# ENTROPY AND DRIFT IN WORD HYPERBOLIC GROUPS 

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#### Abstract

The fundamental inequality of Guivarc'h relates the entropy and the drift of random walks on groups. It is strict if and only if the random walk does not behave like the uniform measure on balls. We prove that, in any nonelementary hyperbolic group which is not virtually free, endowed with a word distance, the fundamental inequality is strict for symmetric measures with finite support, uniformly for measures with a given support. This answers a conjecture of S. Lalley. For admissible measures, this is proved using previous results of Ancona and Blachère-Haïssinsky-Mathieu. For non-admissible measures, this follows from a counting result, interesting in its own right: we show that, in any infinite index subgroup, the proportion of non-distorted points is exponentially small. The uniformity is obtained by studying the behavior of measures that degenerate towards a measure supported on an elementary subgroup.


## 1. Main Results

Let $\Gamma$ be a finitely generated infinite group. Although the following discussion makes sense in a much broader context, we will assume that $\Gamma$ is hyperbolic since all results of this article are devoted to this setting. There are two natural ways to construct random elements in $\Gamma$ :

- Let $d$ be a proper left-invariant distance on $\Gamma$ (for instance a word distance). For large $n$, one can pick an element at random with respect to the uniform measure $\rho_{n}$ on the ball $B_{n}=B(e, n)$ (where $e$ denotes the identity of $\Gamma$ ).
- Let $\mu$ be a probability measure on $\Gamma$. For large $n$, one can pick an element at random with respect to the measure $\mu^{* n}$ (the $n$-th convolution of the measure $\mu$ ). Equivalently, let $g_{1}, g_{2}, \ldots$ be a sequence of random elements of $\Gamma$ that are distributed independently according to $\mu$. Form the random walk $X_{n}=g_{1} \cdots g_{n}$. Then the distribution of $X_{n}$ is $\mu^{* n}$.
From a theoretical point of view, these methods share a lot of properties. From a computational point of view, the second method is much easier to implement in general groups since it does not require the computation of the ball $B_{n}$ (note however that, in hyperbolic groups, simulating the uniform measure is very easy thanks to the automatic structure of the group). It is therefore of interest to find probability measures $\mu$ such that these two methods give equivalent results, in a sense that will be made precise below. This is the main question of Vershik in [Ver00]. In free groups (with the word distance coming from the usual set of generators), everything can be computed: if $\mu$ is the uniform measure on the generators, then $\mu^{* n}$ and $\rho_{n}$ behave essentially in the same way. The situation is the

[^0]same in free products of finite groups, again thanks to the underlying tree structure. However, in more complicated groups, explicit computations are essentially impossible, and it is expected that the methods always differ. Our main result confirms this intuition in a special class of groups: In hyperbolic groups which are not virtually free (i.e., there is no finite index free subgroup), if $d$ is a word distance, the two methods are always different, in a precise quantitative way.

Remark 1.1. We emphasize that the question really depends on the choice of the distance $d$, since the shape of the balls $B_{n}$ depends on $d$. For instance, for any symmetric probability measure $\mu$ on $\Gamma$ whose support is finite and generates $\Gamma$, there exists a distance $d$ (called the Green distance, see [BHM11]) for which the measures $\rho_{n}$ and $\mu^{* n}$ behave in the same way. A famous open problem (to which our methods do not apply) is to understand what happens when $\Gamma$ acts cocompactly on the hyperbolic space $\mathbb{H}^{k}$, and the distance $d$ is given by $d(e, \gamma)=d_{\mathbb{H}^{k}}(O, \gamma \cdot O)$ where $O$ is a base point in $\mathbb{H}^{k}$. In this case, it is also expected that the two methods are always different. Here are the main partial results in this context:
(1) The two methods are different for some symmetric measures with finite support ([LP07], see also Theorem 5.9 below).
(2) If, instead of a cocompact lattice, one considers a lattice with cusps, the two methods are always different [GLJ93].
(3) If, instead of a lattice, one considers a nice dense subgroup, there exist symmetric measures with finite support for which the two methods are equivalent [Bou12].

This question also makes sense in continuous time, for negatively curved manifolds. A conjecture of Sullivan asserts that, in this setting, the two methods coincide if and only if the manifold is locally symmetric, see [Led95].

One can give several meanings to the question "are the two methods equivalent?" Let us first discuss an interpretation in terms of behavior at infinity. The measures $\mu^{* n}$ converge in the geometric compactification $\Gamma \cup \partial \Gamma$ to a measure $\mu_{\infty}$, supported on the boundary, called the exit measure of the random walk, or its stationary measure. Geometrically, the random walk $\left(X_{n}\right)_{n \geqslant 1}$ converges almost surely to a random point on the boundary $\partial \Gamma$, the measure $\mu_{\infty}$ is its distribution. On the other hand, let $\rho_{\infty}$ be the Patterson-Sullivan measure on $\partial \Gamma$ associated to the distance $d$, constructed in [Coo93] in this context. One should think of it as the uniform measure on the boundary (it is equivalent to the Hausdorff measure of maximal dimension on the boundary, for any visual distance coming from $d$ ). The measures $\rho_{n}$ do not always converge to $\rho_{\infty}$, but all their limit points are equivalent to $\rho_{\infty}$, with a density bounded from above and from below (this follows from the arguments of [Coo93], see Lemma 2.13 below). A version of the question is then to ask if the measures $\mu_{\infty}$ and $\rho_{\infty}$ are mutually singular: in this case, the random walk mainly visits parts of the groups that are not important from the point of view of the uniform measure.

Another version of the same question is quantitative: Does the random walk visit parts of the groups that are exponentially negligible from the point of view of the uniform measure? This is made precise through the notions of drift and entropy. Define

$$
\begin{equation*}
L(\mu)=\sum_{g \in \Gamma} \mu(g)|g|, \quad H(\mu)=\sum_{g \in \Gamma} \mu(g)(-\log \mu(g)), \tag{1.1}
\end{equation*}
$$

where $|g|=d(e, g)$. The quantity $L(\mu)$ is the average distance of an element to the identity. The quantity $H(\mu)$, called the time one entropy of $\mu$, is the average logarithmic weight of the points. They can both be finite or infinite. The functions $L$ and $H$ both behave in a subadditive way with respect to convolution: $L\left(\mu_{1} * \mu_{2}\right) \leqslant L\left(\mu_{1}\right)+L\left(\mu_{2}\right)$ and $H\left(\mu_{1} * \mu_{2}\right) \leqslant$ $H\left(\mu_{1}\right)+H\left(\mu_{2}\right)$. It follows that the sequences $L\left(\mu^{* n}\right)$ and $H\left(\mu^{* n}\right)$ are subadditive. Hence, the following quantities are well defined:

$$
\begin{equation*}
\ell(\mu)=\lim L\left(\mu^{* n}\right) / n, \quad h(\mu)=\lim H\left(\mu^{* n}\right) / n . \tag{1.2}
\end{equation*}
$$

They are called respectively the drift and the asymptotic entropy of the random walk. They also admit characterizations along typical trajectories. If $L(\mu)$ is finite, then almost surely $\ell(\mu)=\lim \left|X_{n}\right| / n$. In the same way, if $H(\mu)$ is finite, then almost surely $h(\mu)=$ $\lim \left(-\log \mu^{* n}\left(X_{n}\right)\right) / n$. The most intuitive characterization of the entropy is probably the following one: at time $n$, the random walk is essentially supported by $e^{h(\mu) n}$ points (see Lemma 2.4 for a precise statement). Let us also define the exponential growth rate of the group with respect to $d$, i.e.,

$$
\begin{equation*}
v=\liminf _{n \rightarrow \infty} \frac{\log \left|B_{n}\right|}{n} \tag{1.3}
\end{equation*}
$$

where $B_{n}$ is the ball of radius $n$ around $e$. In hyperbolic groups, it satisfies the apparently stronger inequality $C^{-1} e^{n v} \leqslant\left|B_{n}\right| \leqslant C e^{n v}$, by [Coo93]. For large $n$, most points for $\mu^{* n}$ are contained in a ball $B_{(1+\varepsilon) \ell n}$, which has cardinality at most $e^{(1+2 \varepsilon) \ell n v}$. Since the random walk at time $n$ essentially visits $e^{h n}$ points, we deduce the fundamental inequality of Guivarc'h [Gui80]

$$
h \leqslant \ell v .
$$

If this inequality is an equality, this means that the walk visits most parts of the group. Otherwise, it is concentrated in an exponentially small subset. Another version of our main question is therefore: Is the inequality $h \leqslant \ell v$ strict?

In hyperbolic groups, it turns out that the two versions of the question are equivalent, at least for finitely supported measures, and that they also have a geometric interpretation in terms of Hausdorff dimension. If $\mu$ is a probability measure on a group, we write $\Gamma_{\mu}^{+}$ for the semigroup generated by the support of $\mu$, and $\Gamma_{\mu}$ for the group it generates. When $\mu$ is symmetric, they coincide. We say that $\mu$ is admissible if $\Gamma_{\mu}^{+}=\Gamma$. The following result is Corollary 1.4 and Theorem 1.5 in [BHM11] (see also [Haï13]) when the measure is symmetric, and is proved in [Tan14] when $\mu$ is not necessarily symmetric and $d$ is a word distance.

Theorem 1.2. Let $\Gamma$ be a non-elementary hyperbolic group, endowed with a left-invariant distance $d$ which is hyperbolic and quasi-isometric to a word distance. Let $v$ be the exponential growth rate of $(\Gamma, d)$. Let $d_{\partial \Gamma}$ be a visual distance on $\partial \Gamma$ associated to d. Consider an admissible probability measure $\mu$ on $\Gamma$, with finite support. Assume additionally either that the measure $\mu$ is symmetric, or that the distance $d$ is a word distance. The following conditions are equivalent:
(1) The equality $h=\ell v$ holds.
(2) The Hausdorff dimension of the exit measure $\mu_{\infty}$ on $\left(\partial \Gamma, d_{\partial \Gamma}\right)$ is equal to the Hausdorff dimension of this space.
(3) The measure $\mu_{\infty}$ is equivalent to the Patterson-Sullivan measure $\rho_{\infty}$.
(4) The measure $\mu_{\infty}$ is equivalent to the Patterson-Sullivan measure $\rho_{\infty}$, with density bounded from above and from below.
(5) There exists $C>0$ such that, for any $g \in \Gamma$,

$$
\left|v d(e, g)-d_{\mu}(e, g)\right| \leqslant C
$$

where $d_{\mu}$ is the "Green distance" associated to $\mu$, i.e., $d_{\mu}(e, g)=-\log \mathbb{P}\left(\exists n, X_{n}=\right.$ $g$ ) where $X_{n}$ is the random walk given by $\mu$ starting from the identity (it is an asymmetric distance in general, and a genuine distance if $\mu$ is symmetric).

The different statements in this theorem go from the weakest to the strongest: since entropy is an asymptotic quantity, an assumption on $h$ seems to allow subexponential fluctuations, so the assumption (1) is rather weak. On the other hand, (3) says that two measures are equivalent, so most points are controlled. Finally, in (5), all points are uniformly controlled. The equivalence between these statements is a strong rigidity theorem. The equivalence between (1) and (2) follows from a formula for the respective dimensions. The definition of a visual distance at infinity $d_{\partial \Gamma}$ involves a small parameter $\varepsilon$. In terms of this parameter, one has $H D\left(\mu_{\infty}\right)=h /(\varepsilon \ell)$ and $H D\left(\rho_{\infty}\right)=H D(\partial \Gamma)=v / \varepsilon$, so that these dimensions coincide if and only if $h=\ell v$.

In this theorem, the finite support assumption can be weakened to an assumption of superexponential moment (i.e., for all $M>0, \sum_{g \in \Gamma} \mu(g) e^{M|g|}<\infty$ ), thanks to [Gou13]. The assumption that $\mu$ is symmetric or that $d$ is a word distance is probably not necessary. However, the most important assumption in Theorem 1.2 is admissibility: it ensures that the random walk can see the geometry of the whole group (which is hyperbolic). For a random walk living in a strict (maybe distorted) subgroup, one would not be expecting the same nice behavior.

Our main theorem follows. It states that, in hyperbolic groups which are not virtually free, endowed with a word distance, the different equivalent conditions of Theorem 1.2 are never satisfied, uniformly on measures with a fixed support.
Theorem 1.3. Let $\Gamma$ be a hyperbolic group which is not virtually free, endowed with a word distance $d$. Let $\Sigma$ be a finite subset of $\Gamma$. There exists $c<1$ such that, for any symmetric probability measure $\mu$ supported in $\Sigma$,

$$
h(\mu) \leqslant c \ell(\mu) v
$$

where $v$ is the exponential growth rate of balls in $(\Gamma, d)$.
This theorem gives a positive answer to a conjecture of S. Lalley [Lal14, slide 16]. In the language of Vershik [Ver00], this theorem says that no finite subset of $\Gamma$ is extremal. On the other hand, if one lets $\Sigma$ grow, $h / \ell$ can converge to $v$ :

Theorem 1.4. Let $\Gamma$ be a hyperbolic group, endowed with a left invariant distance $d$ which is hyperbolic and quasi-isometric to a word distance. Let $\rho_{i}$ be the uniform measure on the ball of radius $i$. Then $h\left(\rho_{i}\right) / \ell\left(\rho_{i}\right) \rightarrow v$, where $v$ is the exponential growth rate of balls in $(\Gamma, d)$.

More precisely, we prove that $\ell\left(\rho_{i}\right) \sim i$ and $h\left(\rho_{i}\right) \sim i v$. The only difficulty is to prove the lower bound on $h\left(\rho_{i}\right)$ : since $h$ is defined in (1.2) using a subadditive sequence, upper bounds
are automatic, but to get lower bounds one should show that additional cancellations do not happen later on. This difficulty already appears in [EK13], where the authors prove that the entropy depends continuously on the measure. Our proof of Theorem 1.4, given in Paragraph 2.5, also applies to this situation and gives a new proof of their result, under slightly weaker assumptions. There is nothing special about the uniform measure on balls, our proof also gives the same conclusion for the uniform measure on spheres, or for the measures $\sum e^{-s|g|} \delta_{g} / \sum e^{-s|g|}$ when $s \searrow v$.

Our main result is Theorem 1.3. It is a consequence of the three following results. Since their main aim is Theorem 1.3, they are designed to handle finitely supported symmetric measures. However, these theorems are all valid under weaker assumptions, which we specify in the statements as they carry along implicit information on the techniques used in the proofs.

The first result deals with admissible (or virtually admissible) measures.
Theorem 1.5. Let $\Gamma$ be a hyperbolic group which is not virtually free, endowed with a word distance. Let $\mu$ be a probability measure with a superexponential moment, such that $\Gamma_{\mu}^{+}$is a finite index subgroup of $\Gamma$. Then $h(\mu)<\ell(\mu) v$.

The second result deals with non-admissible measures.
Theorem 1.6. Let $\Gamma$ be a hyperbolic group endowed with a word distance. Let $\mu$ be a probability measure with a moment of order 1 (i.e., $L(\mu)<\infty$ ). Assume that $\ell(\mu)>0$ and that $\Gamma_{\mu}$ has infinite index in $\Gamma$. Then $h(\mu)<\ell(\mu) v$.

Finally, the third result is a kind of continuity statement, to get the uniformity.
Theorem 1.7. Let $\Gamma$ be a hyperbolic group, endowed with a left-invariant distance which is hyperbolic and quasi-isometric to a word distance. Let $\Sigma$ be a subset of $\Gamma$ which does not generate an elementary subgroup. There exists a probability measure $\mu_{\Sigma}$ with finite support such that $\ell\left(\mu_{\Sigma}\right)>0$ and

$$
\sup \{h(\mu) / \ell(\mu): \mu \text { probability }, \operatorname{Supp}(\mu) \subset \Sigma, \ell(\mu)>0\}=h\left(\mu_{\Sigma}\right) / \ell\left(\mu_{\Sigma}\right) .
$$

The same statement holds if the maximum is taken over symmetric probability measures, the resulting maximizing measure being symmetric.

Theorem 1.3 is a consequence of these three statements.
Proof of Theorem 1.3 using the three auxiliary theorems. As in the statement of the theorem, consider a finite subset $\Sigma$ of $\Gamma$. If $\Sigma$ generates an elementary subgroup of $\Gamma$, all measures supported on $\Sigma$ have zero entropy. Hence, one can take $c=0$ in the statement of the theorem. Otherwise, by Theorem 1.7, there exists a symmetric measure $\mu_{\Sigma}$ with finite support that maximizes the quantity $h(\mu) / \ell(\mu)$ over $\mu$ symmetric supported by $\Sigma$. If $\Gamma_{\mu_{\Sigma}}=\Gamma_{\mu_{\Sigma}}^{+}$has finite index, $h\left(\mu_{\Sigma}\right) / \ell\left(\mu_{\Sigma}\right)<v$ by Theorem 1.5. If it has infinite index, the same conclusion follows from Theorem 1.6.

The three auxiliary theorems are non-trivial. Their proofs are independent, and use completely different tools. Here are some comments about them.

- At first sight, Theorem 1.5 seems to be the most delicate (this is the only one with the assumption that $\Gamma$ is not virtually free). However, this is also the setting that has been mostly studied in the literature. Hence, we may use several known results, including most notably results of Ancona [Anc87], of Blachère, Haïssinsky and Mathieu [BHM11] and Tanaka [Tan14] (Theorem 1.2 above) and of Izumi, Neshveyev and Okayasu [INO08] on rigidity results for cocycles. The proof relies mainly on the fact that the word distance is integer valued, contrary to the Green distance (more precisely, we use the fact that the stable translation length of hyperbolic elements is rational with bounded denominator).
- In Theorem 1.6, the difficulty comes from the lack of information on the subgroup $\Gamma_{\mu}$. If it has good geometric properties (for instance if it is quasi-convex), one may use the same kind of techniques as for Theorem 1.5. Otherwise, the random walk does not really see the hyperbolicity of the ambient group. The fundamental inequality always gives $h \leqslant \ell v_{\Gamma_{\mu}}$, where $v_{\Gamma_{\mu}}$ is the growth rate of the subgroup $\Gamma_{\mu}$ (for the initial word distance on $\Gamma$ ). If $v_{\Gamma_{\mu}}<v$, the result follows. Unfortunately, there exist non-quasi-convex subgroups of some hyperbolic groups with the same growth as the ambient group. However, a random walk does not typically visit all points of $\Gamma_{\mu}$, it concentrates on those points that are not distorted (i.e., their distances to the identity in $\Gamma$ and $\Gamma_{\mu}$ are comparable). To prove Theorem 1.6, we will show that in any infinite index subgroup of a hyperbolic group, the number of non-distorted points is exponentially smaller than $e^{n v}$.
- Theorem 1.7 is less simple than it may seem at first sight: it does not claim that $\mu_{\Sigma}$ is supported by $\Sigma$, and indeed this is not the case in general (see Example 5.4). Hence, the proof is not a simple continuity argument: We need to understand precisely the behavior of sequences of measures that degenerate towards a measure supported on an elementary subgroup. The proof will show that $\mu_{\Sigma}$ is supported by $K \cdot(\Sigma \cup\{e\}) \cdot K$, where $K$ is a finite subgroup generated by some elements in $\Sigma$.

