G \gamma \ast \mu \, d\nu(g).
\]

In this appendix, we provide a construction of a measure \( \nu \) on \( \Gamma \) such that \( \mu \) is \( \nu \)-stationary and \( \nu \) has a finite exponential moment, that is there exists \( \epsilon > 0 \) such that \( \int_G \|\gamma\|^\alpha \, d\nu(\gamma) < \infty \).

For a measure \( \nu \) on \( G \), we let \( \Gamma_\nu \) be the subgroup generated by the support of \( \nu \).

Theorem A.1. Let \( \Gamma \) be a convex cocompact subgroup of \( G \), and let \( \mu \) be the Patterson-Sullivan measure on the limit set \( \Lambda_\Gamma \). Then there exists a probability measure \( \nu \) on \( \Gamma \) with a finite exponential moment such that \( \mu \) is \( \nu \)-stationary and \( \Gamma_\nu = \Gamma \).

Remark A.2. 1. In [30] and [31], Lalley announced the existence of such a \( \nu \) for Schottky groups. But Lalley’s proof only works for Schottky semigroups. In [12], the authors proved the existence of such a \( \nu \) without the moment condition in the geometrically finite case. Our construction combines the methods of Connell-Muchnik and Lalley.

2. For cocompact lattices, the construction is due to Furstenberg [20]. When the Hausdorff dimension \( \delta \geq (n - 1)/2 \), the construction is due to Sullivan [43]. Their methods are based on the discretization of Brownian motions on hyperbolic spaces.

3. On the other hand, for the geometrically finite with cusps case, it is impossible to find such a measure \( \nu \) with exponential moment, because the finite exponential moment condition is impossible for noncompact lattice \( \Gamma \) in \( SL(2, \mathbb{R}) \). It is shown that if \( \nu \) is a measure on \( \Gamma \) with a finite first moment, then the \( \nu \)-stationary measure \( \mu \) is singular with respect to the Lebesgue measure (This fact is due to Guivarc’h and Le Jan [22]. See also [15] and [4]).

4. The result for \( SO_0(1,2) \) has already been announced and used in [32].

A.1 Basic properties and cover

We will use the ball model for the hyperbolic \( n \)-space and fix the origin point \( o \) in \( X = \mathbb{B}^n = \{ x \in \mathbb{R}^n | ||x||^2 = x_1^2 + \cdots + x_n^2 < 1 \} \). The hyperbolic riemannian metric \( d \) at \( x \) in \( X \) is given by

\[
\frac{4dx^2}{(1-\|x\|^2)^2}.
\]
The infinity $\partial X$ is isomorphic to the sphere $S^{n-1}$. For $x, y$ in $X$ and $C > 0$, let $\mathcal{O}_C(x, y)$ be the shadow of a ball centred at $y$ of radius $C$ seen from $x$, that is, the set of $\xi$ in the boundary $\partial X$ such that the geodesic ray issued from $x$ with limit point $\xi$ intersects the ball $B(y, C)$.

Let $\text{hull}(\Lambda_\Gamma)$ be the convex hull of the limit set $\Lambda_\Gamma$ in $X \cup \partial X$. Without loss of generality, we suppose that $o$ is in the convex hull. Since the group $\Gamma$ is convex cocompact, the quotient $C(\Gamma) = \Gamma \backslash (\text{hull}(\Lambda_\Gamma) \cap X)$ is compact. Let

$$C_0 = 6 \max\{ \text{the diameter of the quotient set } C(\Gamma), 2C_1, 2, \log C_2 \},$$

where $C_1, C_2$ are defined in Lemma A.3 and Lemma A.7 respectively.

For an element $\gamma \in G$, we write $x^\gamma_o$ for the intersection of the ray $o, \gamma^{-1}o$ with the boundary $\partial X$. For $\gamma \in \Gamma$, let $\kappa(\gamma) = d(o, \gamma o)$ and let $r_\gamma = e^{-\kappa(\gamma)}$. Set $B_o = \mathcal{O}_{C_0}(o, \gamma^{-1}o)$.

Recall the Busemann function and the Patterson-Sullivan measure. Recall (3), that is for $\gamma$ in $\Gamma$ and $\xi \in \partial X$, we have

$$\frac{d\gamma_\star \mu}{d\mu}(\xi) = e^{-\delta B_\xi(\gamma^{-1}o, o)}.$$ (63)

Let $f_\gamma(\xi) = e^{-\delta B_\xi(\gamma^{-1}o, o)}$. For the stationary equation $\sum_{\gamma \in \Gamma} \nu(\gamma) \gamma_\star \mu = \mu$, it is sufficient to verify

$$\sum_{\gamma \in \Gamma} \nu(\gamma) f_\gamma = 1 \text{ on } \Lambda_\Gamma.$$ (64)

Now we start to establish some properties of $f_\gamma$.