A natural question is whether Theorem 1.3 holds for non-symmetric measures. For admissible measures, (i.e., $\Gamma_{\mu}^{+}=\Gamma$ ), Theorem 1.5 holds. For non-symmetric measures such that $\Gamma_{\mu}$ has infinite index, Theorem 1.6 applies directly. However, since $\Gamma_{\mu} \neq \Gamma_{\mu}^{+}$for general non-symmetric measures, there is another case to consider: the case of measures $\mu$ such that $\Gamma_{\mu}=\Gamma$ ( or $\Gamma_{\mu}$ has finite index in $\Gamma$ ), but $\Gamma_{\mu}^{+}$is much smaller than $\Gamma$. In this case, it seems that our arguments do not suffice. We give in Section 6 two examples illustrating the new difficulties:
(1) One can not rely on growth arguments, as for Theorem 1.6. Indeed, there are subsemigroups $\Lambda^{+}$with bad asymptotic behavior, for instance such that $\lim \inf \mid B_{n} \cap$ $\Lambda^{+}\left|/\left|B_{n}\right|=0\right.$ and $\left.\lim \sup \right| B_{n} \cap \Lambda^{+}\left|/\left|B_{n}\right|>0\right.$.
(2) The arguments of Theorem 1.5 work for finitely supported measures, or for measures with a superexponential moment, but also more generally for measures with a nice geometric behavior (they should satisfy so-called Ancona inequalities). In the nonsymmetric situation, we give in Proposition 6.2 explicit examples of (non-admissible) measures with an exponential moment and a very nice geometric behavior, and such that nevertheless $h=\ell v$. So, arguments similar to those of Theorem 1.5 can not
suffice, one needs a new argument that distinguishes in a finer way between measures with finite support and measures with infinite support.
This article is organized as follows. In Section 2, we give more details on the notions of hyperbolic group, drift and entropy. We also prove Theorem 1.4 on the asymptotic entropy and drift of the uniform measure on large balls. The following three sections are then devoted to the proofs of the three auxiliary theorems. Finally, we describe in Section 6 what can happen in the non-symmetric setting. In particular, we show that in any torsionfree group with infinitely many ends, there exist (non-admissible, non-symmetric) measures with an exponential moment satisfying $h=\ell v$.

## 2. General properties of entropy and drift in hyperbolic groups

2.1. Hyperbolic spaces. In this paragraph, we recall classical properties of hyperbolic spaces. See for instance [GdlH90] or [BH99].

Consider a metric space $(X, d)$. The Gromov product of two points $y, y^{\prime} \in X$, based at $x_{0} \in X$, is by definition

$$
\begin{equation*}
\left(y \mid y^{\prime}\right)_{x_{0}}=(1 / 2)\left[d\left(x_{0}, y\right)+d\left(x_{0}, y^{\prime}\right)-d\left(y, y^{\prime}\right)\right] \tag{2.1}
\end{equation*}
$$

The space $(X, d)$ is hyperbolic if there exists $\delta \geqslant 0$ such that, for any $x_{0}, y_{1}, y_{2}, y_{3}$, the following inequality holds:

$$
\left(y_{1} \mid y_{3}\right)_{x_{0}} \geqslant \min \left(\left(y_{1} \mid y_{2}\right)_{x_{0}},\left(y_{2} \mid y_{3}\right)_{x_{0}}\right)-\delta
$$

The main intuition to have is that, in hyperbolic spaces, configurations of finitely many points look like configurations in trees: for any $k$, for any subset $F$ of $X$ with cardinality at most $k$, there exists a map $\Phi$ from $F$ to a tree such that, for all $x, y \in F$,

$$
d(x, y)-2 k \delta \leqslant d(\Phi(x), \Phi(y)) \leqslant d(x, y)
$$

Hence, a lot of distance computations can be reduced to equivalent computations in trees (which are essentially combinatorial), up to a bounded error. Up to $\delta$, the Gromov product $\left(y \mid y^{\prime}\right)_{x_{0}}$ is, in the approximating tree, the length of the part that is common to the geodesics from $x_{0}$ to $y$ and from $x_{0}$ to $y^{\prime}$.

A space $(X, d)$ is geodesic if there exists a geodesic between any pair of points. For such spaces, there is a convenient characterization of hyperbolicity. A geodesic space $(X, d)$ is hyperbolic if and only if there exists $\delta \geqslant 0$ such that its geodesic triangles are $\delta$-thin, i.e., each side is included in the $\delta$-neighborhood of the union of the two other sides.

Assume that $\left(X, d_{X}\right)$ and $\left(Y, d_{Y}\right)$ are two geodesic metric spaces, and that they are quasiisometric. If $\left(X, d_{X}\right)$ is hyperbolic, then so is $\left(Y, d_{Y}\right)$. Note however that this equivalence only holds for geodesic spaces.

Let $(X, d)$ be a geodesic hyperbolic metric space. A subset $Y$ of $X$ is quasi-convex if there exists a constant $C$ such that, for any $y, y^{\prime} \in Y$, the geodesics from $y$ to $y^{\prime}$ stay in the $C$-neighborhood of $Y$.

We will sometimes encounter hyperbolic spaces which are not geodesic, but only quasigeodesic: there exist constants $C>0$ and $\lambda$ such that any two points can be joined by a $(\lambda, C)$-quasi-geodesic, i.e., a map $f$ from a real interval to $X$ such that $\lambda^{-1}\left|t^{\prime}-t\right|-C \leqslant$ $d\left(f(t), f\left(t^{\prime}\right)\right) \leqslant \lambda\left|t^{\prime}-t\right|+C$. When the space is geodesic, a quasi-geodesic stays a bounded distance away from a true geodesic. Most properties that hold or can be defined using
geodesics (for instance the notion of quasi-convexity) can be extended to this setting, simply replacing geodesics with quasi-geodesics in the statements.

Let $(X, d)$ be a proper geodesic hyperbolic space. Its boundary at infinity $\partial X$ is by definition the set of geodesics originating from a base point $x_{0}$, where two such geodesics are identified if they remain a bounded distance away. It is a compact space, which does not depend on $x_{0}$. The space $X \cup \partial X$ is also compact. If $X$ is only quasi-geodesic, all these definitions extend using quasi-geodesics instead of geodesics.

Any isometry (or, more generally, quasi-isometry) of a hyperbolic space extends continuously to its boundary, giving a homeomorphism of $\partial X$.

The Gromov product may be extended to $X \cup \partial X$ : we define $(\xi \mid \eta)_{x_{0}}$ as the infimum limit of $\left(x_{n} \mid y_{n}\right)_{x_{0}}$ for $x_{n}$ and $y_{n}$ converging respectively to $\xi$ and $\eta$. The choice to take the infimum is arbitrary, one could also take the supremum or any accumulation point, those quantities differ by at most a constant only depending on $\delta$. Hence, one should think of the Gromov product at infinity to be canonically defined only up to an additive constant. Heuristically, $(\xi \mid \eta)_{x_{0}}$ is the time after which two geodesics from $x_{0}$ to $\xi$ and to $\eta$ start diverging.

Let ( $X, d$ ) be a proper geodesic (or quasi-geodesic) hyperbolic space. For any small enough $\varepsilon>0$, one may define a visual distance $d_{\partial X, \varepsilon}$ on $\partial X$ such that $d_{\partial X, \varepsilon}(\xi, \eta) \asymp e^{-\varepsilon(\xi \mid \eta)_{x_{0}}}$ (meaning that the ratio between these quantities is uniformly bounded from above and from below).

Let $(X, d)$ be a proper hyperbolic metric space. One can define another boundary of $X$, the Busemann boundary (or horoboundary), as follows. Let $x_{0}$ be a fixed basepoint in $X$. To $x \in X$, one associates its horofunction $h_{x}(y)=d(y, x)-d\left(x_{0}, x\right)$, normalized so that $h_{x}\left(x_{0}\right)=0$. The map $\Phi: x \mapsto h_{x}$ is an embedding of $X$ into the space of 1 Lipschitz functions on $X$, with the topology of uniform convergence on compact sets. The horoboundary is obtained by taking the closure of $\Phi(X)$. In other words, a sequence $x_{n} \in X$ converges to a boundary point if $h_{x_{n}}(y)$ converges, uniformly on compact sets. Its limit is the horofunction $h_{\xi}$ associated to the corresponding boundary point $\xi$ (it is also called the Busemann function associated to $\xi$ ). We denote by $\partial_{B} X$ the Busemann boundary of $X$. There is a continuous projection $\pi_{B}: \partial_{B} X \rightarrow \partial X$, which is onto but not injective in general. The boundary $\partial_{B} X$ is rather sensitive to fine scale details of the distance $d$, while $\partial X$ only depends on its quasi-isometry class.

Any isometry $\varphi$ of $X$ acts on horofunctions, by the formula $h_{\varphi(x)}(y)=h_{x}\left(\varphi^{-1} y\right)-$ $h_{x}\left(\varphi^{-1} x_{0}\right)$. This implies that $\varphi$ extends to a homeomorphism on $\partial_{B} X$, given by the same formula $h_{\varphi(\xi)}(y)=h_{\xi}\left(\varphi^{-1} y\right)-h_{\xi}\left(\varphi^{-1} x_{0}\right)$. Note that, contrary to the action on the geometric boundary, this only works for isometries of $X$, not quasi-isometries.
2.2. Hyperbolic groups. Let $\Gamma$ be a finitely generated group, with a finite symmetric generating set $S$. Denote by $d=d_{S}$ the corresponding word distance. The group $\Gamma$ is hyperbolic if the metric space ( $\Gamma, d_{S}$ ) is hyperbolic. Since hyperbolicity is invariant under quasi-isometry for geodesic spaces, this notion does not depend on the choice of the generating set $S$. However, if one considers another left-invariant distance on $\Gamma$ which is equivalent to $d_{S}$ but not geodesic, its hyperbolicity is not automatic. Hence, one should postulate its hyperbolicity if it is needed, as in the statement of Theorem 1.2. We say that the pair $(\Gamma, d)$ is a metric hyperbolic group if the group $\Gamma$ is hyperbolic for one (or, equivalently, for any)
word distance, and if the distance $d$ is left-invariant, hyperbolic, and quasi-isometric to one (or equivalently, any) word distance. Such a distance $d$ does not have to be geodesic, but it is quasi-geodesic since geodesics for a given word distance form a system of quasi-geodesics for $d$, going from any point to any point.

Let $(\Gamma, d)$ be a metric hyperbolic group. The left-multiplication by elements of $\Gamma$ is isometric. Hence, $\Gamma$ acts by homeomorphisms on its compactifications $\Gamma \cup \partial \Gamma$ and $\Gamma \cup \partial_{B} \Gamma$. Moreover, any infinite order element $g \in \Gamma$ acts hyperbolically on $\Gamma \cup \partial \Gamma$ : it has two fixed points at infinity $g^{-}$and $g^{+}$, the points in $\Gamma \cup \partial \Gamma \backslash\left\{g^{-}\right\}$are attracted to $g^{+}$by forward iteration of $g$, and the points in $\Gamma \cup \partial \Gamma \backslash\left\{g^{+}\right\}$are attracted to $g^{-}$by backward iteration of $g$.

Definition 2.1. Consider an action of a group $\Gamma$ on a space $Z$. A function $c: \Gamma \times Z \rightarrow \mathbb{R}$ is a cocycle if, for any $g, h \in \Gamma$ and any $\xi \in Z$,

$$
\begin{equation*}
c(g h, \xi)=c(g, h \xi)+c(h, \xi) \tag{2.2}
\end{equation*}
$$

The cocycle is Hölder-continuous if $Z$ is a metric space and each function $\xi \mapsto c(g, \xi)$ is Hölder-continuous.

There is a choice to be made in the definition of cocycles, since one may compose with $g$ or $g^{-1}$. Our definition is the most customary. With this definition, the map $c_{B}: \Gamma \times \partial_{B} \Gamma \rightarrow \mathbb{R}$ given by $c_{B}(g, \xi)=h_{\xi}\left(g^{-1}\right)$ is a cocycle, called the Busemann cocycle.

A subgroup $H$ of $\Gamma$ is nonelementary if its action on $\partial \Gamma$ does not fix a finite set. Equivalently, $H$ is not virtually the trivial group or $\mathbb{Z}$. We say that a probability measure $\mu$ on $\Gamma$ is nonelementary if the subgroup $\Gamma_{\mu}$ generated by its support is itself nonelementary.

Let $\mu$ be a probability measure on $\Gamma$. Since $\Gamma$ acts by homeomorphisms on the compact space $\partial \Gamma$, it admits a stationary measure: there exists a probability measure $\nu$ on $\partial \Gamma$ such that $\mu * \nu=\nu$, i.e., $\sum_{g \in \Gamma} \mu(g) g_{*} \nu=\nu$. If $\mu$ is nonelementary, this measure is unique, and has no atom (see [Kai00]). It is also the exit measure of the corresponding random walk $X_{n}=g_{1} \cdots g_{n}$ : almost every trajectory $X_{n}(\omega)$ converges to a point $X_{\infty}(\omega) \in \partial \Gamma$, and moreover the distribution of $X_{\infty}$ is precisely $\nu$.

In the same way, since $\Gamma$ acts on $\partial_{B} \Gamma$, it admits a stationary measure $\nu_{B}$ there. This measure is not unique in general, even if $\mu$ is nonelementary. However, all such measures project under $\pi_{B}$ to the unique stationary measure on $\partial \Gamma$.
2.3. The drift. Let $(\Gamma, d)$ be a metric hyperbolic group. Consider a probability measure $\mu$ on $\Gamma$, with finite first moment $L(\mu)$ (defined in (1.1)). The drift of the random walk has been defined in (1.2) as $\ell(\mu)=\lim L\left(\mu^{* n}\right) / n$. Let $X_{n}=g_{1} \cdots g_{n}$ be the position at time $n$ of the random walk generated by $\mu$ (where the $g_{i}$ are independent and distributed according to $\mu$ ). Then, almost surely, $\ell(\mu)=\lim \left|X_{n}\right| / n$.

The drift also admits a description in terms of the Busemann boundary. The following result is well-known (compare with [KL11, Theorem 18]).

Proposition 2.2. Let $(\Gamma, d)$ be a metric hyperbolic group. Let $\mu$ be a nonelementary probability measure on $\Gamma$ with finite first moment. Let $\nu_{B}$ be a $\mu$-stationary measure on $\partial_{B} \Gamma$. Then

$$
\begin{equation*}
\ell(\mu)=\int_{\Gamma \times \partial_{B} \Gamma} c_{B}(g, \xi) \mathrm{d} \mu(g) \mathrm{d} \nu_{B}(\xi) \tag{2.3}
\end{equation*}
$$

Proof. Let $X_{n}$ be the position of the random walk at time $n$. Using the cocycle property of the Busemann cocycle, we have

$$
\begin{aligned}
\int c_{B}\left(X_{n}(\omega), \xi\right) \mathrm{d} \mathbb{P}(\omega) \mathrm{d} \nu_{B}(\xi) & =\int c_{B}\left(g_{1} \cdots g_{n}, \xi\right) \mathrm{d} \mu\left(g_{1}\right) \cdots \mathrm{d} \mu\left(g_{n}\right) \mathrm{d} \nu_{B}(\xi) \\
& =\sum_{k=1}^{n} \int c_{B}\left(g_{k}, g_{k+1} \cdots g_{n} \xi\right) \mathrm{d} \mu\left(g_{k}\right) \cdots \mathrm{d} \mu\left(g_{n}\right) \mathrm{d} \nu_{B}(\xi)
\end{aligned}
$$

Since the measure $\nu_{B}$ is stationary, the point $g_{k+1} \cdots g_{n} \xi$ is distributed according to $\nu_{B}$. Hence, the terms in the above sum do not depend on $k$. We get

$$
\begin{equation*}
\int_{\Gamma \times \partial_{B} \Gamma} c_{B}(g, \xi) \mathrm{d} \mu(g) \mathrm{d} \nu_{B}(\xi)=\frac{1}{n} \int c_{B}\left(X_{n}(\omega), \xi\right) \mathrm{d} \mathbb{P}(\omega) \mathrm{d} \nu_{B}(\xi) \tag{2.4}
\end{equation*}
$$

We have $\left|c_{B}\left(X_{n}, \xi\right)\right| / n \leqslant\left|X_{n}\right| / n$, which converges in $L^{1}$ and almost surely to $\ell$. Hence, the sequence of functions $c_{B}\left(X_{n}(\omega), \xi\right) / n$ is uniformly integrable on $\Omega \times \partial_{B} \Gamma$. Moreover, $X_{n}$ converges almost surely to a point on the boundary $\partial \Gamma$, distributed according to the exit measure, which has no atom. It follows that, for all $\xi$, the trajectory $X_{n}(\omega)$ converges almost surely to a point different from $\pi_{B}(\xi)$. This implies that, almost surely, one has $c_{B}\left(X_{n}, \xi\right)=\left|X_{n}\right|+O(1)$, giving in particular $c_{B}\left(X_{n}, \xi\right) / n \rightarrow \ell$. The result follows by taking the limit in $n$ in the equality (2.4).

This formula easily implies that the drift depends continuously on the measure, as explained in [EK13].

Proposition 2.3. Let $(\Gamma, d)$ be a metric hyperbolic group. Consider a sequence of probability measures $\mu_{i}$ with finite first moment, converging simply to a nonelementary probability measure $\mu$ (i.e., $\mu_{i}(g) \rightarrow \mu(g)$ for all $g \in \Gamma$ ). Assume moreover that $L\left(\mu_{i}\right) \rightarrow L(\mu)$. Then $\ell\left(\mu_{i}\right) \rightarrow \ell(\mu)$.

Proof. Let $\nu_{i}$ be stationary measures for $\mu_{i}$ on $\partial_{B} \Gamma$. Taking a subsequence if necessary, we may assume that $\nu_{i}$ converges to a limiting measure $\nu$. By continuity of the action on the boundary, it is stationary for $\mu$.

For each $g \in \Gamma$, the quantity $\int_{\partial_{B} \Gamma} c_{B}(g, \xi) \mathrm{d} \nu_{i}(\xi)$ converges to $\int_{\partial_{B} \Gamma} c_{B}(g, \xi) \mathrm{d} \nu(\xi)$ since $\xi \mapsto c_{B}(g, \xi)$ is continuous. Averaging over $g$ (and using the assumption $L\left(\mu_{i}\right) \rightarrow L(\mu)$ to get a uniform domination), we deduce that

$$
\sum_{g \in \Gamma} \mu_{i}(g) \int_{\partial_{B} \Gamma} c_{B}(g, \xi) \mathrm{d} \nu_{i}(\xi) \rightarrow \sum_{g \in \Gamma} \mu(g) \int_{\partial_{B} \Gamma} c_{B}(g, \xi) \mathrm{d} \nu(\xi)
$$

Together with the formula (2.3) for the drift, this completes the proof.
In this proposition, it is important that $\mu$ is nonelementary: the result is wrong otherwise. For instance, in the infinite dihedral group $\mathbb{Z} \rtimes \mathbb{Z} / 2$, the measures $\mu_{i}=\left(1-2^{-i}\right) \delta_{(1,0)}+$ $2^{-i} \delta_{(0,1)}$ have zero drift since the $\mathbb{Z} / 2$ element symmetrizes everything in $\mathbb{Z}$, while the limiting measure $\mu=\delta_{(1,0)}$ has drift 1 . The reason is the non-uniqueness of the stationary measure for $\mu$ on the boundary.
2.4. The entropy. Let $\Gamma$ be a countable group. Consider a probability measure $\mu$ on $\Gamma$, with finite time one entropy $H(\mu)$ (defined in (1.1)). The entropy of the random walk has been defined in (1.2) as $h(\mu)=\lim H\left(\mu^{* n}\right) / n$. Let $X_{n}=g_{1} \cdots g_{n}$ be the position at time $n$ of the random walk generated by $\mu$ (where the $g_{i}$ are independent and distributed according to $\mu$ ). Then, almost surely, $h(\mu)=\lim \left(-\log \mu^{* n}\left(X_{n}\right)\right) / n$. The fundamental inequality (1.3) shows that if $h>0$ then $\ell>0$.