Recall that for two real functions $f$ and $g$, we write $f \ll g$ if there exists a constant $C > 0$ only depending on the group $\Gamma$ such that $f \leq Cg$. We write $f \approx g$ if $f \ll g \ll f$.

**Lemma A.3** (Sullivan). There exists $C_1 > 0$ such that the following holds. For any $C \geq C_1$ there exists $C'$ such that for all $\gamma$ in $\Gamma$

$$\frac{1}{C} r_\gamma^\delta \leq \mu(\mathcal{O}_C(o, \gamma^{-1}o)) \leq C' r_\gamma^\delta.$$ For the proof please see [37, Page 10].

**Lemma A.4** (Triangle rule). Let $ABC$ be a geodesic triangle in $X$. Let $\alpha, \beta, \gamma$ be the three angles of $A, B, C$ and let $a, b, c$ be the length of $BC, CA, AB$. Then

$$\frac{\sin \alpha}{\sinh a} = \frac{\sin \beta}{\sinh b} = \frac{\sin \gamma}{\sinh c}.$$ See for example [2, Page 148]. For $\xi$ in $\partial X$ and $t \in \mathbb{R}_+$, let $\xi_t$ be the point in the geodesic ray $o\xi$ with distance $t$ to $o$. We define distances on the boundary, that is the visual distance: for $x \in X$ and $\xi, \xi'$ on $\partial X$

$$d_x(\xi, \xi') = \lim_{t \to +\infty} e^{\frac{1}{2}(-d(x, \xi_t)-d(x, \xi'_t)+d(\xi_t, \xi'_t))}.$$ We fix the visual distance at the origin $o$, i.e. $d_o$, on the boundary $\partial X$.

**Lemma A.5.** The distance $d_o$ is the sinuous of the angle, that is for $\xi, \xi'$ in $\partial X$ we have $d_o(\xi, \xi') = \sin \frac{1}{2} \angle \xi o \xi'$. The shadow $B_o$ is of radius $\frac{\sinh C_0}{\sinh \kappa(\gamma)}$.

**Proof.** Let $p$ be the midpoint of $\xi_t, \xi'_t$. Then $op$ is orthogonal to the geodesic $\xi_t \xi'_t$. Let $\theta$ be the half of the angle $\xi o \xi'$. Hence by triangle rule

$$\frac{\sin \pi/2}{\sinh d(o, \xi_t)} = \frac{\sin \theta}{\sinh(d(\xi_t, \xi'_t)/2)}.$$
Therefore
\[ \sin \theta = \frac{\sinh(d(\xi, \xi')/2)}{\sinh d(o, \xi)}. \]
The function \( \sinh s \) is almost \( e^s/2 \) when \( s \) is large. When \( t \) tends to infinite, we obtain
\[ \sin \theta = \lim_{t \to +\infty} e^{d(\xi, \xi')/2 - d(o, \xi)} = d_o(\xi, \xi'). \]

Let \( q \) be the tangent point of the ball \( B(\gamma^{-1}o, C_0) \) with a geodesic ray starting from \( o \). Then by triangle rule
\[ \sin \angle qo(\gamma^{-1}o) = \sinh C_0/ \sinh d(o, \gamma^{-1}o). \]
The proof is complete.

We need a Lipschitz property of the Busemann function as a function on \( \partial X \).

Lemma A.6. For \( \xi, \xi' \) in \( \partial X \) and \( x, y \) in \( X \) with \( d_x(\xi, \xi') \leq e^{-d(x,y)-2} \), we have
\[ |B_{\xi}(x, y) - B_{\xi'}(x, y)| \leq 32e^{d(x,y)/2} d_x(\xi, \xi')^{1/2}. \]

This lemma is implicitly contained in the proof of [36, Proposition 3.5]. We summarize the properties of \( f_{\gamma} \) in the following lemma

Lemma A.7. Let \( \gamma \) be an element in \( \Gamma \).
1) Let \( \eta \) be a point in \( B_{\gamma} \). Then we have
\[ f_{\gamma}(\eta) \leq r_{\gamma}^{-\delta} = e^{2C_0 \delta} f_{\gamma}(\eta). \] (65)
2) There exists \( C_2 > 0 \). If \( \xi, \eta \in \partial X \) satisfy \( d_o(\xi, \eta) \leq r_{\gamma}/e^2 \), then we have
\[ |f_{\gamma}(\xi)/f_{\gamma}(\eta) - 1| \leq C_2 d_o(\xi, \eta)^{1/2} r_{\gamma}^{-1/2}. \] (66)
3) Let \( \xi \) be a point in \( \partial X \). Then we have
\[ f_{\gamma}(\xi) \leq r_{\gamma}^{\delta} d_o(\xi, x_{\gamma}^m)^{-2\delta}. \] (67)

Proof. For (65), by definition of \( x_{\gamma}^m \), we have
\[ f_{\gamma}(x_{\gamma}^m) = e^{-\delta B_{\gamma}^{\eta}(\gamma^{-1}o, o)} = e^{\delta d(\gamma^{-1}o, o)} = r_{\gamma}^{-\delta}. \]
For \( \eta \) in \( B_{\gamma} \), by applying triangle inequality, we obtain (65).

For the second statement, by Lemma A.6, we get
\[ |B_{\xi}(\gamma^{-1}o, o) - B_{\eta}(\gamma^{-1}o, o)| \leq 32d_o(\xi, \eta)^{1/2} r_{\gamma}^{-1/2} \leq 32/e. \]
Due to \( |e^t - 1| \ll |t| \) for \( |t| \leq 32/e \), we obtain
\[ |f_{\gamma}(\xi)/f_{\gamma}(\eta) - 1| = |e^{-\delta(B_{\xi}(\gamma^{-1}o, o) - B_{\eta}(\gamma^{-1}o, o))} - 1| \ll d_o(\xi, \eta)^{1/2} r_{\gamma}^{-1/2}. \]

For the third statement, by definition
\[ -B_{\xi}(\gamma^{-1}o, o) + d(\gamma^{-1}o, o) = \lim_{z \to \xi} -d(z, \gamma^{-1}o) + d(z, o) + d(\gamma^{-1}o, o) \]
\[ \leq \lim_{z \to \xi, w \to x_{\gamma}^m} -d(z, w) + d(z, o) + d(w, o) = -2 \log d_o(\xi, x_{\gamma}^m). \]
The proof is complete.

\[ \square \]
For any $n$ in $\mathbb{N}$, let $r_n = e^{-4C_0n}$. We want to construct a cover of $\Lambda_\Gamma$. Let $S_n$ be the set of all $\gamma$ that satisfy

$$e^{-2C_0r_n} \leq r_\gamma < r_n. \quad (68)$$

**Lemma A.8.** For any $n$ in $\mathbb{N}_0$, the family $\{B_\gamma\}_{\gamma \in S_n}$ consists of balls which cover $\Lambda_\Gamma$ with bounded Lebesgue number $C_3$, that is any $\xi \in \Lambda_\Gamma$ is contained in at most $C_3$ balls.

**Proof.** Let $\xi$ be a point in the limit set $\Lambda_\Gamma$, then the ray $o\xi$ is in the convex hull $hull(\Lambda_\Gamma)$. Consider the point $p_n$ in the ray such that $d(p_n, o) = |\log r_n| + C_0$. Since the diameter of $C(\Gamma)$ is less than $C_0$ (62), there exists $\gamma$ in $\Gamma$ such that

$$d(p_n, \gamma^{-1}o) \leq C_0. \quad (69)$$

Hence $d(\gamma^{-1}o, o) \in [|\log r_n|, |\log r_n| + 2C_0]$, which implies $\gamma \in S_n$. The inequality (69) also implies that the distance from $\gamma^{-1}o$ to the ray $o\xi$ is less than $C_0$, i.e. $d(\gamma^{-1}o, o\xi) \leq C_0$. By the definition of shadow, we obtain $\xi \in B_\gamma$. The family $\{B_\gamma\}_{\gamma \in S_n}$ is a cover of the limit set $\Lambda_\Gamma$.