The entropy has several equivalent characterizations. The first one is in terms of the size of the typical support of the random walk: This support has size roughly $e^{h n}$. The following lemma follows from [Haï13, Proposition 1.13].

Lemma 2.4. Consider a probability measure $\mu$ with $H(\mu)<\infty$ on a countable group. Let $h=h(\mu)$ be its asymptotic entropy. Let $\eta>0$ and $\varepsilon>0$.
(1) For large enough $n$, there exists a subset $K_{n}$ of $\Gamma$ with $\mu^{* n}\left(K_{n}\right) \geqslant 1-\eta$ and $\left|K_{n}\right| \leqslant$ $e^{(h+\varepsilon) n}$.
(2) For large enough $n$, there exists no subset $K_{n}$ of $\Gamma$ with $\mu^{* n}\left(K_{n}\right) \geqslant \eta$ and $\left|K_{n}\right| \leqslant$ $e^{(h-\varepsilon) n}$.

Another description is in terms of the Poisson boundary of the walk. To avoid general definitions, let us only state this description for measures on hyperbolic groups. The following proposition is a consequence of [Kai00].

Proposition 2.5. Let $\Gamma$ be a hyperbolic group. Let $\mu$ be a nonelementary probability measure on $\Gamma$ with $H(\mu)<\infty$. Let $\nu$ be its unique stationary measure on $\partial \Gamma$. Define the Martin cocycle on $\Gamma \times \partial \Gamma$ by $c_{M}(g, \xi)=-\log \left(\mathrm{d} g_{*}^{-1} \nu / \mathrm{d} \nu\right)(\xi)$. Then

$$
\begin{equation*}
h(\mu) \geqslant \int_{\Gamma \times \partial \Gamma} c_{M}(g, \xi) \mathrm{d} \mu(g) \mathrm{d} \nu(\xi), \tag{2.5}
\end{equation*}
$$

with equality if $\mu$ has a logarithmic moment.
When $\mu$ has a logarithmic moment, this proposition has a very similar flavor to Proposition 2.2 expressing the drift of a random walk. Indeed, for symmetric measures, [BHM11] interprets Proposition 2.5 as a special case of Proposition 2.2, for a distance $d=d_{\mu}$ related to the random walk, the Green distance, which we defined in Theorem 1.2. This distance is hyperbolic if $\mu$ is admissible and has a superexponential moment, by [Anc87, Gou13]. It is not geodesic in general, but this is not an issue since we were careful enough to state Proposition 2.2 without this assumption. The Busemann cocycle for the Green distance is precisely the Martin cocycle.

An important difference between the formulas (2.3) for the drift and (2.5) for the entropy is that, in the latter situation, the cocycle $c_{M}$ depends on the measure $\nu$ (and, therefore, on $\mu)$. This makes it more complicated to prove continuity statements such as Proposition 2.3 for the entropy. Nevertheless, Erschler and Kaimanovich proved in [EK13] that, in hyperbolic groups, the entropy also depends continuously on the measure. As $h(\mu)=\inf H\left(\mu^{* n}\right) / n$ by subadditivity, it is easy to prove that when $\mu_{i} \rightarrow \mu$ one has $\lim \sup h\left(\mu_{i}\right) \leqslant h(\mu)$. The main difficulty to prove the continuity is to get lower bounds. We will need a slightly stronger (and more pedestrian) version of the results of [EK13] to prove Theorem 1.4. Although our argument may seem very different at first sight from the arguments in [EK13], the techniques are in fact closely related (an illustration is that we can recover with our
techniques the result of Kaimanovich that, for measures with finite logarithmic moment, equality holds in (2.5), i.e., the Poisson boundary coincides with the geometric boundary, see Remark 2.11). Our main criterion to get lower bounds on the entropy is the following. We write $\mathbb{S}^{k}=\{g \in \Gamma:|g| \in(k-1, k]\}$ for the thickened sphere, so that the union of these spheres covers the whole group.

Theorem 2.6. Let $(\Gamma, d)$ be a metric hyperbolic group. Let $\mu_{i}$ be a sequence of nonelementary probability measures on $\Gamma$ with $H\left(\mu_{i}\right)<\infty$. Let $\nu_{i}$ be the unique stationary measure for $\mu_{i}$ on $\partial \Gamma$. Assume that:
(1) The limit points of $\nu_{i}$ have no atom.
(2) The sequence

$$
\begin{equation*}
h_{i}=\sum_{k} \sum_{g \in \mathbb{S}^{k}} \mu_{i}(g)\left(-\log \left(\mu_{i}(g) / \mu_{i}\left(\mathbb{S}^{k}\right)\right)\right) \tag{2.6}
\end{equation*}
$$

tends to infinity.
Then $\lim \inf h\left(\mu_{i}\right) / h_{i} \geqslant 1$.
The quantity $h_{i}$ can be written

$$
h_{i}=\sum_{g \in \Gamma} \mu_{i}(g)\left(-\log \mu_{i}(g)\right)-\sum_{k} \mu_{i}\left(\mathbb{S}^{k}\right)\left(-\log \mu_{i}\left(\mathbb{S}^{k}\right)\right)
$$

The first term is the time one entropy $H\left(\mu_{i}\right)$ of the measure $\mu_{i}$. In most reasonable cases, the second term is negligible. The theorem then states that the asymptotic entropy $h\left(\mu_{i}\right)$ is comparable to the time one entropy $H\left(\mu_{i}\right)$. In other words, if the measure is supported close to infinity, and sufficiently spread out in the group (this is the meaning of the assumption that the limit points of $\nu_{i}$ have no atom), then there are few coincidences and the entropy does not decrease significantly with time.

To prove this theorem, we will use the following technical lemma.
Lemma 2.7. On a probability space $(X, \mu)$, consider a nonnegative function $f$ with average 1. For any subset $A$ of $X$,

$$
\int_{X}(-\log f) \geqslant \mu(A)\left(-\log \int_{A} f\right)-2 e^{-1}
$$

Proof. As the function $x \mapsto-\log x$ is convex, Jensen's inequality gives $\int(-\log f) \geqslant$ $-\log \left(\int f\right)$. The last quantity vanishes when $\int f=1$.

Let $B \subset X$. Write $a=\int_{B} f \mathrm{~d} \mu / \mu(B)$. The measure $\mathrm{d} \mu / \mu(B)$ is a probability measure on $B$, and the function $f / a$ has integral 1 for this measure. The previous inequality gives $\int_{B}(-\log (f / a)) \mathrm{d} \mu / \mu(B) \geqslant 0$, that is,

$$
\int_{B}(-\log f) \mathrm{d} \mu \geqslant-\mu(B) \log a=-\mu(B) \log \left(\int_{B} f\right)+\mu(B) \log \mu(B)
$$

The quantity $\mu(B) \log \mu(B)$ is bounded from below by $\inf _{[0,1]} x \log x=-e^{-1}$. Therefore,

$$
\int_{B}(-\log f) \mathrm{d} \mu \geqslant-\mu(B) \log \left(\int_{B} f\right)-e^{-1}
$$

We apply this inequality to the complement $A^{c}$ of $A$. As $-\log \left(\int_{A^{c}} f\right) \geqslant 0$, we get a lower bound $-e^{-1}$. Let us also apply this inequality to $A$, and add the results. We obtain

$$
\int_{X}(-\log f) \mathrm{d} \mu \geqslant-\mu(A) \log \left(\int_{A} f\right)-2 e^{-1}
$$

We will use the notion of shadow, due to Sullivan and considered in this context by Coornaert [Coo93]. Let $C>0$ be large enough. The shadow $\mathcal{O}(g, C)$ of $g \in \Gamma$ is $\{\xi \in \partial \Gamma$ : $\left.(g \mid \xi)_{e} \geqslant|g|-C\right\}$. In geometric terms (and assuming the space is geodesic), this is essentially the trace at infinity of geodesics originating from $e$ and going through the ball $B(g, C)$. We will use the following properties of shadows [Coo93]:
(1) Their covering number is finite. More precisely, there exists $D>0$ (depending on $C)$ such that, for any integer $k$, for any $\xi \in \partial \Gamma$,

$$
\left|\left\{g \in \mathbb{S}^{k}: \xi \in \mathcal{O}(g, C)\right\}\right| \leqslant D
$$

(2) The preimages of shadows are large. More precisely, for any $\eta>0$, there exists $C>0$ such that, for all $g \in \Gamma$, the complement of $g^{-1} \mathcal{O}(g, C)$ has diameter at most $\eta$ (for a fixed visual distance on the boundary).

Proof of Theorem 2.6. Fix $\varepsilon>0$. As the limit points of $\nu_{i}$ have no atom, there exists $\eta>0$ such that any ball of radius $\eta$ in $\partial \Gamma$ has measure at most $\varepsilon$ for $\nu_{i}$, for $i$ large enough. We can then choose a shadow size $C$ so that $g^{-1} \mathcal{O}(g, C)$ has for all $g$ a complement with diameter at most $\eta$. This yields $\nu_{i}\left(g^{-1} \mathcal{O}(g, C)\right) \geqslant 1-\varepsilon$.

By (2.5), the entropy of $\mu_{i}$ satisfies

$$
h\left(\mu_{i}\right) \geqslant \sum_{g \in \Gamma} \mu_{i}(g) \int_{\partial \Gamma}\left(-\log \frac{\mathrm{d} g_{*}^{-1} \nu_{i}}{\mathrm{~d} \nu_{i}}(\xi)\right) \mathrm{d} \nu_{i}(\xi)
$$

The function $f_{i, g}=\frac{\mathrm{d} g_{*}^{-1} \nu_{i}}{\mathrm{~d} \nu_{i}}(\xi)$ is nonnegative and has integral 1. For any $A \subset \partial \Gamma$, Lemma 2.7 gives

$$
\begin{aligned}
\int_{\partial \Gamma}\left(-\log \frac{\mathrm{d} g_{*}^{-1} \nu_{i}}{\mathrm{~d} \nu_{i}}(\xi)\right) \mathrm{d} \nu_{i}(\xi) & \geqslant-\nu_{i}(A) \log \left(\int_{A} \frac{\mathrm{~d} g_{*}^{-1} \nu_{i}}{\mathrm{~d} \nu_{i}}(\xi) \mathrm{d} \nu_{i}(\xi)\right)-2 e^{-1} \\
& =-\nu_{i}(A) \log \left(g_{*}^{-1} \nu_{i}(A)\right)-2 e^{-1} \\
& =-\nu_{i}(A) \log \left(\nu_{i}(g A)\right)-2 e^{-1}
\end{aligned}
$$

Let us take $A=g^{-1} \mathcal{O}(g, C)$, so that $\nu_{i}(A) \geqslant 1-\varepsilon$. Summing over $g$, we get

$$
\begin{equation*}
h\left(\mu_{i}\right) \geqslant(1-\varepsilon) \sum_{g \in \Gamma} \mu_{i}(g)\left(-\log \nu_{i}(\mathcal{O}(g, C))\right)-2 e^{-1} \tag{2.7}
\end{equation*}
$$

We split the sum according to the spheres $\mathbb{S}^{k}$. Let $\Sigma_{k}=\sum_{g \in \mathbb{S}^{k}} \nu_{i}(\mathcal{O}(g, C))$, it is at most $D$ since the shadows have a covering number bounded by $D$. We have

$$
\begin{aligned}
\sum_{g \in \mathbb{S}^{k}} \mu_{i}(g) & \left(-\log \nu_{i}(\mathcal{O}(g, C))\right) \\
& =-\mu_{i}\left(\mathbb{S}^{k}\right) \sum_{g \in \mathbb{S}^{k}} \frac{\mu_{i}(g)}{\mu_{i}\left(\mathbb{S}^{k}\right)}\left[\log \left(\frac{\nu_{i}(\mathcal{O}(g, C))}{\Sigma_{k} \mu_{i}(g) / \mu_{i}\left(\mathbb{S}^{k}\right)}\right)+\log \Sigma_{k}+\log \left(\mu_{i}(g) / \mu_{i}\left(\mathbb{S}^{k}\right)\right)\right]
\end{aligned}
$$

The point of this decomposition is that the function on $\mathbb{S}^{k}$ given by $\varphi: g \mapsto \frac{\nu_{i}(\mathcal{O}(g, C))}{\Sigma_{k} \mu_{i}(g) / \mu_{i}\left(\mathbb{S}^{k}\right)}$ has integral 1 for the probability measure $\mu_{i}(g) / \mu_{i}\left(\mathbb{S}^{k}\right)$. By Jensen's inequality, the integral of $-\log \varphi$ is nonnegative. This yields

$$
\sum_{g \in \mathbb{S}^{k}} \mu_{i}(g)\left(-\log \nu_{i}(\mathcal{O}(g, C))\right) \geqslant-\mu_{i}\left(\mathbb{S}^{k}\right) \log D+\sum_{g \in \mathbb{S}^{k}} \mu_{i}(g)\left(-\log \left(\mu_{i}(g) / \mu_{i}\left(\mathbb{S}^{k}\right)\right) .\right.
$$

Summing over $k$, we deduce from (2.7) the inequality

$$
h\left(\mu_{i}\right) \geqslant(1-\varepsilon) h_{i}-2 e^{-1}-\log D
$$

As $h_{i}$ tends to infinity, this gives $h\left(\mu_{i}\right) \geqslant(1-2 \varepsilon) h_{i}$ for large enough $i$, completing the proof.

To apply the previous theorem, we need to estimate $h_{i}$. In this respect, the following lemma is often useful.

Lemma 2.8. Let $R_{i} \geqslant 1$. The quantity $h_{i}$ defined in (2.6) satisfies

$$
h_{i} \geqslant \sum_{|g| \leqslant R_{i}} \mu_{i}(g)\left(-\log \mu_{i}(g)\right)-\log \left(2+R_{i}\right) .
$$

Proof. In the definition of $h_{i}$, all the terms are nonnegative. Restricting the sum to those $g$ with $|g| \leqslant R_{i}$, we get

$$
\begin{aligned}
h_{i} & \geqslant \sum_{k \leqslant R_{i}} \sum_{g \in \mathbb{S}^{k}} \mu_{i}(g)\left(-\log \left(\mu_{i}(g) / \mu_{i}\left(\mathbb{S}^{k}\right)\right)\right) \\
& =\sum_{|g| \leqslant R_{i}} \mu_{i}(g)\left(-\log \mu_{i}(g)\right)-\sum_{k \leqslant R_{i}} \mu_{i}\left(\mathbb{S}^{k}\right)\left(-\log \mu_{i}\left(\mathbb{S}^{k}\right)\right) .
\end{aligned}
$$

A probability measure supported on a set with $N$ elements has entropy at most $\log N$. The number $\mu_{i}\left(\mathbb{S}^{k}\right)$ for $0 \leqslant k \leqslant R_{i}$ are not a probability measure in general, let us add a last atom with mass $m=\mu_{i}\left(\bigcup_{k>R_{i}} \mathbb{S}^{k}\right)$. We are considering a space of cardinality $R_{n}+2$, hence

$$
m(-\log m)+\sum_{k \leqslant R_{i}} \mu_{i}\left(\mathbb{S}^{k}\right)\left(-\log \mu_{i}\left(\mathbb{S}^{k}\right)\right) \leqslant \log \left(2+R_{i}\right),
$$

completing the proof.
Let us see how Theorem 2.6 implies a slightly stronger version of the continuity result for the entropy of Erschler and Kaimanovich [EK13].

Theorem 2.9. Let $\Gamma$ be a hyperbolic group. Consider a probability measure $\mu$ with finite time one entropy and finite logarithmic moment. Let $\mu_{i}$ be a sequence of probability measures converging simply to $\mu$ with $H\left(\mu_{i}\right) \rightarrow H(\mu)$. Then $h\left(\mu_{i}\right) \rightarrow h(\mu)$.

The assumption $H\left(\mu_{i}\right) \rightarrow H(\mu)$ ensures that there is no additional entropy in $\mu_{i}$ coming from neighborhoods of infinity that would disappear in the limit. It is automatic if the support of $\mu_{i}$ is uniformly bounded or if $\mu_{i}$ satisfies a uniform $L^{1}$ domination, but it is much weaker. For instance, it is allowed that the $\mu_{i}$ have no finite logarithmic moment.

The main lemma for the proof is a lower bound on the entropy, following from Theorem 2.6.

Lemma 2.10. Let $\Gamma$ be a hyperbolic group. Consider a probability measure $\mu$ with finite time one entropy and finite logarithmic moment. Let $\mu_{i}$ be a sequence of measures converging simply to $\mu$. Then $\lim \inf h\left(\mu_{i}\right) \geqslant h(\mu)$.
Proof. Since the result is trivial if $h(\mu)=0$, we can assume that $h(\mu)>0$.
Let $\varepsilon>0$. For large $n$, most atoms for $\mu^{* n}$ have a probability at most $e^{-(1-\varepsilon) n h(\mu)}$. Moreover, since $\mu$ has a finite logarithmic moment, $\log \left|X_{n}\right| / n$ tends almost surely to 0 by [Aar97, Proposition 2.3.1]. Therefore, the set

$$
K_{n}=\left\{g: \mu^{* n}(g) \leqslant e^{-(1-\varepsilon) n h(\mu)},|g| \leqslant e^{\varepsilon n}\right\}
$$

has measure tending to 1 . In particular $\mu^{* n}\left(K_{n}\right) \geqslant 1-\varepsilon$ for large $n$. We get

$$
\begin{aligned}
\sum_{|g| \leqslant \varepsilon^{\varepsilon n}} \mu^{* n}(g)\left(-\log \mu^{* n}(g)\right) & \geqslant \sum_{g \in K_{n}} \mu^{* n}(g)\left(-\log \mu^{* n}(g)\right) \geqslant \sum_{g \in K_{n}} \mu^{* n}(g)(1-\varepsilon) n h(\mu) \\
& =\mu^{* n}\left(K_{n}\right)(1-\varepsilon) n h(\mu) \geqslant(1-\varepsilon)^{2} n h(\mu) .
\end{aligned}
$$

For each fixed $n$, the measures $\mu_{i}^{* n}$ converge to $\mu^{* n}$ when $i$ tends to infinity. Hence, we get for large enough $i$ the inequality

$$
\sum_{|g| \leqslant e^{\varepsilon n}} \mu_{i}^{* n}(g)\left(-\log \mu_{i}^{* n}(g)\right) \geqslant(1-\varepsilon)^{3} n h(\mu) .
$$

Letting $\varepsilon$ tend to 0 (and, therefore, $n$ to infinity), we deduce the existence of sequences $n_{i} \rightarrow \infty$ and $\varepsilon_{i} \rightarrow 0$ such that, for any $i$,

$$
\sum_{|g| \leqslant e^{\varepsilon_{i} n_{i}}} \mu_{i}^{* n_{i}}(g)\left(-\log \mu_{i}^{* n_{i}}(g)\right) \geqslant\left(1-\varepsilon_{i}\right)^{3} n_{i} h(\mu) .
$$

Let $\tilde{\mu}_{i}=\mu_{i}^{* n_{i}}$. Its stationary measure $\nu_{i}$ is also the stationary measure of $\mu_{i}$, by uniqueness. Any limit point of $\nu_{i}$ is stationary for $\mu$, and is therefore atomless since $\mu$ is nonelementary as $h(\mu)>0$. The assumptions of Theorem 2.6 are satisfied by the sequence $\tilde{\mu}_{i}$. Moreover, Lemma 2.8 yields

$$
h_{i} \geqslant\left(1-\varepsilon_{i}\right)^{3} n_{i} h(\mu)-2 \varepsilon_{i} n_{i} \geqslant\left(1-C \varepsilon_{i}\right) n_{i} h(\mu) .
$$

Theorem 2.6 ensures that $\lim \inf h\left(\tilde{\mu}_{i}\right) / h_{i} \geqslant 1$. As $h\left(\tilde{\mu}_{i}\right)=n_{i} h\left(\mu_{i}\right)$, this gives $\lim \inf h\left(\mu_{i}\right) \geqslant$ $h(\mu)$ as desired.