It remains to prove that each point $\xi \in \Lambda_\Gamma$ is covered by a bounded number of balls. Let $q_n, q_n'$ be two points in the ray $o\xi$ with $d(q_n, o) = |\log r_n| - C_0$ and $d(q_n', o) = |\log r_n| + 3C_0$. Let $J$ be the geodesic segment connecting $q_n$ and $q_n'$. Let

$$S_n(\xi) = \{\gamma \in S_n : \xi \in B_\gamma\}.$$ 

Due to $\gamma \in S_n(\xi)$ and the definition of shadow, we obtain $d(\gamma^{-1}o, J) \leq C_0$, that is $\gamma^{-1}o$ are in $JC_0$, the $C_0$ neighbourhood of $J$. The group $\Gamma$ is discrete without torsion, there exists $c > 0$ such that $\min_{\gamma \neq e} d(o, \gamma^{-1}o) > c$. Then the set $\{\gamma^{-1}o\}_{\gamma \in S_n(\xi)}$ is a discrete set in $JC_0$ and is $c$ separated, that is any two different points has distance greater than $c$. The volume of $JC_0$ is uniformly bounded. Hence there is upper bound of the number of elements in $S_n(\xi)$. The proof is complete.

**Remark A.9.** This is a key lemma where we need the hypothesis of convex cocompactness. When $\Gamma$ is a Schottky subgroup, the construction is easier. We can find a cover of the limit set $\Lambda_\Gamma$ with no overlap.

A.2 Properties of operator $P_n$

We will use the family of covers $S_n$ and basic properties of $f_\gamma$ to give properties of operator $P_n$, which will be constructed later.

Let $C_4 = 4C_0$ and let $\beta$, $\epsilon$ and $C_5$ be positive numbers defined subsequently such that

$$1 - \beta > \beta + e^{-C_4}, \quad r_n^* = (1 - \beta)^n, \quad C_5 = \frac{2e^{C_4}}{1 - \beta}. \quad (70)$$

For every $\gamma$ in $\Gamma$, let $B_\gamma' = \mathcal{O}_{C_1}(o, \gamma^{-1}o)$. By Lemma A.3, we know that $B_\gamma' \cap \Lambda_\Gamma \neq \emptyset$. We fix $\eta_\gamma$ in the set $B_\gamma' \cap \Lambda_\Gamma$ for each $\gamma$. Let $R$ be a continuous function on $\partial X$, which is positive on $\Lambda(\Gamma)$. For any $n$ in $\mathbb{N}$ and $\eta \in \partial X$, we define

$$P_nR(\eta) = \sum_{\gamma \in S_n} R(\eta_\gamma) r_\gamma^* f_\gamma(\eta). \quad (70)$$

This construction of $P_nR$ will inherit the Lipschitz property of $f_\gamma$.

**Lemma A.10.** For any $n$ in $\mathbb{N}_0$, if $\xi, \eta \in \partial X$ satisfy $d_o(\xi, \eta) \leq r_{n+1}$, then

$$\left| \frac{P_nR(\xi)}{P_nR(\eta)} - 1 \right| \leq (d_o(\xi, \eta)/r_{n+1})^{1/2}. \quad (70)$$

27
Proof. For $\gamma \in S_n$, by (68) we have $d_o(\xi, \eta) \leq r_{n+1} = e^{-4C_0r_n} \leq e^{-2C_0r_{\gamma}} \leq r_{\gamma}/e^2$. By (68) and (62), we get $r_{n+1} \leq r_{\gamma}e^{-2C_0} \leq r_{\gamma}/C_2^2$, then we use (66) to obtain
\[ \left| \frac{f_{\gamma}(\xi)}{f_{\gamma}(\eta)} - 1 \right| \leq C_2(d_o(\xi, \eta)/r_{\gamma})^{1/2} \leq (d_o(\xi, \eta)/r_{n+1})^{1/2}. \]
As $P_nR$ is a positive linear combination of $f_\gamma$ with $\gamma \in S_n$, we have $|P_nR(\xi)/P_nR(\eta) - 1| \leq (d_o(\xi, \eta)/r_{n+1})^{1/2}$.

Lemma A.11. For any $n$ in $\mathbb{N}$, if a positive function $R$ on $\Lambda(\Gamma)$ satisfies the following condition:
(1) For $\xi, \eta$ in $\Lambda_\Gamma$, if $d_o(\xi, \eta) \leq r_{n+1}$, then we have
\[ |R(\xi)/R(\eta) - 1| \leq (d_o(\xi, \eta)/r_{n+1})^{1/2}. \]  
(2) For $\xi, \eta$ in $\Lambda_\Gamma$, if $d_o(\xi, \eta) > r_{n+1}$, then
\[ |R(\xi)/R(\eta)| \leq C_5d_o(\xi, \eta)^{1/2}/(1 - \beta)^n. \]
Then there exist $C_6, C_7$ independent of $n, R$ such that for all $\eta \in \Lambda_\Gamma$
\[ R(\eta)/C_6 \leq P_{n+1}R(\eta) \leq C_7R(\eta). \]  
Proof. Since $\{B_\gamma\}_{\gamma \in S_{n+1}}$ is a cover of $\Lambda_\Gamma$, there is a $\gamma \in S_{n+1}$ such that $\eta \in B_\gamma$. By definition $P_{n+1}R(\eta) \geq R(\eta)/f_{\gamma}(\eta)r_{\gamma}^\delta$. Thanks to $r_{\gamma} \leq r_{n+1} = e^{-4C_0n+1} \leq e^{-4}$, we get $\sinh(\kappa(\gamma)) \geq r_{\gamma}^{-1/4}$. Due to Lemma A.5 and (68), we obtain $B_\gamma$ is of radius $\sinh C_0/\sinh(\kappa(\gamma)) \leq 2e^{C_0r_{\gamma}} \leq 2e^{C_0r_{n+1}}$.
Applying inequality (71) or (72) implies
\[ R(\eta)/C_7 \leq R(\eta). \]  
Due to $\eta$ in $B_\gamma$, by (65), we obtain $f_{\gamma}(\eta) \geq f(\tilde{x}_\gamma^m) = r_{\gamma}^{-\delta}$. Putting it all together, we get $P_{n+1}R(\eta) \gg R(\eta)$.

By Lemma A.8, there is at most $C_3$ element $\gamma$ such that $B_\gamma$ contains $\eta$. For these $\gamma$, by (74), we have
\[ \sum_{\gamma \in S_{n+1}, \eta \in B_\gamma} R(\eta)/r_{\gamma}^{\delta}f_{\gamma}(\eta) \ll C_3R(\eta). \]  
For the rest of $\gamma$’s, recall that $B_\gamma' = O_{C_1}(a, \gamma^{-1}0)$ is a smaller ball in $B_\gamma$. Due to Lemma A.5, the radius of $B_\gamma'$ is $r(B_\gamma') := \sinh C_1/\sinh(\kappa(\gamma))$. For $\gamma$ such that $\eta \notin B_\gamma$, we know there exists $C > 0$ such that
\[ d_o(\eta, B_\gamma') > r_{n+1}/C. \]  
This is due to $r(B_\gamma) - r(B_\gamma') \gg r_{\gamma} \gg r_{n+1}$. By (76) and (71), we know that even $d(\eta, \eta') < r_{n+1}$ we also have $R(\eta') \ll R(\eta)d_o(\eta, \eta')^{1/2}/(1 - \beta)^n$. Together with (72) and (67), we have
\[ \sum_{\gamma \in S_{n+1}, \eta \notin B_\gamma} R(\eta)/r_{\gamma}^{\delta}f_{\gamma}(\eta) \ll R(\eta)(1 - \beta)^{-n} \sum_{\gamma \in S_{n+1}, \eta \notin B_\gamma} r_{\gamma}^{\delta}f_{\gamma}(\eta)d_o(\eta, \eta')^{1/2} \leq R(\eta)(1 - \beta)^{-n} \sum_{\gamma \in S_{n+1}, \eta \notin B_\gamma} r_{\gamma}^{2\delta}d_o(\eta, \eta')^{1/2}d_o(\eta, \tilde{x}_\gamma^m)^{-2\delta}. \]  
Due to (62), we get $r(B_\gamma) - r(B_\gamma') \geq 4r(B_\gamma')$. This implies for $\xi$ in $B_\gamma'$, we have $d_o(\eta, \eta') \geq d_o(\eta, \xi) - d_o(\xi, \eta) \geq \frac{1}{2}d_o(\eta, \xi)$, which is also true if we replace $\eta$ by $\tilde{x}_\gamma^m$. Together with Lemma A.3, that is $r_{\gamma}^{\delta} \approx \mu(B_\gamma)$, and (75), (77), (68), we obtain
\[ P_{n+1}R(\eta) \ll R(\eta) \left( 1 + (1 - \beta)^{-n} \sum_{\gamma \in S_{n+1}, \eta \notin B_\gamma} \frac{1}{(B_\gamma')^2(\mu(\tilde{x}_\gamma^m)} \right). \]  
28
By Lemma A.8, the union of balls $B_\gamma$ covers at most $C_3$ times, which is also true for smaller covers $B'_\gamma$. By (76), this implies
\[
P_{n+1}R(\eta) \ll R(\eta) \left(1 + (1 - \beta)^{-n}r_{n+1}^{\delta + x_1} \frac{1}{d_\theta(\eta, \xi)^{2\delta + x}} \mu(\xi) \right).
\] (78)