Proof of Theorem 2.9. For fixed $n$, the sequence $\mu_{i}^{* n}$ converges simply to $\mu^{* n}$. Moreover, $H\left(\mu_{i}^{* n}\right) \rightarrow H\left(\mu^{* n}\right)$ since there is no loss of entropy at infinity by assumption. Choose $n$ such that $H\left(\mu^{* n}\right) \leqslant n(1+\varepsilon) h(\mu)$. We get $H\left(\mu_{i}^{* n}\right) / n \leqslant(1+2 \varepsilon) h(\mu)$ for large enough $i$. As $h\left(\mu_{i}\right) \leqslant H\left(\mu_{i}^{* n}\right) / n$, this shows that $\lim \sup h\left(\mu_{i}\right) \leqslant h(\mu)$ (this is the classical semi-continuity property of entropy, valid in any group).

For the reverse inequality $\lim \inf h\left(\mu_{i}\right) \geqslant h(\mu)$, we apply Lemma 2.10.
Remark 2.11. Let $h(\mu, \partial \Gamma)=\int_{\Gamma \times \partial \Gamma}\left(-\log \mathrm{d} g_{*}^{-1} \nu / \mathrm{d} \nu\right)(\xi) \mathrm{d} \mu(g) \mathrm{d} \nu(\xi)$ where $\nu$ is the stationary measure for $\mu$ on $\partial \Gamma$. In general, $h(\mu) \geqslant h(\mu, \partial \Gamma)$ with equality if and only if ( $\partial \Gamma, \nu$ ) is the Poisson boundary of $(\Gamma, \mu)$. A theorem of Kaimanovich [Kai00] asserts that, when $\mu$ has finite entropy and finite logarithmic moment, $h(\mu, \partial \Gamma)=h(\mu)$. We can recover this theorem using the previous arguments. Indeed, what the proof of Theorem 2.6 really shows is that $\lim \inf h\left(\mu_{i}, \partial \Gamma\right) / h_{i} \geqslant 1$. Hence, Lemma 2.10 proves that $\liminf h\left(\mu_{i}, \partial \Gamma\right) \geqslant h(\mu)$ if $\mu_{i}$ converges simply to a measure $\mu$ with a logarithmic moment. Taking $\mu_{i}=\mu$ for all $i$, we obtain in particular $h(\mu, \partial \Gamma) \geqslant h(\mu)$, as desired.
2.5. A criterion to bound the entropy from below. In order to prove Theorem 1.4 on the entropy of the uniform measure on balls, we want to apply Theorem 2.6. Thus, we need a criterion to check that limit points of stationary measures have no atom.

Lemma 2.12. Let $\Gamma$ be a hyperbolic group. Let $\mu_{i}$ be a sequence of probability measures on $\Gamma$. Assume that, on the space $\Gamma \cup \partial \Gamma$, the sequence $\mu_{i}$ converges to a limit $\nu$ which is supported on $\partial \Gamma$. Assume moreover that the limit points of $\check{\mu}_{i}$ (defined by $\check{\mu}_{i}(g)=\mu_{i}\left(g^{-1}\right)$ ) have no atom. Then the stationary measures $\nu_{i}$ associated to $\mu_{i}$ also converge to $\nu$.

Proof. We fix a word distance $d$ on $\Gamma$. Let $f$ be a continuous function on $\Gamma \cup \partial \Gamma$. Let us show that, uniformly in $\xi \in \partial \Gamma$, the integral $\int f(g \xi) \mathrm{d} \mu_{i}(g)$ is close to $\int f(g) \mathrm{d} \mu_{i}(g)$. We estimate the difference as

$$
\begin{aligned}
\left|\int(f(g \xi)-f(g)) \mathrm{d} \mu_{i}(g)\right| \leqslant & \int|f(g \xi)-f(g)| 1\left((g \xi \mid g)_{e}>C\right) \mathrm{d} \mu_{i}(g) \\
& +2\|f\|_{\infty} \int 1\left((g \xi \mid g)_{e} \leqslant C\right) \mathrm{d} \mu_{i}(g)
\end{aligned}
$$

where $C$ is a fixed constant. If $C$ is large enough, $|f(x)-f(y)| \leqslant \varepsilon$ when $(x \mid y)_{e}>C$, by uniform continuity of $f$. Hence, the first integral is bounded by $\varepsilon$. For the second integral, we use the formula $(g x \mid g)_{e}=|g|-\left(x \mid g^{-1}\right)_{e}$, valid for any $x \in \Gamma$ (it follows readily from the definition (2.1) of the Gromov product). This equality does not extend to the boundary since the Gromov product there is only well defined up to an additive constant $D$. Nevertheless, we get $(g \xi \mid g)_{e} \geqslant|g|-\left(\xi \mid g^{-1}\right)_{e}-D$. Hence, the second integral is bounded by

$$
\begin{equation*}
\mu_{i}\left\{g:|g|-C-D \leqslant\left(\xi \mid g^{-1}\right)_{e}\right\} . \tag{2.8}
\end{equation*}
$$

If $|g|$ is large, the points $g$ with $\left(\xi \mid g^{-1}\right)_{e} \geqslant|g|-C-D$ are such that $g^{-1}$ belongs to a small neighborhood of $\xi$ in $\Gamma \cup \partial \Gamma$. As the limit points of $\check{\mu}_{i}$ are supported on $\partial \Gamma$ and have no atom, it follows that (2.8) converges to 0 when $i$ tends to infinity, uniformly in $\xi$.

We have proved that

$$
\sup _{\xi \in \partial \Gamma}\left|\int f(g \xi) \mathrm{d} \mu_{i}(g)-\int f(g) \mathrm{d} \mu_{i}(g)\right| \rightarrow 0
$$

By stationarity,

$$
\int_{\xi \in \partial \Gamma} f(\xi) \mathrm{d} \nu_{i}(\xi)=\int_{\xi \in \partial \Gamma}\left(\int f(g \xi) \mathrm{d} \mu_{i}(g)\right) \mathrm{d} \nu_{i}(\xi)
$$

Combining these equations, we get $\int f(\xi) \mathrm{d} \nu_{i}(\xi)-\int f(g) \mathrm{d} \mu_{i}(g) \rightarrow 0$. This shows that the limit points of $\nu_{i}$ and $\mu_{i}$ are the same.

Let us now consider the uniform measure $\mu_{i}$ on the ball of radius $i$, as in Theorem 1.4. The next lemma follows from the techniques of [Coo93].

Lemma 2.13. Let $(\Gamma, d)$ be a metric hyperbolic group. Let $\rho_{i}$ be the uniform measure on the ball of radius $i$. Let $\rho_{\infty}$ be the Patterson-Sullivan of $(\Gamma, d)$ constructed in [Coo93] (it is supported on $\partial \Gamma$ and atomless). Then the limit points of $\rho_{i}$ are equivalent to $\rho_{\infty}$, with a density bounded from above and from below.

Proof. Let $C$ be large enough. We will use the shadows $\mathcal{O}(g, C)$ as defined before the proof of Theorem 2.6. The main property of $\rho_{\infty}$ is that it satisfies

$$
\begin{equation*}
K_{0}^{-1} e^{-v|g|} \leqslant \rho_{\infty}(\mathcal{O}(g, C)) \leqslant K_{0} e^{-v|g|} \tag{2.9}
\end{equation*}
$$

where $K_{0}$ is a constant only depending on $C$ and $v$ is the growth of ( $\Gamma, d$ ) (Proposition 6.1 in [Coo93]).

Let $\mu_{i}$ be the uniform measure on thickened spheres $S_{i}=\{g: i \leqslant|g| \leqslant i+L\}$, where $L$ is large enough so that the cardinality of $S_{i}$ grows like $e^{i v}$, see the proof of Theorem 7.2 in [Coo93]. Let us push $\mu_{i}$ to a measure $\tilde{\mu}_{i}$ on $\partial \Gamma$, by choosing for each $g \in S_{i}$ a corresponding point in its shadow. It is clear that $\mu_{i}$ and $\tilde{\mu}_{i}$ have the same limit points, since the diameter of the shadows tends uniformly to 0 when $i \rightarrow \infty$. We will prove that the limit points of $\tilde{\mu}_{i}$ are equivalent to $\rho_{\infty}$. The same result follows for $\mu_{i}$ and then $\rho_{i}$.

The shadows of $g \in S_{i}$ have a covering number which is bounded from above by a constant $D$, and from below by 1 if $C$ is large enough. Hence, the measures $\tilde{\mu}_{i}$ satisfy

$$
K_{1}^{-1} e^{-i v} \leqslant \tilde{\mu}_{i}(\mathcal{O}(g, C)) \leqslant K_{1} e^{-i v}
$$

for any $g \in S_{i}$. This is comparable to $\rho_{\infty}(\mathcal{O}(g, C))$ by (2.9), up to a multiplicative constant $K_{2}$. Consider a limit $\tilde{\mu}$ of a sequence $\tilde{\mu}_{i_{n}}$, let us prove that it is uniformly equivalent to $\rho_{\infty}$. We will only prove that $\tilde{\mu} \leqslant D K_{2} \rho_{\infty}$, the other inequality is proved in the same way. By regularity of the measures, it suffices to check this inequality on compact sets.

Let $A$ be a compact subset of $\partial \Gamma$, and $\varepsilon>0$. By regularity of the measure $\rho_{\infty}$, there is an open neighborhood $U$ of $A$ with $\rho_{\infty}(U) \leqslant \rho_{\infty}(A)+\varepsilon$. Consider $B$ a compact neighborhood of $A$, included in $U$, with $\tilde{\mu}(\partial B)=0$ (such a set exists, since among the sets $B_{r}=\{\xi$ : $d(\xi, A) \leqslant r\}$, at most countably of them many have a boundary with nonzero measure). For large enough $i$, the shadows $\mathcal{O}(g, C)$ with $g \in S_{i}$ which intersect $B$ are contained in $U$. Therefore,

$$
\tilde{\mu}_{i}(B) \leqslant \sum_{g \in S_{i}, \mathcal{O}(g, C) \cap B \neq \emptyset} \tilde{\mu}_{i}(\mathcal{O}(g, C)) \leqslant K_{2} \sum_{g \in S_{i}, \mathcal{O}(g, C) \cap B \neq \emptyset} \rho_{\infty}(\mathcal{O}(g, C)) \leqslant D K_{2} \rho_{\infty}(U)
$$

As $\tilde{\mu}(\partial B)=0$, the sequence $\tilde{\mu}_{i_{n}}(B)$ tends to $\tilde{\mu}(B)$. We obtain $\tilde{\mu}(B) \leqslant D K_{2} \rho_{\infty}(U)$. As $A$ is included in $B$, we get $\tilde{\mu}(A) \leqslant D K_{2}\left(\rho_{\infty}(A)+\varepsilon\right)$. Letting $\varepsilon$ tend to 0 , this gives $\tilde{\mu}(A) \leqslant D K_{2} \rho_{\infty}(A)$, as desired.

Proof of Theorem 1.4. Let $\rho_{i}$ be the uniform measure on the ball of radius $i$ (which has cardinality in $\left[C^{-1} e^{i v}, C e^{i v}\right]$ ). We wish to apply Theorem 2.6 to this sequence of measures. First, by Lemmas 2.12 and 2.13 , the limit points of the stationary measures $\nu_{i}$ are equivalent to the Patterson-Sullivan measure. Therefore, they have no atom. Second, Lemma 2.8 shows that the quantity $h_{i}$ in (2.6) satisfies $h_{i} \geqslant i v-\log C-\log (2+i)$. This tends to infinity. Hence, Theorem 2.6 applies, and gives $h\left(\rho_{i}\right) \geqslant(1-\varepsilon) i v$ for large $i$.

Using the fundamental inequality $h \leqslant \ell v$ and the trivial bound $\ell\left(\rho_{i}\right) \leqslant L\left(\rho_{i}\right) \leqslant i$, we get

$$
(1-\varepsilon) i v \leqslant h\left(\rho_{i}\right) \leqslant \ell\left(\rho_{i}\right) v \leqslant i v .
$$

It follows that $h\left(\rho_{i}\right) \sim i v$ and $\ell\left(\rho_{i}\right) \sim i$.
Remark 2.14. Our technique also applies to estimate the entropy of other measures, for instance the measure $\mu_{s}=\sum e^{-s|g|} \delta_{g} / \sum e^{-s|g|}$ classically used in the construction of the Patterson-Sullivan measure. Indeed, $\mu_{s}$ converges when $s \searrow v$ to $\rho_{\infty}$, which has no atom. Moreover, writing $Z_{s}=\sum e^{-s|g|}$, we have $H\left(\mu_{s}\right)=s L\left(\mu_{s}\right)+\log Z_{s}$. One checks that $\log Z_{s}$ is negligible with respect to $H\left(\mu_{s}\right)$, and that the quantity $h_{s}$ from (2.6) is also equivalent to $H\left(\mu_{s}\right)$. Hence, Theorem 2.6 gives

$$
H\left(\mu_{s}\right)(1+o(1)) \leqslant h_{s}(1+o(1)) \leqslant h\left(\mu_{s}\right) \leqslant \ell\left(\mu_{s}\right) v \leqslant L\left(\mu_{s}\right) v \leqslant H\left(\mu_{s}\right)(1+o(1)) .
$$

These inequalities show that $h\left(\mu_{s}\right) / \ell\left(\mu_{s}\right) \rightarrow v$.
Remark 2.15. One could imagine another strategy to find finitely supported measures $\mu_{i}$ for which $h\left(\mu_{i}\right) / \ell\left(\mu_{i}\right) \rightarrow v$. First, find a nice measure $\mu$ for which the stationary measure $\nu$ at infinity is precisely the Patterson-Sullivan measure (which implies that $h(\mu)=\ell(\mu) v$ since the Martin cocycle and the Busemann cocycle coincide). Let $\mu_{i}$ be a truncation of $\mu$. Since it converges to $\mu$, the continuity results for the drift and the entropy imply that $h\left(\mu_{i}\right) / \ell\left(\mu_{i}\right) \rightarrow h(\mu) / \ell(\mu)=v$.

We were not able to implement successfully this strategy. Given a measure $\nu$, there is a general technique due to Connell and Muchnik [CM07] to get a measure $\mu$ on $\Gamma$ with $\mu * \nu=\nu$. This technique requires a continuity assumption on $\xi \mapsto\left(\mathrm{d} g_{*} \nu / \mathrm{d} \nu\right)(\xi)$, which is not satisfied in our setting for $\nu=\rho_{\infty}$. However, in nice groups such as surface groups, this function is, for every $g$, continuous at all but finitely many points. The technique of [CM07] can be adapted to such a situation (in the proof of their Theorem 6.2, one should just take sets $Y_{n}$ that avoid the discontinuities of the spikes we have already used). Unfortunately, the resulting measure $\mu$ (which satisfies $\mu * \nu=\nu$ ) has infinite moment and infinite entropy, and is therefore useless for our purposes.

## 3. Rigidity for admissible measures

In this section, we prove Theorem 1.5. Assume that $(\Gamma, d)$ is a hyperbolic group endowed with a word distance, which is not virtually free. Let $\mu$ be a probability measure on $\Gamma$, with a superexponential moment, such that $\Gamma_{\mu}^{+}$is a finite index subgroup of $\Gamma$. We want to prove that $h(\mu)<\ell(\mu) v$. We argue by contradiction, assuming that $h(\mu)=\ell(\mu) v$. Assume first that $\Gamma_{\mu}^{+}=\Gamma$.

Since we are assuming the equality $h(\mu)=\ell(\mu) v$, Theorem 1.2 implies that

$$
\left|d_{\mu}(e, g)-v d(e, g)\right| \leqslant C .
$$

As a warm-up, let us first deal with the baby case $C=0$. Then the distances $d_{\mu}$ and $d$ are proportional, hence they define the same Busemann boundary. The Busemann boundary $\partial_{B} \Gamma$ corresponding to $d$ is totally discontinuous since the distance $d$ takes integer values (it is a word distance). On the other hand, the Busemann boundary associated to the Green metric $d_{\mu}$ is known as the Martin boundary of the random walk ( $\Gamma, \mu$ ). By [Anc87] and [Gou13], it is homeomorphic to the boundary $\partial \Gamma$ of $\Gamma$. Since the group $\Gamma$ is not virtually free, its boundary $\partial \Gamma$ is not totally discontinuous (see [KB02, Theorem 8.1]), hence a contradiction.

Let us now go back to the general situation, when $C$ is nonzero (but still assuming $\Gamma_{\mu}^{+}=\Gamma$ ). The argument is more complicated, but it still relies on the same facts: the boundary is not totally disconnected, while the word distance is integer valued (we will not use directly this fact, rather the fact that stable translation lengths are rational, see Lemma 3.4). These two opposite features will give rise to a contradiction.

In order to get rid of the constant $C$, we will need an homogenized version of the inequality $\left|d_{\mu}(e, g)-v d(e, g)\right| \leqslant C$. This is Lemma 3.1 below. The homogenized quantity associated to the distance $d$ is called the stable translation length. For an element $g$ of $\Gamma$, it is defined by $l(g)=\lim \left|g^{n}\right| / n$ (it exists by subadditivity).

Recall that we write $c_{M}(g, \xi)$ for the Martin cocycle associated to the random walk, defined in Proposition 2.5. It satisfies the cocycle relation of Definition 2.1. We will not use its probabilistic definition, but rather the fact that the Martin cocycle is the Busemann cocycle associated to the Green distance $d_{\mu}$ of Theorem 1.2. In other words, $c_{M}(g, \xi)=$ $\lim _{x \rightarrow \xi} d_{\mu}\left(g^{-1}, x\right)-d_{\mu}(e, x)$ (and this limit exists).

Lemma 3.1. For $g \in \Gamma$ with infinite order, $c_{M}\left(g, g^{+}\right)=v l(g)$.
Proof. Recall that we are assuming that the equality $h(\mu)=\ell(\mu) v$ holds, therefore we have $\left|d_{\mu}(e, g)-v d(e, g)\right| \leqslant C$. It follows that the cocycle $c_{M}$ corresponding to $d_{\mu}$ and the cocycle $c_{B}$ corresponding to the distance $d$ satisfy $\left|c_{M}-v c_{B}\right| \leqslant 2 C$. Note that $c_{B}$ is not defined on the geometric boundary, but on the horoboundary, so the proper way to write this inequality is $\left|c_{M}\left(g, \pi_{B}(\xi)\right)-v c_{B}(g, \xi)\right| \leqslant 2 C$ for any $g \in \Gamma$ and any $\xi \in \partial_{B} \Gamma$.

Let $\xi \in \partial_{B} \Gamma$ with $\pi_{B}(\xi) \neq g^{-}$. Then $\lim c_{B}\left(g^{n}, \xi\right) / n=\lim h_{\xi}\left(g^{-n}\right) / n=l(g)$. We choose $\xi$ with $\pi_{B}(\xi)=g^{+}$, to get

$$
\lim c_{M}\left(g^{n}, g^{+}\right) / n=\lim v c_{B}\left(g^{n}, \xi\right) / n \pm 2 C / n=v l(g)
$$

As $g^{+}$is $g$-invariant, the cocycle equation for $c_{M}$ on $\partial \Gamma$ gives $c_{M}\left(g, g^{+}\right)=c_{M}\left(g^{n}, g^{+}\right) / n$. This converges to $v l(g)$ when $n \rightarrow \infty$ by the previous equation.

The proof of Theorem 1.5 uses the following general result on cocycles.
Proposition 3.2. Let $\Gamma$ be a hyperbolic group which is not virtually free. Let $c: \Gamma \times \partial \Gamma \rightarrow \mathbb{R}$ be a Hölder cocycle, such that any hyperbolic element $g$ satisfies $c\left(g, g^{+}\right) \in \mathbb{Z}$. Then there exists a hyperbolic element $g \in \Gamma$ with $c\left(g, g^{-}\right)=c\left(g, g^{+}\right)$.

Applied to the Busemann cocycle, this proposition implies that if a convex cocompact negatively curved manifold has a fundamental group which is not virtually free, then its length spectrum is not arithmetic, i.e., the lengths of its closed geodesics generate a dense subgroup of $\mathbb{R}$. This result is already known, see [Dal99, Page 205]. It is proved in this article
using crossratios. This argument based on crossratios can be used to prove Proposition 3.2 in full generality. However, we will give a different, more direct, proof.

We will use the following topological lemma.
Lemma 3.3. Let $g$ be a hyperbolic element in a hyperbolic group $\Lambda$ with connected boundary. There exists an arc $I$ (i.e., a subset of $\partial \Lambda$ homeomorphic to $[0,1]$ ) joining $g^{-}$and $g^{+}$, invariant under an iterate $g^{i}$ of $g$.

Proof. We will use nontrivial results on the topology of $\partial \Lambda$. When it is connected, then it is also locally connected by [Swa96]. Hence, it is also path connected and locally path connected, see [HY61, Theorem 3-16]. Moreover, for any $\xi \in \partial \Lambda$, the space $\partial \Lambda \backslash\{\xi\}$ has finitely many ends by [Bow98b].

Consider $g$ as in the statement of the lemma. Its action permutes the ends of $\partial \Lambda \backslash\left\{g^{-}\right\}$. Taking an iterate of $g$, we can assume it stabilizes the ends. If $\xi$ is close to $g^{-}$, it is also the case of $g \xi$. As they belong to the same end, one can join them by a small arc $J$ that avoids $g^{-}$(and $\left.g^{+}\right)$. Then $\bigcup_{n \in \mathbb{Z}} g^{n} J$ joins $g^{-}$to $g^{+}$, and it is invariant under $g$. However, it is not necessarily an arc if $g^{i} J$ intersects $J$ in a nontrivial way for $i \neq 0$. To get a real arc, we will shorten $J$ as follows.

As $g^{n} J$ converges to $g^{ \pm}$when $n$ tends to $\pm \infty$, the arc $J$ can only intersect finitely many $g^{i} J$. Let us fix a parametrization $u:[0,1] \rightarrow J$. The quantity

$$
\inf \left\{|t-s|: s, t \in[0,1] \text { and } \exists i \neq 0, u(t)=g^{i} u(s)\right\}
$$

is realized by compactness (since $i$ remains bounded), for some parameters $s, t, i$. Replacing $s, t, i$ with $t, s,-i$ if necessary, we may assume $i>0$. As $g^{-}$and $g^{+}$are the only fixed points of $g^{i}$, we have $s \neq t$. Let $K=u([s, t])$, this is an arc between $\eta=u(s)$ and $g^{i} \eta=u(t)$. Moreover, $g^{j} K$ does not intersect $K$, except maybe at its endpoints for $j= \pm i$ : otherwise, there exists $x$ in the interior of $K$ such that $g^{j} x$ also belongs to $K$, contradicting the minimality of $|s-t|$.

It follows that $\bigcup_{n \in \mathbb{Z}} g^{n i} K$ is an arc from $g^{-}$to $g^{+}$, invariant under $g^{i}$.
Proof of Proposition 3.2. Let us consider the cocycle $\bar{c}=c \bmod \mathbb{Z}$. The assumption of the proposition ensures that $\bar{c}\left(g, g^{+}\right)=0$ for all hyperbolic elements $g$. In geometric terms, this would correspond to an assumption that the cocycle has vanishing average on all closed orbits. Hence, we may apply a version of Livsic's theorem, due in this context to [INO08] (Theorem 5.1). It ensures that the cocycle $\bar{c}$ is a coboundary: there exists a Hölder continuous function $\bar{b}: \partial \Gamma \rightarrow \mathbb{R} / \mathbb{Z}$ such that, for all $\xi \in \partial \Gamma$, for all $g \in \Gamma$,

$$
\begin{equation*}
\bar{c}(g, \xi)=\bar{b}(g \xi)-\bar{b}(\xi) \tag{3.1}
\end{equation*}
$$

Recall that, since the group $\Gamma$ is not virtually free, its boundary is not totally discontinuous (see [KB02, Theorem 8.1]). The stabilizer of a nontrivial component $L$ of $\partial \Gamma$ is a subgroup $\Lambda$ of $\Gamma$, quasi-convex hence hyperbolic, whose boundary is $L$ (see the discussion on top of Page 55 in [Bow98a]).

Let us consider an infinite order element $g \in \Lambda$. Lemma 3.3 constructs an arc $I$ from $g^{-}$ to $g^{+}$in $\partial \Lambda \subset \partial \Gamma$, invariant under an iterate $g^{i}$ of $g$. Replacing $g$ with $g^{i}$, we may assume $i=1$.

The restriction of the function $\bar{b}$ to the arc $I$ admits a continuous lift $b: I \rightarrow \mathbb{R}$, as $I$ is simply connected. The function $F: \xi \mapsto c(g, \xi)-b(g \xi)+b(\xi)$ is well defined on
$I$, continuous, and it vanishes modulo $\mathbb{Z}$ by (3.1). Hence, it is constant. In particular, $c\left(g, g^{-}\right)=F\left(g^{-}\right)=F\left(g^{+}\right)=c\left(g, g^{+}\right)$.

In order to apply Proposition 3.2, we will need the following result on stable translation lengths in hyperbolic groups ([BH99, Theorem III.Г.3.17]).

Lemma 3.4. Let $(\Gamma, d)$ be a hyperbolic group with a word distance. Then there exists an integer $N$ such that, for any $g \in \Gamma$, one has $N l(g) \in \mathbb{Z}$.

The combination of Lemma 3.1 and Lemma 3.4 shows that the cocycle $c^{\prime}=N c_{M} / v$ satisfies $c^{\prime}\left(g, g^{+}\right) \in \mathbb{Z}$ for any hyperbolic element $g$. Moreover, this cocycle is Hölder-continuous since the Martin cocycle $c_{M}$ is itself Hölder-continuous. This follows from [INO08] if $\mu$ has finite support, and from [Gou13] if it has a superexponential moment. Now, Proposition 3.2 implies the existence of a hyperbolic element $g$ such that $c_{M}\left(g, g^{+}\right)=c_{M}\left(g, g^{-}\right)$. This is a contradiction since $c\left(g, g^{+}\right)=v l(g)>0$ and $c\left(g, g^{-}\right)=-c\left(g^{-1}, g^{-}\right)=-v l(g)<0$ again by Lemma 3.1. This concludes the proof of Theorem 1.5 when $\Gamma_{\mu}^{+}=\Gamma$.

If $\Gamma_{\mu}^{+}$is a finite index subgroup of $\Gamma$, the same proof almost works in $\Gamma_{\mu}^{+}$to conclude that $\Gamma_{\mu}^{+}$is virtually free if $h=\ell v$, implying that $\Gamma$ is also virtually free. The only difficulty is that the distance we are considering on $\Gamma_{\mu}^{+}$is not a word distance for a system of generators of $\Gamma_{\mu}^{+}$. However, the only properties of the distance we have really used are:
(1) It is hyperbolic and quasi-isometric to a word distance (to apply Theorem 1.2).
(2) The stable translation lengths are rational numbers with bounded denominators.

These two properties are clearly satisfied for the restriction of the distance $d$ to $\Gamma_{\mu}^{+}$. Hence, the above proof also works in this case. This completes the proof of Theorem 1.5.

Remark 3.5. If $\Lambda$ is a quasi-convex subgroup of a hyperbolic group $\Gamma$, then the restriction to $\Lambda$ of a word distance on $\Gamma$ also satisfies the above two properties. Hence, Theorem 1.5 also holds in $\Lambda$ for such a distance.

## 4. Growth of non-distorted points in subgroups

Our goal in this section is to prove Theorem 1.6 on the entropy of a random walk on an infinite index subgroup $\Lambda$ of a hyperbolic group $\Gamma$. Since the geometry of such random walks is complicated to describe in general, our argument is indirect: we will show that, in any infinite index subgroup, the number of points that the random walk effectively visits is exponentially small compared to the growth of $\Gamma$. This is trivial if the growth $v_{\Lambda}=$ $\liminf _{n \rightarrow \infty} \frac{\log \left|B_{n} \cap \Lambda\right|}{n}$ is strictly smaller than $v=v_{\Gamma}$. When $v_{\Lambda}=v$, on the other hand, we will argue that the random walk does not typically visit all of $\Lambda$, but only a subset made of non-distorted points. To prove Theorem 1.6 , the main step is to show that, even when $v_{\Lambda}=v$, the number of such non-distorted points is exponentially smaller than $e^{n v}$. We introduce the notion of non-distorted points in Paragraph 4.1, prove this main geometric estimate in Paragraph 4.2, and apply this to random walks in Paragraph 4.4. Paragraph 4.3 is devoted to the case $v_{\Lambda}<v$, where unexpected phenomena happen even in distorted subgroups.
4.1. Non-distorted points. There are at least two different ways to define a notion of non-distorted point.

Definition 4.1. Let $\Gamma$ be a finitely generated group endowed with a word distance $d=d_{\Gamma}$, and let $\Lambda$ be a subgroup of $\Gamma$.

- For $\varepsilon>0$ and $M>0$, we say that $g \in \Lambda$ is $(\varepsilon, M)$-quasi-convex if any geodesic $\gamma$ from e to $g$ spends at least a proportion $\varepsilon$ of its time in the $M$-neighborhood of $\Lambda$, i.e.,

$$
|\{i \in[1,|g|]: d(\gamma(i), \Lambda) \leqslant M\}| \geqslant \varepsilon|g| .
$$

We write $\Lambda_{Q C(\varepsilon, M)}$ for the set of points in $\Lambda$ which are $(\varepsilon, M)$-quasi-convex.

- Assume additionally that $\Lambda$ is finitely generated, and endowed with a word distance $d_{\Lambda}$. For $D>0$, we say that $g \in \Lambda$ is $D$-undistorted if $d_{\Lambda}(e, g) \leqslant D d_{\Gamma}(e, g)$. We write $\Lambda_{U D(D)}$ for the set of $D$-undistorted points.

Up to a change in the constants, these notions do not depend on the choice of the distance $d$. The first definition has the advantage to work for infinitely generated subgroups, but it may seem less natural than the second one. If $\Lambda$ is a quasi-convex subgroup of a hyperbolic group $\Gamma$, then all its points are $(1, M)$-quasi-convex if $M$ is large enough, and all its points are also $D$-undistorted for large enough $D$. In the general case, a quasi-convex point does not have to be undistorted: it may happen that the times $i$ such that $d(\gamma(i), \Lambda) \leqslant M$ are all included in $[1,|g| / 2]$, while between $|g| / 2$ and $|g|$ one needs to make a huge detour to follow $\Lambda$, making $d_{\Lambda}(e, g)$ much larger than $d_{\Gamma}(e, g)$. On the other hand, an undistorted point is automatically quasi-convex, at least in hyperbolic groups:

Proposition 4.2. Let $\Gamma$ be a hyperbolic group, let $\Lambda$ be a finitely generated subgroup of $\Gamma$, and let $D>0$. There exist $\varepsilon>0$ and $M>0$ such that any $D$-undistorted point is also $(\varepsilon, M)$-quasi-convex, i.e., $\Lambda_{U D(D)} \subset \Lambda_{Q C(\varepsilon, M)}$.

Proof. Consider $g \in \Lambda$ which is not $(\varepsilon, M)$-quasi-convex, we have to show that $d_{\Lambda}(e, g)$ is much bigger than $n=d_{\Gamma}(e, g)$. The intuition is that, away from a $\Gamma$-geodesic from $e$ to $g$, the progress towards $g$ is much slower by hyperbolicity.

Let us consider a geodesic from $e$ to $g$ in $\Lambda$, with length $d_{\Lambda}(e, g)$. Replacing each generator of $\Lambda$ by the product of a uniformly bounded number of generators of $\Gamma$, we obtain a path $\gamma_{\Lambda}$ in the Cayley graph of $\Gamma$, remaining in the $C_{0}$-neighborhood of $\Lambda$ (for some $C_{0}>0$ ) and with length $\left|\gamma_{\Lambda}\right| \leqslant C_{0} d_{\Lambda}(e, g)$.

Let us consider a geodesic $\gamma_{\Gamma}$ from $e$ to $g$ for the distance $d_{\Gamma}$. For each $x \in \Gamma$, we can consider its projection $\pi(x)$ on $\gamma_{\Gamma}$, i.e., the point on $\gamma_{\Gamma}$ that is closest to $x$ (if several points correspond, we take the closest one to $e$ ). This projection is 1-Lipschitz. In particular, the projection of $\gamma_{\Lambda}$ covers the whole geodesic $\gamma_{\Gamma}$. For each $x_{i} \in \gamma_{\Gamma}$, let us consider the first point $y_{i} \in \gamma_{\Lambda}$ projecting to $x_{i}$.

Let us fix an integer $L$, large enough with respect to the hyperbolicity constant of $\Gamma$. Along $\gamma_{\Gamma}$, let us consider the points at distance $k L$ from $e$, i.e., $x_{0}=e, x_{L}, x_{2 L}, \ldots, x_{m L}$ with $m=\lfloor n / L\rfloor$. In particular, $\left|\gamma_{\Lambda}\right| \geqslant \sum_{i} d_{\Gamma}\left(y_{i L}, y_{(i+1) L}\right)$. Moreover, a tree approximation shows that $d_{\Gamma}\left(y_{i L}, y_{(i+1) L}\right) \geqslant d_{\Gamma}\left(y_{i L}, x_{i L}\right)+L+d_{\Gamma}\left(x_{(i+1) L}, y_{(i+1) L}\right)-C_{1}$ (where $C_{1}$ only
depends on the hyperbolicity constant of $\Gamma$ ). Choosing $L \geqslant C_{1}$, we get

$$
\left|\gamma_{\Lambda}\right| \geqslant \sum_{i=0}^{m} d_{\Gamma}\left(x_{i L}, y_{i L}\right) \geqslant \sum_{i=0}^{m}\left(d_{\Gamma}\left(x_{i L}, \Lambda\right)-C_{0}\right)
$$

Since we assume that $g$ is not $(\varepsilon, M)$-quasi-convex, the set of indices $i$ with $d\left(x_{i}, \Lambda\right) \leqslant M$ has cardinality at most $\varepsilon n$. Taking $M \geqslant C_{0}$, the previous equation is bounded from below by

$$
(m+1-\varepsilon n) M-(m+1) C_{0} \geqslant(n / L-\varepsilon n) M-n C_{0} / L
$$

Finally, we get

$$
d_{\Lambda}(e, g) \geqslant\left|\gamma_{\Lambda}\right| / C_{0} \geqslant n(1 / L-\varepsilon) M / C_{0}-n / L
$$

If $\varepsilon$ is small enough and $M$ is large enough so that $(1 / L-\varepsilon) M / C_{0}-1 / L>D$, we obtain $d_{\Lambda}(e, g)>D n$, i.e., $g \notin \Lambda_{U D(D)}$, as desired.

From this point on, we will mainly work with the notion of quasi-convex points, since counting results on such points imply results on undistorted points by the previous proposition.
4.2. Non-distorted points in subgroups with $v_{\Lambda}=v$. In this section, we show that there are exponentially few quasi-convex points in infinite-index subgroups of hyperbolic groups.

Theorem 4.3. Let $\Gamma$ be a nonelementary hyperbolic group endowed with a word distance. Let $\Lambda$ be an infinite index subgroup of $\Gamma$. Then

$$
\begin{equation*}
\left|B_{n} \cap \Lambda\right|=o\left(\left|B_{n}\right|\right) \tag{4.1}
\end{equation*}
$$

Moreover, for all $\varepsilon>0$ and $M>0$, there exists $\eta>0$ such that, for all large enough $n$,

$$
\begin{equation*}
\left|B_{n} \cap \Lambda_{Q C(\varepsilon, M)}\right| \leqslant e^{-\eta n}\left|B_{n}\right| \tag{4.2}
\end{equation*}
$$

One may wonder why we put the estimate (4.1) in the statement of the theorem, while the main emphasis is on counting quasi-convex points. It turns out that this estimate is not trivial, and that its proof uses the same techniques as for the proof of (4.2). To illustrate that it is not trivial, let us remark that this estimate is not true without the hyperbolicity assumption. For instance, in $\Gamma=\mathbb{F}_{2} \times \mathbb{Z}$ (with its canonical generating system, and the corresponding word distance), the infinite index subgroup $\Lambda=\mathbb{F}_{2}$ satisfies $\left|\Lambda \cap B_{n}\right| /\left|B_{n}\right| \geqslant c>0$.

Theorem 4.3 is trivial if the growth rate $v_{\Lambda}$ of $\Lambda$ is strictly smaller than the growth rate $v$ of $\Gamma$, since in this case $\left|B_{n} \cap \Lambda\right|$ itself is exponentially smaller than $\left|B_{n}\right|$. However, this is not always the case, even for finitely generated subgroups.

Consider for instance a compact hyperbolic 3-manifold which fibers over the circle, obtained as a suspension of a hyperbolic surface with a pseudo-Anosov. Its fundamental group $\Gamma$ surjects into $\mathbb{Z}=\pi_{1}\left(\mathbb{S}^{1}\right)$. The kernel $\Lambda$ of this morphism $\varphi$ is the fundamental group of the fiber. It is finitely generated, with infinite index, and $\left|B_{n} \cap \Lambda\right| \sim c\left|B_{n}\right| / \sqrt{n}$, see [Sha98].

Heuristically, one can understand in this case why there are exponentially few quasiconvex points in $\Lambda$. Let us consider a geodesic of length $n$ in $\Gamma$. It projects under $\varphi$ to a path in $\mathbb{Z}$, which behaves roughly like a random walk. In particular, $e^{-n v}\left|\mathbb{S}^{n} \cap \Lambda\right|$ behaves like the probability that a random walk on $\mathbb{Z}$ comes back to the identity at time $n$. This
is of order $1 / \sqrt{n}$, in accordance with the rigorous results of [Sha98]. Such an element is quasi-convex if the random walk in $\mathbb{Z}$ spends a big proportion of its time close to the origin. A large deviation estimate shows that this is exponentially unlikely.

The proof of the theorem consists in making this heuristic precise, in the general case where the subgroup $\Lambda$ is not normal (so that there is no morphism $\varphi$ at hand). An important point in the proof is that a hyperbolic group is automatic, i.e., there exists a finite state automaton that recognizes a system of geodesics parameterizing bijectively the points in the group. Counting points in the group then amounts to a random walk on the graph of this automaton, while counting points in $\Lambda$ amounts to a fibred random walk, on this graph times $\Lambda \backslash \Gamma$. As this space is infinite, the random walk spends most of its time outside of finite sets, i.e., far away from $\Lambda$.

To formalize this argument, we will reduce the question to Markov chains on graphs, where we will use the following probabilistic lemma.