**Lemma A.12.** Let $\theta$ be a positive number. For all $r > 0$ and $\eta \in \Lambda_\Gamma$, we have
\[
\int_{B(\eta, r)} \frac{1}{d_\theta(\eta, \xi)^{\delta + \theta}} d\mu(\xi) \ll \frac{1}{r^\theta}. \quad (79)
\]

Therefore, Lemma A.12 and (78) imply
\[
P_{n+1}R(\eta) \ll R(\eta) (1 + (1 - \beta)^{-n}r_{n+1}^{\delta - x} = R(\eta)(1 + (r_{n+1}/r_n)^x).
\]

The proof is complete.

It remains to prove Lemma A.12.

**Proof of Lemma A.12.** Due to [43, Theorem 7], we have that $\mu(B(\eta, r)) \ll r^\delta$ for all balls in $\partial X$ with $\eta \in \Lambda_\Gamma$ and $r > 0$.

Then
\[
\int_{B(\eta, r)} \frac{1}{d_\theta(\eta, \xi)^{\delta + \theta}} d\mu(\xi) = \sum_{1 \leq n \leq 1/r} \int_{B(\eta, (n+1)r) - B(\eta, nr)} \frac{1}{d_\theta(\eta, \xi)^{\delta + \theta}} d\mu(\xi)
\]
\[
\leq \sum_{1 \leq n \leq 1/r} \int_{B(\eta, (n+1)r) - B(\eta, nr)} \frac{1}{(nr)^{\delta + \theta}} d\mu(\xi)
\]
\[
\leq r^{-(\delta + \theta)} \left( \sum_{1 \leq n \leq 1/r} \mu(B(\eta, (n+1)r)) \left( \frac{1}{n^{\delta + \theta}} - \frac{1}{(n+1)^{\delta + \theta}} \right) - \mu(B(\eta, r)) \right)
\]
\[
\ll r^{-(\delta + \theta)} \left( \sum_{n \geq 1} ((n+1)r)^{\delta} \left( \frac{1}{n^{\delta + \theta}} - \frac{1}{(n+1)^{\delta + \theta}} \right) \right) \ll \theta r^{-\theta}.
\]

The proof is complete.

**A.3 Proof of Theorem A.1**

We start to prove our main theorem in this section. We will construct $\{u_n\}_{n \in \mathbb{N}}$ by induction such that
\[
|1 - \sum_{n \leq M} u_n(\eta)| \to 0 \text{ as } M \to \infty, \text{ uniformly for all } \eta \in \Lambda_\Gamma,
\]
where $u_n$ is a finite linear combination of $f_\gamma$. The main idea is the same as that in [12]. Once we have a function on $\Lambda(\Gamma)$ which satisfies the conditions in Lemma A.11, we can use the operator $P_{n+1}$ to drop some mass for elements in $S_{n+1}$.

Let $R_0 = 1$ be the constant function on $\partial X$. We now proceed by induction. For $n$ in $\mathbb{N}$, let
\[
u_{n+1} = \frac{\beta}{C_\gamma} P_{n+1}R_n, \ R_{n+1} = R_n - u_{n+1}.
\]

The following lemma is similar to [30, Lemma 3].

---

\footnote{We write $f \ll \theta g$ for two real functions if there exists a constant $C$ depending on the group and $\theta$ such that $f \leq Cg$.}
Lemma A.13. For any \( n \) in \( \mathbb{N} \), the following holds. The function \( R_n \) is positive on \( \Gamma \) and for \( \xi, \eta \in \Gamma \), if \( d_0(\xi, \eta) \leq r_{n+1} \), then we have
\[
|R_n(\xi)/R_n(\eta) - 1| \leq (d_0(\xi, \eta)/r_{n+1})^{1/2}. \tag{80}
\]
For \( \xi, \eta \) in \( \Gamma \), if \( d_0(\xi, \eta) > r_{n+1} \), then
\[
|R_n(\xi)/R_n(\eta)| \leq C_5 d_0(\xi, \eta)^r/(1-\beta)^n. \tag{81}
\]

Proof. The proof is by induction on \( n \). For \( n = 0 \), two inequalities hold trivially. Suppose they hold for \( n \), we will prove they also hold for \( n+1 \). By the induction hypothesis and Lemma A.11, we know
\[
u_{n+1}(\eta) \leq \beta R_n(\eta) \text{ for } \eta \in \Gamma, \tag{82}
\]
which implies that \( R_{n+1} \) is always a positive function on \( \Gamma \).

Due to Lemma A.10, if \( d_0(\xi, \eta) < r_{n+2} \), then
\[
|u_{n+1}(\xi)/u_{n+1}(\eta) - 1| = |P_{n+1} R_n(\xi)/P_{n+1} R_n(\eta) - 1| \leq (d_0(\xi, \eta)/r_{n+2})^{1/2}. \tag{83}
\]
Hence, for \( \xi, \eta \) such that \( d_0(\xi, \eta) < r_{n+2} \), by (82), (80), (83) we have
\[
\left|\frac{R_{n+1}(\xi)}{R_{n+1}(\eta)} - 1\right| = \left|\frac{R_n(\xi) - u_{n+1}(\xi)}{R_n(\eta) - u_{n+1}(\eta)} - 1\right| = \left|\frac{R_n(\xi)/R_n(\eta) - 1 - u_{n+1}(\xi)/u_{n+1}(\eta) - 1}{1 - u_{n+1}(\xi)/u_{n+1}(\eta)} - 1\right|
\leq \frac{d_0(\xi, \eta)^{1/2}}{r_{n+1}^{1/2}(1 - \beta)} + \frac{d_0(\xi, \eta)^{1/2}}{r_{n+2}^{1/2}}(e^{-C_4/2} + \beta)/(1 - \beta) \leq \left(\frac{d_0(\xi, \eta)}{r_{n+2}}\right)^{1/2}.
\]
It remains to prove (81). By construction and (82), we have
\[
R_{n+1}(\xi)/R_{n+1}(\eta) = (R_n(\xi) - u_{n+1}(\xi))/(R_n(\eta) - u_{n+1}(\eta)) \leq |R_n(\xi)/R_n(\eta)|/(1 - \beta). \tag{84}
\]
If \( d_0(\xi, \eta) \geq r_{n+1} \), then due to (84), the inequality (81) holds for \( n+1 \) is a direct consequence of case \( n \). If else, we have \( r_{n+2} < d_0(\xi, \eta) \leq r_{n+1} \). By (80) we have \( R_n(\xi)/R_n(\eta) \leq 2 \), then by (84)
\[
R_{n+1}(\xi)/R_{n+1}(\eta) \leq 2/(1 - \beta) = 2r_{n+1}/(1 - \beta)^{n+2} \leq \frac{2e^{-C_4}}{1 - \beta (1 - \beta)^{n+1} + C_5 d_0(\xi, \eta)^r/(1-\beta)^n}. \tag{85}
\]
The proof is complete. \( \square\)

Proof of Theorem A.1. We start to prove our theorem. Let \( C = C_6 C_7 \), where \( C_6, C_7 \) are constants in Lemma A.11. Lemma A.13 implies that the constructed \( R_n \) is positive on \( \Gamma \) and always satisfies the condition in Lemma A.11 for \( n \in \mathbb{N} \). Hence for a point \( \eta \) in \( \Gamma \), we apply Lemma A.11 to obtain
\[
R_{n+1}(\eta) = R_n(\eta)(1 - u_{n+1}(\eta)/R_n(\eta)) = R_n(\eta)\left(1 - \frac{\beta P_{n+1} R_n(\eta)}{C_7 R_n(\eta)}\right) \leq R_n(\eta)(1 - \beta/C).
\]
Iterating the above inequality, we get \( R_n(\eta) \leq (1 - \beta/C)^n \). Therefore, \( R_n \to 0 \) uniformly on \( \Gamma \) as \( n \to \infty \).