Lemma 4.4. Consider a Markov chain $\left(X_{n}\right)$ on a countable set $V$, with a stationary measure $m$ (i.e., $m(x)=\sum_{y} m(y) p(y, x)$ for all $x$ ). Let $\tilde{V}$ be the set of points $x \in V$ such that $\sum_{x \rightarrow y} m(y)=+\infty$, where we write $x \rightarrow y$ if there exists a positive probability path from $x$ to $y$. Then, for all $x \in V$ and $x^{\prime} \in \tilde{V}$,

$$
\begin{equation*}
\mathbb{P}_{x}\left(X_{n}=x^{\prime}\right) \rightarrow 0 \text { when } n \rightarrow \infty . \tag{4.3}
\end{equation*}
$$

Take $x \in \tilde{V}$ and $\varepsilon>0$. There exists $\eta>0$ such that, for all large enough $n$,

$$
\begin{equation*}
\mathbb{P}_{x}\left(X_{n}=x \text { and } X_{i} \text { visits } x \text { at least } \varepsilon n \text { times in between }\right) \leqslant e^{-\eta n} . \tag{4.4}
\end{equation*}
$$

Proof. In countable state Markov chains, a point $x$ can be either transient, or null recurrent, or positive recurrent. Let us first show that points in $\tilde{V}$ are not positive recurrent, by contradiction. Otherwise, the points that can be reached from $x$ form an irreducible class $\mathcal{C}$, which admits a stationary probability measure $p$. The restriction of $m$ to $\mathcal{C}$ is an excessive measure. By uniqueness (see [Rev84, Theorem 3.1.9]), the measure $m$ is proportional on $\mathcal{C}$ to $p$. In particular, it has finite mass there. This contradicts the assumption $\sum_{x \rightarrow y} m(y)=$ $+\infty$.

Let us now show that, for all $x \in V$ and $x^{\prime} \in \tilde{V}$, the probability $\mathbb{P}_{x}\left(X_{n}=x^{\prime}\right)$ tends to 0 . Otherwise, conditioning on the first visit to $x^{\prime}$, we deduce that $\mathbb{P}_{x^{\prime}}\left(X_{n}=x^{\prime}\right)$ does not tend to 0 . This implies that $x^{\prime}$ is positive recurrent, a contradiction.

Let us now prove (4.4). Consider $x \in \tilde{V}$, it is either transient or null recurrent. If it is transient, the probability $p$ to come back to $x$ is $<1$. Hence, the probability to come back $\varepsilon n$ times is bounded by $p^{\varepsilon n}$, and is therefore exponentially small as desired.

Assume now that $x$ is null recurrent: almost surely, the Markov chain comes back to $x$, but the waiting time $\tau$ has infinite expectation. Let $\tau_{1}, \tau_{2}, \ldots$ be the length of the successive excursions based at $x$. They are independent and distributed like $\tau$, by the Markov property. The probability in (4.4) is bounded by $\mathbb{P}\left(\sum_{i=1}^{\varepsilon n} \tau_{i} \leqslant n\right)$, which is bounded for any $M$ by $\mathbb{P}\left(\sum_{i=1}^{\varepsilon n} \tau_{i} 1_{\tau_{i} \leqslant M} \leqslant n\right)$. The random variables $\tau_{i} 1_{\tau_{i} \leqslant M}$ are bounded, independent and identically distributed. If $M$ is large enough, they have expectation $>1 / \varepsilon$. A standard large deviation result then shows that $\mathbb{P}\left(\sum_{i=1}^{\varepsilon n} \tau_{i} 1_{\tau_{i} \leqslant M} \leqslant n\right)$ is exponentially small, as desired.

We will also need the following technical lemma, which was explained to us by B. Bekka.

Lemma 4.5. Let $\Lambda$ be a subgroup of a group $\Gamma$. Assume that there exists a finite subset $B$ of $\Gamma$ such that $B \Lambda B=\Gamma$. Then $\Lambda$ has finite index in $\Gamma$.

Proof. We have by assumption $\Gamma=\bigcup_{i, j} b_{i} \Lambda b_{j}=\bigcup_{i, j} \Lambda_{i} b_{i} b_{j}$, where $\Lambda_{i}=b_{i} \Lambda b_{i}^{-1}$ is a conjugate of $\Lambda$ (and has therefore the same index). A theorem of Neumann [Neu54] ensures that a group is never a finite union of right cosets of infinite index subgroups. Hence, one of the $\Lambda_{i}$ has finite index in $\Gamma$, and so has $\Lambda$.

Let $\Gamma$ be a hyperbolic group, with a finite generating set $S$. Consider a finite directed graph $\mathcal{A}=\left(V, E, x_{*}\right)$ with vertex set $V$, edges $E$, a distinguished vertex $x_{*}$, and a labeling $\alpha: E \rightarrow S$. We associate to any path $\gamma$ in the graph (i.e., a sequence of edges $\sigma_{0}, \sigma_{1}, \ldots, \sigma_{m-1}$ where the endpoint of $\sigma_{i}$ is the beginning of $\sigma_{i+1}$ ) a path in the Cayley graph starting from the identity and following the edges labeled $\alpha\left(\sigma_{0}\right)$, then $\alpha\left(\sigma_{1}\right)$, and so on. The endpoint of this path is $\alpha_{*}(\gamma):=\alpha\left(\sigma_{0}\right) \cdots \alpha\left(\sigma_{m-1}\right)$. We always assume that any point can be reached by a path starting at $x_{*}$.

A hyperbolic group is automatic (see, for instance, [Cal13]): there exists such a graph with the following properties.
(1) For any path $\gamma$ in the graph, the corresponding path $\alpha(\gamma)$ is geodesic in the Cayley graph.
(2) The map $\alpha_{*}$ induces a bijection between the set of paths in the graph starting from $x_{*}$ and the group $\Gamma$.
In particular, the paths of length $n$ in the graph originating from $x_{*}$ parameterize the sphere $\mathbb{S}^{n}$ of radius $n$ in the group. The existence of such a structure makes it for instance possible to prove that the growth series of a hyperbolic group is rational. We will use such an automaton to count the points in the subgroup $\Lambda$, and in particular the quasi-convex points.

We define a transition matrix $A$, indexed by $V$. By definition, $A_{x y}$ is the number of edges from $x$ to $y$. Hence, $\left(A^{n}\right)_{x y}$ is the number of paths of length $n$ from $x$ to $y$. In particular, the number of paths of length $n$ starting from $x_{*}$ is $\sum_{y}\left(A^{n}\right)_{x * y}$. Write $u$ for the line vector with 1 at position $x_{*}$ and 0 elsewhere, and $\tilde{u}$ for the column vector with 1 everywhere. This number of paths reads $u A^{n} \tilde{u}$. Therefore, $\left|\mathbb{S}^{n}\right|=u A^{n} \tilde{u}$, proving the rationality of the growth function of the group. Let $v$ be the growth rate of balls in $\Gamma$. It satisfies $\left|B_{n}\right| \leqslant C e^{n v}$, by [Coo93]. Hence, the spectral radius of $A$ is $e^{v}$, and $A$ has no Jordan block for this maximal eigenvalue.

To understand the points of the infinite index subgroup $\Lambda$ of $\Gamma$, we consider an extension $\mathcal{A}_{\Lambda}$ of $\mathcal{A}$, with fibers $\Lambda \backslash \Gamma$. Its vertex set $V_{\Lambda}$ is made of the pairs $(x, \Lambda g) \in V \times \Lambda \backslash \Gamma$. For any edge $\sigma$ in $\mathcal{A}$, going from $x$ to $y$ and with label $\alpha(\sigma)$, we put for any $g \in \Gamma$ an edge in $\mathcal{A}_{\Lambda}$ from $(x, \Lambda g)$ to $(y, \Lambda g \alpha(\sigma))$. A path $\gamma$ in $\mathcal{A}$, from $x$ to $y$, lifts to a path $\tilde{\gamma}$ in $\mathcal{A}_{\Lambda}$ originating from $(x, \Lambda e)$. By construction, its endpoint is $\left(y, \Lambda \alpha_{*}(\gamma)\right)$. This shows that the paths in the graph $\mathcal{A}_{\Lambda}$ remember the current right coset of $\Lambda$.

The next lemma proves that the relevant components of this fibred graph are infinite.
Lemma 4.6. Let $\tilde{x}_{0}=\left(x_{0}, \Lambda g_{0}\right)$ belong to $\mathcal{A}_{\Lambda}$. Let $\mathcal{C}$ be the component of $x_{0}$ in $\mathcal{A}$ (i.e., the set of points that can be reached from $x_{0}$ and from which one can go back to $x_{0}$ ). Let $A_{\mathcal{C}}$ be the restriction of the matrix $A$ to the points in $\mathcal{C}$. Assume that its spectral radius $\rho\left(A_{\mathcal{C}}\right)$ is
equal to $e^{v}$. Then, starting from $\tilde{x}_{0}$ in the graph $\mathcal{C}_{\Lambda}$ (the restriction of $\mathcal{A}_{\Lambda}$ to $\mathcal{C} \times \Lambda \backslash \Gamma$ ), one can reach infinitely many different points of $\mathcal{C}_{\Lambda}$.

Proof. It suffices to show that one can reach infinitely many points whose component in $\mathcal{C}$ is $x_{0}$. Assume by contradiction that one can only reach a finite number of classes $\left(x_{0}, \Lambda g_{i}\right)$.

Given $w \in \Gamma$ and $C>0$, let $Y_{w, C}$ be the set of points in $\Gamma$ that have a geodesic expression in which, for any subword $\tilde{w}$ of this expression and for any $a, b$ with length at most $C$, one has $w \neq a \tilde{w} b$. In other words, the points in $Y_{w, C}$ are those that never see $w$ (nor even a thickening of $w$ of size $C$ ) in their geodesic expressions. Theorem 3 in [AL02] proves the existence of $C_{0}$ such that, for any $w$, the quantity $\left|B_{n} \cap Y_{w, C_{0}}\right| /\left|B_{n}\right|$ tends to 0 (the important point is that $C_{0}$ does not depend on $w$ ).

The number of paths in $\mathcal{C}$ originating from $x_{0}$ grows at least like $c\left|B_{n}\right|$ since the spectral radius of $A_{\mathcal{C}}$ is $e^{v}$. These paths give rise to distinct points in $\Gamma$. Hence, there exists such a path $\gamma_{0}$ such that $\alpha_{*}\left(\gamma_{0}\right) \notin Y_{w, C_{0}}$. In particular, there exists a subpath $\gamma_{1}$ such that $\alpha_{*}\left(\gamma_{1}\right)$ can be written as $a_{1} w b_{1}$ with $\left|a_{1}\right| \leqslant C_{0}$ and $\left|b_{1}\right| \leqslant C_{0}$. We can choose a path from $x_{0}$ to the starting point of $\gamma_{1}$, with fixed length (since $\mathcal{C}$ is finite), and another path from the endpoint of $\gamma_{1}$ to $x_{0}$. Concatenating them, we get a path $\gamma_{2}$ from $x_{0}$ to itself with $\alpha_{*}\left(\gamma_{2}\right)=a_{2} w b_{2}$ and $\left|a_{2}\right|,\left|b_{2}\right| \leqslant C_{1}=C_{0}+2 \operatorname{diam}(\mathcal{C})$. By assumption, $\Lambda g_{0} \alpha_{*}\left(\gamma_{2}\right)$ is one of the finitely many $\Lambda g_{i}$ since we are returning to $x_{0}$. Hence, there exists $\lambda \in \Lambda$ such that $g_{0} a_{2} w b_{2}=\lambda g_{i}$. This shows that $w \in B \Lambda B$, where $B$ is the ball of radius $C_{1}+\max _{i} d\left(e, g_{i}\right)$. As this holds for any $w$, we have proved that $B \Lambda B=\Gamma$. By Lemma 4.5, this shows that $\Lambda$ has finite index in $\Gamma$, a contradiction.

Lemma 4.7. Let $K\left(n, \tilde{x}_{0}, \varepsilon_{0}\right)$ denote the set of paths in $\mathcal{A}_{\Lambda}$ starting at a point $\tilde{x}_{0}$, of length $n$, coming back to $\tilde{x}_{0}$ at time $n$, and spending a proportion at least $\varepsilon_{0}$ of the time at $\tilde{x}_{0}$. Consider $\tilde{x}_{0} \in \mathcal{A}_{\Lambda}$ and $\varepsilon_{0}>0$. Then there exist $\eta>0$ and $C>0$ such that, for all $n \in \mathbb{N}$,

$$
\left|K\left(n, \tilde{x}_{0}, \varepsilon_{0}\right)\right| \leqslant C e^{n(v-\eta)}
$$

Proof. Write $\tilde{x}_{0}=\left(x_{0}, \Lambda g_{0}\right)$, let $\mathcal{C}$ be the component of $x_{0}$ in $\mathcal{A}$. If the spectral radius of the restricted transition matrix $A_{\mathcal{C}}$ is $<e^{v}$, we simply bound $\left|K\left(n, \tilde{x}_{0}, \varepsilon_{0}\right)\right|$ by the number of paths in $\mathcal{C}$ from $x_{0}$ to itself. This is at most $\left\|A_{\mathcal{C}}^{n}\right\|$, which is exponentially smaller than $e^{n v}$ as desired.

Assume now that $\rho\left(A_{\mathcal{C}}\right)=e^{v}$. We will understand the number of paths in $\mathcal{C}$ (and in its lift $\mathcal{C}_{\Lambda}$ ) in terms of a Markov chain. The matrix $A_{\mathcal{C}}$ has a unique eigenvector $q$ corresponding to the eigenvalue $e^{v}$, it is positive by Perron-Frobenius's theorem. By definition, $p(x, y)=$ $e^{-v} A_{x y} q(y) / q(x)$ satisfies, for any $x \in \mathcal{C}$,

$$
\sum_{y \in \mathcal{C}} p(x, y)=\frac{e^{-v}}{q(x)} \sum A_{x y} q(y)=1
$$

This means that $p(x, y)$ is a transition kernel on $\mathcal{C}$. Denote by $\left(X_{n}\right)_{n \in \mathbb{N}}$ the corresponding Markov chain. By construction,

$$
\mathbb{P}_{x}\left(X_{n}=y\right)=e^{-n v}\left(A^{n}\right)_{x y} q(y) / q(x)
$$

Moreover, $\left(A^{n}\right)_{x y}$ is the number of paths of length $n$ in $\mathcal{A}$ from $x$ to $y$. Hence, up to a bounded multiplicative factor $q(y) / q(x)$, the transition probabilities of the Markov chain $X_{n}$
count the number of paths in the graph $\mathcal{C}$. Let $m$ denote the unique stationary probability for the Markov chain on $\mathcal{C}$.

We lift everything to $\mathcal{C}_{\Lambda}$, assigning to an edge the transition probability of its projection in $\mathcal{C}$. The stationary measure $m$ lifts to a stationary measure $m_{\Lambda}$, which is simply the product of $m$ and of the counting measure in the direction $\Lambda \backslash \Gamma$. Denoting by $X_{n}^{\Lambda}$ the Markov chain in $\mathcal{C}_{\Lambda}$, we have

$$
e^{-n v}\left|K\left(n, \tilde{x}_{0}, \varepsilon_{0}\right)\right|=\mathbb{P}_{\tilde{x}_{0}}\left(X_{n}^{\Lambda}=\tilde{x}_{0} \text { and } X_{i}^{\Lambda} \text { visits } \tilde{x}_{0} \text { at least } \varepsilon_{0} n \text { times in between }\right) .
$$

By Lemma 4.6 , the Markov chain starting from $\tilde{x}_{0}$ can reach infinitely many points. Equivalently, since $m$ is bounded from below, it can reach a set of infinite $m_{\Lambda}$-measure. Therefore, Lemma 4.4 applies, and shows that the above quantity is exponentially small.

Proof of Theorem 4.3. Let us first prove (4.2). Counting the points in $\mathbb{S}^{n} \cap \Lambda_{Q C(\varepsilon, M)}$ amounts to counting the paths of length $n$ in $\mathcal{A}_{\Lambda}$, starting from $\left(x_{*}, \Lambda e\right)$ and spending a proportion at least $\varepsilon$ of their time in the finite subset $F=V \times \Lambda B_{M} \subset V_{\Lambda}$. Such a path spends a proportion at least $\varepsilon_{0}=\varepsilon /|F|$ of its time at a given point $\tilde{x} \in F$. Let $k$ and $k+m$ denote the first and last visits to $\tilde{x}$ (with $m \geqslant \varepsilon_{0} n$ since there are at least $\varepsilon_{0} n$ visits). Such a path is the concatenation of a path from $\left(x_{*}, \Lambda e\right)$ to $\tilde{x}$ of length $k$ (their number is bounded by the corresponding number of paths in $\mathcal{A}$, at most $\left\|A^{k}\right\| \leqslant C e^{k v}$, of a path in $K\left(m, \tilde{x}, \varepsilon_{0}\right)$, and of a path starting from $\tilde{x}$ of length $n-k-m$ (their number is again bounded by the number of corresponding paths in $\mathcal{A}$, at most $\left.C e^{(n-k-m) v}\right)$. Hence, their number is at most $C e^{(n-m) v}\left|K\left(m, \tilde{x}, \varepsilon_{0}\right)\right|$. Summing over the points $\tilde{x} \in F$, over the at most $n$ possible values of $k$, and the values of $m$, we get the inequality

$$
\left|\mathbb{S}^{n} \cap \Lambda_{Q C(\varepsilon, M)}\right| \leqslant C n e^{n v} \sum_{\tilde{x} \in F} \sum_{m=\varepsilon_{0} n}^{n} e^{-m v}\left|K\left(m, \tilde{x}, \varepsilon_{0}\right)\right|
$$

Lemma 4.7 shows that this is exponentially smaller than $e^{n v}$.
Let us now prove (4.1), using similar arguments. A point in $\mathbb{S}^{n} \cap \Lambda$ corresponds to a path of length $n$ in $\mathcal{A}_{\Lambda}$, starting from $\left(x_{*}, \Lambda e\right)$ and ending at a point $(x, \Lambda e)$. We say that a component $\mathcal{C}$ in the graph $\mathcal{A}$ is maximal if the spectral radius of the corresponding restricted matrix $A_{\mathcal{C}}$ is $e^{v}$. Since the matrix $A$ has no Jordan block corresponding to the eigenvalue $e^{v}$, a path in the graph encounters at most one maximal component. The paths in $\mathcal{A}_{\Lambda}$ whose projection in $\mathcal{A}$ spends a time $k$ in non-maximal components give an overall contribution to $\left|\mathbb{S}^{n} \cap \Lambda\right|$ bounded by $C e^{(n-k) v+k(v-\eta)} \leqslant C e^{-\eta k}\left|B_{n}\right|$. Given $\varepsilon>0$, their contribution for $k \geqslant k_{0}(\varepsilon)$ is bounded by $\varepsilon\left|B_{n}\right|$. Hence, it suffices to control the paths for fixed $k$. Let us fix the beginning of such a path, from $\left(x_{*}, \Lambda e\right)$ to a point $\left(x_{0}, \Lambda g_{0}\right)$ where $x_{0}$ is in a maximal component $\mathcal{C}$, and its end from $\left(x_{1}, \Lambda g_{1}\right)$ with $x_{1} \in \mathcal{C}$ to a point $(x, \Lambda e)$. To conclude, one should show that the number of paths of length $n$ from $\left(x_{0}, \Lambda g_{0}\right)$ to $\left(x_{1}, \Lambda g_{1}\right)$ is $o\left(e^{n v}\right)$. This follows from the probabilistic interpretation in the proof of Lemma 4.7 and from (4.3).
4.3. Non-distorted points in subgroups with $v_{\Lambda}<v$. Let $\Lambda$ be a subgroup of a hyperbolic group $\Gamma$. Let $v_{\Lambda}$ and $v_{\Gamma}$ be their respective growths, for a word distance on $\Gamma$. If $v_{\Lambda}=v_{\Gamma}$, Theorem 4.3 proves that there is a dichotomy:
(1) Either $\Lambda$ is quasi-convex (equivalently, $\Lambda$ has finite index in $\Gamma$ ). Then $\left|B_{n} \cap \Lambda\right| \geqslant$ $c e^{n v_{\Lambda}}$, and all points in $\Lambda$ are quasi-convex.
(2) Or $\Lambda$ is not quasi-convex (equivalently, it has infinite index in $\Gamma$ ). Then $\left|B_{n} \cap \Lambda\right|=$ $o\left(e^{n v_{\Lambda}}\right)$, and there are exponentially few quasi-convex points in $\Lambda$.
Consider now a general subgroup $\Lambda$ with $v_{\Lambda}<v_{\Gamma}$. If it is quasi-convex, then (1) above is still satisfied: $\left|B_{n} \cap \Lambda\right| \geqslant c e^{n v_{\Lambda}}$ by [Coo93], and all points in $\Lambda$ are quasi-convex. One may ask if these properties are equivalent, and if they characterize quasi-convex subgroups. This question is reminiscent of a question of Sullivan in hyperbolic geometry: Are convex cocompact groups the only ones to have finite Patterson-Sullivan measure? Peigné showed in [Pei03] that the answer to this question is negative. His counterexamples adapt to our situation, giving also a negative answer to our question.