We set
\[
\nu(\gamma) = \begin{cases} R_{n-1}(\eta_\gamma) r_7^{\delta} \beta/C_7 & \text{for } n \in \mathbb{N}_0, \gamma \in S_n, \\ 0 & \text{for } \gamma \notin \bigcup_{n \in \mathbb{N}_0} S_n. \end{cases} \tag{85}
\]
Then \( R_n - R_{n+1} = \sum_{\gamma \in S_{n+1}} \nu(\gamma) f_\gamma \). It follows that \( 1 = R_0 = \sum_{\gamma \in \Gamma} \nu(\gamma) f_\gamma \) on \( \Gamma \), which means that \( \mu \) is \( \nu \)-stationary by (64).
Next we verify the moment condition. Let \( \epsilon_1 \) be a positive number. Let \( \| \gamma \| \) be the operator norm of its action on \( \mathbb{R}^{n+1} \) equipped with euclidean norm. By the Cartan decomposition, we obtain \( \| \gamma \| = r_\gamma^{-1} \) (see for example [3, Remark 6.28 and Lemma 6.33]). We can compute the exponential moment
\[
\sum_{\gamma \in \Gamma} \nu(\gamma) \| \gamma \|^{\epsilon_1} = \sum_{n \in \mathbb{N}_0} \sum_{\gamma \in S_n} \nu(\gamma) \| \gamma \|^{\epsilon_1} \leq \frac{\beta}{C_7} \sum_{n \in \mathbb{N}_0} \sum_{\gamma \in S_n} R_{n-1}(\eta_\gamma) r_\gamma^{-\epsilon_1}.
\]

While \( r_\gamma^\delta \approx \mu(B_\gamma) \) (Lemma A.3), we have \( \sum_{\gamma \in S_n} r_\gamma^\delta \ll \sum_{\gamma \in S_n} \mu(B_\gamma) = 1. \) Due to \( r_\gamma \geq e^{-C_4(n+1)} \) and \( R_n \leq (1 - \beta/C)^n, \) we get
\[
\sum_{\gamma \in \Gamma} \nu(\gamma) \| \gamma \|^{\epsilon_1} \ll \sum_{n \in \mathbb{N}} (1 - \beta/C)^n e^{\epsilon_1 C_4(n+1)}.
\]

Take \( \epsilon_1 \) small enough, the above sum is finite.

Lastly we prove \( \Gamma_\nu = \Gamma. \) Since the diameter of \( C(\Gamma) \) is less than \( C_0/2, \) there exists \( \gamma_1 \) in \( S_1 \) such that \( d(o, \gamma_1 o) \in [\| \log r_1 \| + C_0/2, \| \log r_1 \| + 3C_0/2] \). By construction of \( S_1 \) (68) and (85) we know that the set \( \Gamma_{C_0} := \{ \gamma \in \Gamma | d(o, \gamma o) \leq C_0/2 \} \) is contained in \( \gamma_1^{-1} S_1 \subset S_1^{-1} S_1 \subset \Gamma_\nu. \)

**Lemma A.14.** If \( C_0 \) is greater than 6 times the diameter of the quotient set \( C(\Gamma), \) then the set \( \Gamma_{C_0} \) generates the group \( \Gamma. \)

By Lemma A.14, the proof is complete. \( \square \)

It remains to prove Lemma A.14.

**Proof of Lemma A.14.** This is a classical lemma. We give a proof here. Let \( C_\Gamma \) be the diameter of the quotient \( C(\Gamma) \). For any \( \gamma \) in \( \Gamma, \) we will find a sequence \( \beta_j, 0 \leq j \leq k \) in \( \Gamma_{C_0} \) such that \( \gamma = \beta_0 \cdots \beta_k, \) which finishes the proof.

In the geodesic \( o(\gamma o), \) let \( p_j \) be the point with distance \( jC_\Gamma \) to \( o. \) Suppose that \( kC_\Gamma \leq d(o, \gamma o) \leq (k + 1)C_\Gamma \) and let \( p_{k+1} = \gamma o. \) Since \( o(\gamma o) \) is in the convex hull, for every \( p_j \) with \( 1 \leq j \leq k, \) by the definition of \( C_\Gamma, \) we can find \( \gamma_j \) in \( \Gamma \) such that \( d(\gamma_j o, p_j) \leq C_\Gamma. \) Let \( \gamma_0 = e \) and \( \gamma_{k+1} = \gamma. \) Hence for \( 0 \leq j \leq k \)
\[
d(\gamma_j o, \gamma_{j+1} o) \leq d(\gamma_j o, p_j) + d(p_j, p_{j+1}) + d(p_{j+1}, \gamma_{j+1} o) \leq 3C_\Gamma.
\]

Let \( \beta_j = \gamma_j^{-1} \gamma_{j+1}. \) Then \( d(\beta_j o, o) = d(\gamma_j^{-1} \gamma_{j+1} o, o) = d(\gamma_{j+1} o, \gamma_j o) \leq 3C_\Gamma \leq C_0, \) which implies \( \beta_j \in \Gamma_{C_0}. \) Therefore \( \gamma = \beta_0 \cdots \beta_k. \) \( \square \)

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