Proposition 4.8. There exists a finitely generated subgroup $\Lambda$ of a hyperbolic group $\Gamma$ endowed with a word distance, which is not quasi-convex, but for which $C^{-1} e^{n v_{\Lambda}} \leqslant\left|B_{n} \cap \Lambda\right| \leqslant$ $C e^{n v_{\Lambda}}$. Moreover, most points of $\Lambda$ are quasi-convex: there exist $\varepsilon$ and $\eta$ such that

$$
\begin{equation*}
\left|B_{n} \cap \Lambda \backslash \Lambda_{Q C(\varepsilon, 0)}\right| \leqslant C e^{n\left(v_{\Lambda}-\eta\right)} \tag{4.5}
\end{equation*}
$$

Proof. The example is the same as in [Pei03], but his geometric proofs are replaced by combinatorial arguments based on generating series.

Let $G$ be a finitely generated non-quasi-convex subgroup of a hyperbolic group $\tilde{G}$ (take for instance for $\tilde{G}$ the fundamental group of a hyperbolic 3-manifold which fibers over the circle, and for $G$ the fundamental group of the fiber of this fibration). Let $H=\mathbb{F}_{k}$, with $k$ large enough so that $v_{H} \geqslant v_{G}$. We take $\Lambda=G * H \subset \Gamma=\tilde{G} * H$. It is not quasi-convex, because of the factor $G$. Writing $v_{\Lambda}$ for its growth, we claim that, for some $c>0$,

$$
\begin{equation*}
\left|\mathbb{S}^{n} \cap \Lambda\right| \sim c e^{n v_{\Lambda}} \tag{4.6}
\end{equation*}
$$

We compute with generating series. Let $F_{G}(z)$ be the growth series for $G$, given by $F_{G}(z)=$ $\sum_{n \geqslant 0}\left|\mathbb{S}^{n} \cap G\right| z^{n}$. Likewise, we define $F_{H}$ and $F_{\Lambda}$. Since any word in $\Lambda$ has a canonical decomposition in terms of words in $G$ and $H$, a classical computation (see [dlH00, Prop. VI.A.4]) gives

$$
\begin{equation*}
F_{\Lambda}=\frac{F_{G} F_{H}}{1-\left(F_{G}-1\right)\left(F_{H}-1\right)} \tag{4.7}
\end{equation*}
$$

Let $z_{G}=e^{-v_{G}} \geqslant z_{H}=e^{-v_{H}}$ be the convergence radii of $F_{G}$ and $F_{H}$. At $z_{H}$, we have $F_{H}\left(z_{H}\right)=+\infty$, since the cardinality of spheres in the free group is exactly of the order of $e^{n v_{H}}$. When $z$ increases to $z_{H}$, the function $\left(F_{G}(z)-1\right)\left(F_{H}(z)-1\right)$ takes the value 1 , at a number $z=z_{\Lambda}$. Since this is the first singularity of $F_{\Lambda}$, we have $z_{\Lambda}=e^{-v_{\Lambda}}$. Moreover, the function $F_{\Lambda}$ is meromorphic at $z_{\Lambda}$, with a pole of order 1 (since the function $\left(F_{G}-1\right)\left(F_{H}-1\right)$ has positive derivative, being a power series with nonnegative coefficients). It follows from a simple tauberian theorem (see, for instance, [FS09, Theorem IV.10]) that the coefficients of $F_{\Lambda}$ behave like $c z_{\Lambda}^{-n}$, proving (4.6).

Let us estimate the number of non-quasi-convex points in $\Lambda$. Consider a word $w \in \Lambda$ of length $n$, for instance starting with a factor in $G$ and ending with a factor in $H$. It can be written as $g_{1} h_{1} g_{2} h_{2} \cdots h_{s}$. Along a geodesic from $e$ to $w$, all the words $g_{1} h$ (with $h$ prefix of $h_{1}$ ) belong to $\Lambda$. So do all the words $g_{1} h_{1} g_{2} h$ with $h$ prefix of $h_{2}$, and so on. Therefore, the proportion of time that the geodesic spends outside of $\Lambda$ is at most $\sum\left|g_{i}\right| / n$. Such a point in $\Lambda \backslash \Lambda_{Q C(\varepsilon, 0)}$ satisfies $\sum\left|g_{i}\right| \geqslant(1-\varepsilon) n$ and $\sum\left|h_{i}\right| \leqslant \varepsilon n$. Assuming $\varepsilon \leqslant 1 / 2$, this
gives $\sum\left|h_{i}\right| \leqslant(\varepsilon / 2) \sum\left|g_{i}\right|$. In particular, for any $\alpha>0$, we have $e^{\alpha\left(\sum\left|g_{i}\right|-2 \varepsilon^{-1} \sum\left|h_{i}\right|\right)} \geqslant 1$. Let $u_{n}=\left|\mathbb{S}^{n} \cap \Lambda \backslash \Lambda_{Q C(\varepsilon, 0)}\right|$, its generating series satisfies the following equation (where we only write in details the words starting with $G$ and ending in $H$, the other ones being completely analogous):

$$
\begin{aligned}
\sum u_{n} z^{n} & \leqslant \sum_{\ell \geqslant 1} \sum_{a_{1}+b_{1}+a_{2}+\cdots+b_{\ell}=n} e^{\alpha\left(\sum a_{i}-2 \varepsilon^{-1} \sum b_{i}\right)}\left|\mathbb{S}^{a_{1}} \cap G\right|\left|\mathbb{S}^{b_{1}} \cap H\right| \cdots\left|\mathbb{S}^{b_{\ell}} \cap H\right| z^{n}+\ldots \\
& =\sum_{\ell \geqslant 1}\left[\left(F_{G}\left(e^{\alpha} z\right)-1\right)\left(F_{H}\left(e^{-2 \alpha \varepsilon^{-1}} z\right)-1\right)\right]^{\ell}+\ldots \\
& =\frac{F_{G}\left(e^{\alpha} z\right) F_{H}\left(e^{-2 \alpha \varepsilon^{-1}} z\right)}{1-\left(F_{G}\left(e^{\alpha} z\right)-1\right)\left(F_{H}\left(e^{-2 \alpha \varepsilon^{-1}} z\right)-1\right)}
\end{aligned}
$$

This is the same formula as in (4.7), but the factor $z$ has been shifted in $F_{G}$ and $F_{H}$. Choose $\alpha>0$ such that $e^{\alpha} z_{\Lambda}<z_{G}$, and then $\varepsilon$ small enough so that $\left(F_{G}\left(e^{\alpha} z_{\Lambda}\right)-1\right)\left(F_{H}\left(e^{-2 \alpha \varepsilon^{-1}} z_{\Lambda}\right)-\right.$ $1)<1$. We deduce that the series $\sum u_{n} z^{n}$ converges for $z=z_{\Lambda}$, and even slightly to its right. It follows that $u_{n}$ is exponentially small compared to $z_{\Lambda}^{-n}$. This proves (4.5).
4.4. Application to random walks in infinite index subgroups. In this paragraph, we use Theorem 4.3 to prove Theorem 1.6 on random walks given by a measure $\mu$ on a hyperbolic group $\Gamma$, assuming that $\Gamma_{\mu}$ has infinite index in $\Gamma$.

Before proving Theorem 1.6, we give another easier result, pertaining to the case where $\mu$ has a finite moment for a word distance on $\Gamma_{\mu}$ (which should be finitely generated): In this case, the random walk typically visits undistorted points. This easy statement is not used later on, but it gives a heuristic explanation to Theorem 1.6.

Lemma 4.9. Let $\Lambda$ be a finitely generated subgroup of a finitely generated group $\Gamma$. Let $d_{\Lambda}$ and $d_{\Gamma}$ be the two corresponding word distances. Consider a probability measure $\mu$ on $\Lambda$, with a moment of order 1 for $d_{\Lambda}$ (and therefore for $d_{\Gamma}$ ), with nonzero drift for $d_{\Gamma}$. Let $X_{n}$ denote the corresponding random walk. There exists $D>0$ such that $\mathbb{P}\left(X_{n} \in \Lambda_{U D(D)}\right) \rightarrow 1$.

Proof. Almost surely, $d_{\Gamma}\left(e, X_{n}\right) \sim \ell_{\Gamma} n$, for some nonzero drift $\ell_{\Gamma}$. In the same way, $d_{\Lambda}\left(e, X_{n}\right) \sim \ell_{\Lambda} n$. For any $D>\ell_{\Lambda} / \ell_{\Gamma}$, we get almost surely $d_{\Lambda}\left(e, X_{n}\right) \leqslant D d_{\Gamma}\left(e, X_{n}\right)$ for large enough $n$, i.e., $X_{n} \in \Lambda_{U D(D)}$.

This lemma readily implies Theorem 1.6 under the additional assumption that $\Lambda$ is finitely generated and that $\mu$ has a moment of order 1 for $d_{\Lambda}$. Indeed, for large $n$, with probability at least $1 / 2$, the point $X_{n}$ belongs to $B_{(\ell+\varepsilon) n} \cap \Lambda_{U D(D)}$, whose cardinality is bounded by $C e^{(\ell+\varepsilon) n(v-\eta)}$ according to Theorem 4.3. Lemma 2.4 yields $h \leqslant(\ell+\varepsilon)(v-\eta)$, hence $h \leqslant \ell(v-\eta)<\ell v$, completing the proof.

However, the assumptions of Theorem 1.6 are much weaker: even when $\Lambda$ is finitely generated, it is much more restrictive to require a moment of order 1 on $\Lambda$ than on $\Gamma$, precisely because the $\Gamma$-distance is smaller than the $\Lambda$-distance on distorted points, which make up most of $\Lambda$. The general proof will not use undistorted points (which are not even defined when $\Lambda$ is not finitely generated), but rather quasi-convex points: we will show that, typically, the random walk concentrates on quasi-convex points. With the previous argument, Theorem 1.6 readily follows from the next lemma.

Lemma 4.10. Let $\Lambda$ be a subgroup of a hyperbolic group $\Gamma$ endowed with a word distance $d=d_{\Gamma}$. Let us consider a probability measure $\mu$ on $\Lambda$, with a moment of order 1 for $d_{\Gamma}$. There exist $\varepsilon>0$ and $M>0$ such that $\mathbb{P}\left(X_{n} \in \Lambda_{Q C(\varepsilon, M)}\right) \geqslant 1 / 2$ for large enough $n$.

Proof. The lemma is trivial if $\mu$ is elementary, since all the elements of $\Gamma_{\mu} \subset \Lambda$ are then quasi-convex. We may therefore assume that $\mu$ is non-elementary.

The random walk at time $n$ is given by $X_{n}=g_{1} \cdots g_{n}$, where $g_{i}$ are independent and distributed like $\mu$. We will show that most products $g_{1} \cdots g_{i}$ (which belong to $\Lambda$ ) are within distance $M$ of a geodesic from $e$ to $X_{n}$ (this amounts to the classical fact that trajectories of the random walk follow geodesics in the group), and moreover that they approximate a proportion at least $\varepsilon$ of the points on this geodesic. This will give $X_{n} \in \Lambda_{Q C(\varepsilon, M)}$ as desired. The second point is more delicate: we should for instance exclude the situation where, given a geodesic $\gamma$, one has $X_{n}=\gamma(a(n))$ where $a(n)$ is the smallest square larger than $n$. In this case, $X_{n}$ follows the geodesic $\gamma$ at linear speed, but nevertheless the proportion of $\gamma$ it visits tends to 0 . This behavior will be excluded thanks to the fact that, with high probability, the jumps of the random walk are bounded.

The argument is probabilistic and formulated in terms of the bilateral version of the random walk. On $\Omega=\Gamma^{\mathbb{Z}}$ with the product measure $\mathbb{P}=\mu^{\otimes \mathbb{Z}}$, let $g_{n}$ be the $n$-th coordinate. The $g_{n}$ are independent, identically distributed, and correspond to the increments of a random walk $\left(X_{n}\right)_{n \in \mathbb{Z}}$ with $X_{0}=e$ and $X_{n}^{-1} X_{n+1}=g_{n+1}$. Almost surely, $X_{n}$ converges when $n \rightarrow \pm \infty$ towards two random variables $\xi^{ \pm} \in \partial \Gamma$, with $\xi^{+} \neq \xi^{-}$almost surely since these random variables are independent and atomless. Following Kaimanovich [Kai00], denote by $S\left(\xi^{-}, \xi^{+}\right)$the union of all the geodesics from $\xi^{-}$to $\xi^{+}$. Let $\pi$ be the projection on $S\left(\xi^{-}, \xi^{+}\right)$, i.e., $\pi(g)$ is the closest point to $g$ on $S\left(\xi^{-}, \xi^{+}\right)$. It is not uniquely defined, but two possible choices are within distance $C_{0}$, for some $C_{0}$ only depending on $\Gamma$.

Let us choose $L>0$ large enough (how large will only depend on the hyperbolicity constant of the space). Any measurable function is bounded on sets with arbitrarily large measure. Hence, there exists $K>0$ such that, with probability at least $9 / 10$,
(1) For every $|k| \geqslant K$, the projections $\pi\left(X_{k}\right)$ are distant from $\pi\left(X_{0}\right)$ by at least $L$ (and they are closer to $\xi^{+}$if $k>0$, and to $\xi^{-}$if $k<0$ ).
(2) We have $d\left(e, S\left(\xi^{-}, \xi^{+}\right)\right) \leqslant K$.

As everything is equivariant, we deduce that, for all $i \in \mathbb{Z}$, the point $X_{i}$ satisfies the same properties with probability at least $9 / 10$, i.e.,

$$
\begin{equation*}
d\left(X_{i}, S\left(\xi^{-}, \xi^{+}\right)\right) \leqslant K \text { and, for all }|k| \geqslant K, d\left(\pi\left(X_{i}\right), \pi\left(X_{i+k}\right)\right) \geqslant L . \tag{4.8}
\end{equation*}
$$

Let $n$ be a large integer. Write $m=\lfloor n / K\rfloor$. Among the integers $K, 2 K, \ldots, m K \leqslant n$, we consider the set $I_{n}(\omega)$ of those $i$ such that $X_{i}$ satisfies (4.8). We have $\mathbb{E}\left(\left|I_{n}\right|\right) \geqslant m \cdot 9 / 10$. As $\left|I_{n}\right| \leqslant m$, we get

$$
\frac{9 m}{10} \leqslant \mathbb{E}\left(\left|I_{n}\right|\right) \leqslant \frac{m}{10} \mathbb{P}\left(\left|I_{n}\right|<m / 10\right)+m \mathbb{P}\left(\left|I_{n}\right| \geqslant m / 10\right)=\frac{m}{10}+\frac{9 m}{10} \mathbb{P}\left(\left|I_{n}\right| \geqslant m / 10\right)
$$

This gives $\mathbb{P}\left(\left|I_{n}\right| \geqslant m / 10\right) \geqslant 8 / 9$. Let $\eta=1 /(20 K)$. Let $\Omega_{n}$ be the set of $\omega$ such that $\left|I_{n}(\omega)\right| \geqslant \eta n+1$, and $X_{0}$ and $X_{n}$ satisfy (4.8), and $d\left(X_{n}, e\right) \leqslant 2 \ell n$ (where $\ell$ is the drift of $\mu$ ). It satisfies $\mathbb{P}\left(\Omega_{n}\right) \geqslant 1 / 2$ if $n$ is large enough. This is the set of good trajectories for which we can control the position of many of the $X_{i}$.


Figure 1. The projections on $\gamma$ and $S$

Let $\omega \in \Omega_{n}$. We write $Y_{i}$ for a projection of $X_{i}$ on a geodesic $\gamma$ from $e$ to $X_{n}$. Let $\tilde{I}_{n}=I_{n} \backslash\{m K\}$, so that the elements of $\tilde{I}_{n}$ are at distance at least $K$ of 0 and $n$. As $X_{0}$ and $X_{n}$ satisfy (4.8), the projections $\pi\left(X_{i}\right)$ for $i \in \tilde{I}_{n}$ are located between $\pi\left(X_{0}\right)$ and $\pi\left(X_{n}\right)$, and are at a distance at least $L$ of these points (see Figure 1). If $L$ is large enough, we obtain $d\left(\pi\left(X_{i}\right), Y_{i}\right) \leqslant C_{1}$ by hyperbolicity, where $C_{1}$ only depends on $\Gamma$. This gives

$$
d\left(Y_{i}, \Lambda\right) \leqslant d\left(Y_{i}, \pi\left(X_{i}\right)\right)+d\left(\pi\left(X_{i}\right), X_{i}\right) \leqslant C_{1}+K
$$

thanks to (4.8) for $X_{i}$. When $i \neq j$ belong to $\tilde{I}_{n}$, we have $d\left(\pi\left(X_{i}\right), \pi\left(X_{j}\right)\right) \geqslant L$ again thanks to (4.8), hence $d\left(Y_{i}, Y_{j}\right) \geqslant L-2 C_{1}$. If $L$ was chosen larger than $2 C_{1}+1$, this shows that $Y_{i} \neq Y_{j}$. We have found along $\gamma$ at least $\left|I_{n}\right|-1$ distinct points, within distance $C_{1}+K$ of $\Lambda$. Moreover, for large enough $n$,

$$
\left|I_{n}\right|-1 \geqslant \eta n \geqslant 2 \ell n \cdot(\eta / 2 \ell) \geqslant d\left(e, X_{n}\right) \cdot(\eta / 2 \ell)
$$

Let $\varepsilon=\eta / 2 \ell$ and $M=C_{1}+K$. We have shown that, for $\omega \in \Omega_{n}$ (whose probability is at least $1 / 2)$, the point $X_{n}(\omega)$ belongs to $\Lambda_{Q C(\varepsilon, M)}$.

## 5. Construction of maximizing measures

In this section, we prove Theorem 1.7: Given any finite subset $\Sigma$ in a hyperbolic group $\Gamma$, there exists a measure $\mu_{\Sigma}$ maximizing the quantity $h(\mu) / \ell(\mu)$ over all measures $\mu$ supported on $\Sigma$ with $\ell(\mu)>0$. To prove this result, we start with a sequence of measures $\mu_{i}$ supported on $\Sigma$ such that $h\left(\mu_{i}\right) / \ell\left(\mu_{i}\right)$ converges to the maximum $M$ of these quantities. We are looking for $\mu_{\Sigma}$ with $h\left(\mu_{\Sigma}\right) / \ell\left(\mu_{\Sigma}\right)=M$. Replacing $\mu_{i}$ with $\left(\mu_{i}+\delta_{e}\right) / 2$ (this multiplies entropy and drift by $1 / 2$, and does not change their ratio) and adding $e$ to $\Sigma$, we can always assume $\mu_{i}(e) \geqslant 1 / 2$, to avoid periodicity problems.

Extracting a subsequence, we can ensure that $\mu_{i}$ converges to a limit probability measure $\mu$. We treat separately the two following cases:
(1) $\Gamma_{\mu}$ is non-elementary.
(2) $\Gamma_{\mu}$ is elementary.

Let us handle first the easy case, where $\Gamma_{\mu}$ is non-elementary. In this case, the entropy and the drift are continuous at $\mu$, by Proposition 2.3 and Theorem 2.9, both due to Erschler
and Kaimanovich in [EK13]. Therefore, $h\left(\mu_{i}\right) / \ell\left(\mu_{i}\right)$ tend to $h(\mu) / \ell(\mu)$, since in this case $\ell(\mu)>0$. One can thus take $\mu_{\Sigma}=\mu$.

The case where $\Gamma_{\mu}$ is elementary is much more interesting. Let us describe heuristically what should happen, in a simple case. We assume that $\mu_{i}=(1-\varepsilon) \mu+\varepsilon \nu$ where $\nu$ is a fixed measure, and $\varepsilon$ tends to 0 . The random walk for $\mu_{i}$ can be described as follows. At each jump, one picks $\mu$ (with probability $1-\varepsilon$ ) or $\nu$ (with probability $\varepsilon$ ), then one jumps according to the chosen measure. After time $N$, the measure $\nu$ is chosen roughly $\varepsilon N$ times, with intervals of length $1 / \varepsilon$ in between, where $\mu$ is chosen. Thus, $\mu_{i}^{* N}$ behaves roughly like $\left(\mu^{* 1 / \varepsilon} * \nu\right)^{\varepsilon N}$.

When $\Gamma_{\mu}$ is finite, the measure $\mu^{* 1 / \varepsilon}$ is close, when $\varepsilon$ is small, to the uniform measure $\pi$ on $\Gamma_{\mu}$. Therefore, $\mu_{i}^{* N}$ is close to $(\pi * \nu)^{\varepsilon N}$. We deduce $h\left(\mu_{i}\right) \sim \varepsilon h(\pi * \nu)$ and $\ell\left(\mu_{i}\right) \sim \varepsilon \ell(\pi * \nu)$. In particular, $h\left(\mu_{i}\right) / \ell\left(\mu_{i}\right) \rightarrow h(\pi * \nu) / \ell(\pi * \nu)$. One can take $\mu_{\Sigma}=\pi * \nu$.

When $\Gamma_{\mu}$ is infinite, it is virtually cyclic. Assuming that $\mu$ is centered for simplicity, the walk given by $\mu^{* 1 / \varepsilon}$ arrives essentially at distance $1 / \sqrt{\varepsilon}$ of the origin, by the central limit theorem. Then, one jumps according to $\nu$, in a direction transverse to $\Gamma_{\mu}$, preventing further cancellations. Hence, the walk given by $\left(\mu^{* 1 / \varepsilon} * \nu\right)^{\varepsilon N}$ is at distance roughly $\varepsilon N / \sqrt{\varepsilon}$ from the origin, yielding $\ell\left(\mu_{i}\right) \sim \sqrt{\varepsilon}$. On the other hand, each step $\mu^{* 1 / \varepsilon}$ only visits $1 / \varepsilon$ points, hence the measure $\left(\mu^{* 1 / \varepsilon} * \nu\right)^{\varepsilon N}$ is supported by roughly $(1 / \varepsilon)^{\varepsilon N}$ points, yielding $h\left(\mu_{i}\right) \sim \varepsilon|\log \varepsilon|$. In particular, $h\left(\mu_{i}\right)=o\left(\ell\left(\mu_{i}\right)\right)$. This implies that $h\left(\mu_{i}\right) / \ell\left(\mu_{i}\right)$, which tends to 0 , can not tend to the maximum $M$. Therefore, this case can not happen.

The rigorous argument is considerably more delicate. One difficulty is that $\mu_{i}$ does not decompose in general as $(1-\varepsilon) \mu+\varepsilon \nu$ : there can be in $\mu_{i}$ points with a very small probability (which are not seen by $\mu$ ), but much larger than $\varepsilon$, the probability to visit a nonelementary subset of $\Gamma$. These points will play an important role on the relevant time scale, i.e., $1 / \varepsilon$. Hence, we have to describe the different time scales that happen in $\mu_{i}$.

For each $a \in \Sigma$, we have a weight $\mu_{i}(a)$, which tends to 0 if $a$ is not in the support of $\mu$. Reordering the $a_{k}$ and extracting a subsequence, we can assume that $\Sigma=\left\{a_{1}, \ldots, a_{p}\right\}$ with $\mu_{i}\left(a_{1}\right) \geqslant \cdots \geqslant \mu_{i}\left(a_{p}\right)$ (and $\left.a_{1}=e\right)$. Extracting a further subsequence, we may also assume that $\mu_{i}\left(a_{k}\right) / \mu_{i}\left(a_{k-1}\right)$ converges for all $k$, towards a limit in $[0,1]$.

Let $\Gamma_{k}$ be the subgroup generated by $a_{1}, \ldots, a_{k}$. We consider the smallest $r$ such that $\Gamma_{r}$ is non-elementary. Then, we consider the biggest $s<r$ such that $\mu_{i}(r)=o\left(\mu_{i}(s)\right)$. Roughly speaking, the random walk has enough time to spread on the elementary subgroup $\Gamma_{s}$, before seeing $a_{r}$. It turns out that the asymptotic behavior will depend on the nature of $\Gamma_{s}$ (finite or virtually cyclic infinite).

We will decompose the measure $\mu_{i}$ as the sum of two components $\left(1-\varepsilon_{i}\right) \alpha_{i}+\varepsilon_{i} \beta_{i}$, where $\varepsilon_{i}$ tends to 0 , the measure $\alpha_{i}$ mainly lives on $\Gamma_{s}$, and the measure $\beta_{i}$ corresponds to the remaining part of $\mu_{i}$, on $\left\{a_{s+1}, \ldots, a_{p}\right\}$. The precise construction depends on the nature of $\Gamma_{s}:$

- If $\Gamma_{s}$ is finite. Let $\beta_{i}^{(0)}$ be the normalized restriction of $\mu_{i}$ to $\left\{a_{s+1}, \ldots, a_{p}\right\}$. To avoid periodicity problems, we rather consider $\beta_{i}=\left(\delta_{e}+\beta_{i}^{(0)}\right) / 2$. We decompose $\mu_{i}=\left(1-\varepsilon_{i}\right) \alpha_{i}+\varepsilon_{i} \beta_{i}$, where $\alpha_{i}$ is supported on $a_{1}, \ldots, a_{s}$. By construction, the probability of any element in the support of $\alpha_{i}$ is much bigger than $\varepsilon_{i}$.
- If $\Gamma_{s}$ is virtually cyclic infinite. The group $\Gamma_{s}$ contains a hyperbolic element $g_{0}$, with repelling and attracting points at infinity denoted by $g_{0}^{-}$and $g_{0}^{+}$. The elements of $\Gamma_{s}$ all fix the set $\left\{g_{0}^{-}, g_{0}^{+}\right\}$. We take for $\alpha_{i}$ the normalized restriction of $\mu_{i}$ to those elements in $\Sigma$ that fix $\left\{g_{0}^{-}, g_{0}^{+}\right\}$, and for $\beta_{i}$ the normalized restriction of $\mu_{i}$ to the other elements. Once again, we can write $\mu_{i}=\left(1-\varepsilon_{i}\right) \alpha_{i}+\varepsilon_{i} \beta_{i}$.
In both cases, $\varepsilon_{i}$ is comparable to the probability $\mu_{i}\left(a_{r}\right)$, and is therefore negligible with respect to $\mu_{i}\left(a_{s}\right)$. We will write $\mu_{i}=\mu_{\varepsilon}$ (and, in the same way, we will replace all indices $i$ with $\varepsilon$, since the main parameter is $\varepsilon=\varepsilon_{i}$ ). The measure $\mu_{\varepsilon}$ converges to $\mu$ when $\varepsilon$ tends to 0 , while $\beta_{\varepsilon}$ tends to a probability measure $\beta$, supported on $e, a_{s+1}, \ldots, a_{p}$. If the measures $\mu_{\varepsilon}$ are symmetric to begin with, the measures $\alpha_{\varepsilon}$ and $\beta_{\varepsilon}$ are also symmetric by construction.

To generate the random walk given by $\mu_{\varepsilon}$, one can first independently choose random measures $\rho_{n}$ : one takes $\rho_{n}=\alpha_{\varepsilon}$ with probability $1-\varepsilon$, and $\rho_{n}=\beta_{\varepsilon}$ with probability $\varepsilon$. Then, one chooses elements $g_{n}$ randomly according to $\rho_{n}$, and one multiplies them: the product $g_{1} \cdots g_{n}$ is distributed like the random walk given by $\mu_{\varepsilon}$ at time $n$.

We will group together successive $g_{k}$, into blocks where the equidistribution on $\Gamma_{s}$ can be seen. More precisely, denote by $t_{1}, t_{2}, \ldots$ the successive times where $\rho_{n}=\beta_{\varepsilon}$ (and $t_{0}=0$ ). They are stopping times, the successive differences are independent and identically distributed, with a geometric distribution of parameter $\varepsilon$ (i.e., $\mathbb{P}\left(t_{1}=n\right)=(1-\varepsilon)^{n-1} \varepsilon$ ), with mean $1 / \varepsilon$. Write $L_{N}=g_{t_{N-1}+1} \cdots g_{t_{N}}$. By construction, the $L_{i}$ are independent, identically distributed, and the random walk they define, i.e., $L_{1} \cdots L_{N}$, is a subsequence of the original random walk $g_{1} \cdots g_{n}$. Let $\lambda_{\varepsilon}$ be the distribution of $L_{i}$ on $\Gamma$, i.e.,

$$
\lambda_{\varepsilon}=\sum_{n=0}^{\infty}(1-\varepsilon)^{n} \varepsilon \alpha_{\varepsilon}^{* n} * \beta_{\varepsilon} .
$$

Lemma 5.1. The measure $\lambda_{\varepsilon}$ has finite first moment and finite time one entropy. Moreover, $\ell\left(\mu_{\varepsilon}\right)=\varepsilon \ell\left(\lambda_{\varepsilon}\right)$ and $h\left(\mu_{\varepsilon}\right)=\varepsilon h\left(\lambda_{\varepsilon}\right)$.
Proof. As the mean of $t_{1}$ is $1 / \varepsilon$, the random walk generated by $\lambda_{\varepsilon}$ is essentially the random walk generated by $\mu_{\varepsilon}$, but on a time scale $1 / \varepsilon$. This justifies heuristically the statement.

For the rigorous proof, let us first check that $\lambda_{\varepsilon}$ has finite first moment (and hence finite time one entropy). Since all the measures have finite support, we have $\left|L_{1}\right| \leqslant C t_{1}$. Since a geometric distribution has moments of all order, the same is true for $\left|L_{1}\right|$.

The strong law of large numbers ensures that, almost surely, $t_{N} \sim N / \varepsilon$. Therefore, almost surely,

$$
\ell\left(\lambda_{\varepsilon}\right)=\lim \frac{\left|L_{1} \cdots L_{N}\right|}{N}=\lim \frac{\left|g_{1} \cdots g_{t_{N}}\right|}{N}=\lim \frac{\left|g_{1} \cdots g_{t_{N}}\right|}{t_{N}} \cdot \frac{t_{N}}{N}=\ell\left(\mu_{\varepsilon}\right) \cdot 1 / \varepsilon .
$$

This proves the statement of the lemma for the drift.
For the entropy, we use the characterization of Lemma 2.4. We will show that $h\left(\mu_{\varepsilon}\right) \leqslant$ $\varepsilon h\left(\lambda_{\varepsilon}\right)$ and $h\left(\mu_{\varepsilon}\right) \geqslant \varepsilon h\left(\lambda_{\varepsilon}\right)$. Let $K_{n}$ be a set of cardinality at most $e^{\left(h\left(\mu_{\varepsilon}\right)+\eta\right) n}$ which contains $g_{1} \cdots g_{n}$ with probability at least $1 / 2$. Let $N=\varepsilon n$. With large probability, $t_{N}$ is close to $n$, up to $\eta^{\prime} n$ (where $\eta^{\prime}$ is arbitrarily small). Hence, with probability at least $1 / 3$, the point $L_{1} \cdots L_{N}$ belongs to the $C \eta^{\prime} n$-neighborhood of $K_{n}$, whose cardinality is at most

$$
\left|K_{n}\right| \cdot e^{C^{\prime} \eta^{\prime} n} \leqslant e^{\left(h\left(\mu_{\varepsilon}\right)+\eta+C^{\prime} \eta^{\prime}\right) n}=e^{\left(h\left(\mu_{\varepsilon}\right)+\eta+C^{\prime} \eta^{\prime}\right) N / \varepsilon} .
$$

As $\eta$ and $\eta^{\prime}$ are arbitrary, this shows that $h\left(\lambda_{\varepsilon}\right) \leqslant h\left(\mu_{\varepsilon}\right) / \varepsilon$. The converse inequality is proved in the same way.

The previous lemma shows that we should understand $\lambda_{\varepsilon}$. We define an auxiliary probability measure $\tilde{\alpha}_{\varepsilon}$ so that $\lambda_{\varepsilon}=\tilde{\alpha}_{\varepsilon} * \beta_{\varepsilon}$, by

$$
\begin{equation*}
\tilde{\alpha}_{\varepsilon}=\sum_{n=0}^{\infty}(1-\varepsilon)^{n} \varepsilon \alpha_{\varepsilon}^{* n} . \tag{5.1}
\end{equation*}
$$

In this formula, most weight is concentrated around those $n$ of the order of $1 / \varepsilon$. Hence, we have to understand the iterates of $\alpha_{\varepsilon}$ in time $1 / \varepsilon$. When $\Gamma_{s}$ is finite, we will see that it has enough time to equidistribute on $\Gamma_{s}$ (even though $\alpha_{\varepsilon}$ may give a very small weight to some elements, this weight is by construction much larger than $\varepsilon$, so that $1 / \varepsilon$ iterates are enough to equidistribute). When $\Gamma_{s}$ is virtually cyclic, we will see that the random walk has enough time to drift away significantly from the identity.

In both cases, we will need quantitative results on basic groups, but in weakly elliptic cases (i.e., the transition probabilities are not bounded from below). There are techniques to get quantitative estimates in such settings, especially comparison techniques (due for instance to Varopoulos, Diaconis, Saloff-Coste): one can compare weakly elliptic walks to elliptic ones (which we understand well) thanks to Dirichlet forms arguments: these arguments make it possible to transfer results from the latter to the former (modulo some loss in the constants, due to the lack of ellipticity). We will rely on such results when $\Gamma_{s}$ is infinite. When it is finite, such techniques can also be used, but we will rather give a more elementary argument.

We start with the case where $\Gamma_{s}$ is finite. We need to quantify the speed of convergence to the stationary measure in finite groups, with the following lemma.

Lemma 5.2. Let $\Lambda$ be a finite group. Let $\Sigma_{\Lambda} \subset \Lambda$ be a generating subset (it does not have to be symmetric). Let $\pi_{\Lambda}$ be the uniform measure on $\Lambda$, and let $d\left(\mu, \pi_{\Lambda}\right)$ be the euclidean distance between a measure $\mu$ and $\pi_{\Lambda}$ (i.e., $\left.\left(\sum\left(\mu(g)-\pi_{\Lambda}(g)\right)^{2}\right)^{1 / 2}\right)$. For any $\delta>0$, there exists $K>0$ with the following property. Let $\eta>0$. Consider a probability measure $\mu$ on $\Lambda$ with $\mu(\sigma) \geqslant \eta$ for any $\sigma \in \Sigma_{\Lambda} \cup\{e\}$. Then, for all $n>K / \eta$,

$$
d\left(\mu^{* n}, \pi_{\Lambda}\right) \leqslant \delta
$$

In other words, the time to see the equidistribution towards the stationary measure is bounded by $1 / \eta$, where $\eta$ is the minimum of the transition probabilities on $\Sigma_{\Lambda}$.

Proof. Endow the space $\mathcal{M}(\Lambda)$ of signed measures on $\Lambda$ with the scalar product corresponding to the quadratic form $|\nu|^{2}=\sum \nu(g)^{2}$. Denote by $H=\left\{\nu: \sum \nu(g)=0\right\}$ the hyperplane $\pi_{\Lambda}^{\perp}$ of zero mass measures. For any probability $\rho$, denote by $M_{\rho}$ the left-convolution operator on $\mathcal{M}(\Lambda)$, that is $M_{\rho}(\nu)=\rho * \nu$. Since convolution preserves mass, $H$ is $M_{\rho}$-invariant. Let us prove that the operator norm of $M_{\rho}$ is bounded by 1 . Indeed, put $u_{\rho}(g)=\sum_{h \in \Lambda} \rho(h) \rho(h g)$,
this is a probability on $\Lambda$. We have

$$
\begin{aligned}
\left|M_{\rho} \nu\right|^{2} & =\sum_{g \in \Lambda}\left(M_{\rho} \nu(g)\right)^{2}=\sum_{\left(g, h_{1}, h_{2}\right) \in \Lambda^{3}} \rho\left(g h_{1}^{-1}\right) \rho\left(g h_{2}^{-1}\right) \nu\left(h_{1}\right) \nu\left(h_{2}\right) \\
& =\sum_{\left(h_{1}, h_{2}\right) \in \Lambda^{2}} \nu\left(h_{1}\right) \nu\left(h_{2}\right) u_{\rho}\left(h_{1} h_{2}^{-1}\right)=\sum_{(g, h) \in \Lambda^{2}} \nu(h) \nu\left(g^{-1} h\right) u_{\rho}(g) \\
& \leqslant \sum_{g \in \Lambda}|\nu|^{2} u_{\rho}(g)=|\nu|^{2} .
\end{aligned}
$$

This proves that $\left\|M_{\rho}\right\| \leqslant 1$. Now fix $\rho_{o}$ to be the uniform probability on the set $\Sigma_{\Lambda} \cup\{e\}$. Notice that $u_{\rho_{o}}(g)>0$ for any $g \in \Sigma_{\Lambda} \cup\{e\}$, since $\rho_{o}(e)>0$. We claim that $M_{\rho_{o}}$ restricted to $H$ has an operator norm $c<1$. Would it be not the case, there would exist $\nu \in H-\{0\}$ such that the previous inequalities would be equalities. Thanks to the equality case in the Cauchy-Schwarz inequality, this implies that, for any $g \in \Sigma_{\Lambda}$, the two measures $h \mapsto \nu(h)$ and $h \mapsto \nu\left(g^{-1} h\right)$ are positively proportional. Since their norm are equal, they must be equal. Since $\Sigma_{\Lambda}$ generates $\Lambda, \nu$ is $\Lambda$-invariant and belongs to $H$, so it must be zero.

By assumption, the probability $\mu$ can be decomposed as

$$
\mu=\eta \rho_{o}+(1-\eta) \nu
$$

where $\nu$ is some probability. This implies that $M_{\mu}$ restricted to $H$ has operator norm at most $\eta c+(1-\eta)$. Therefore,

$$
d\left(\mu^{* n}, \pi_{\Lambda}\right)=\left|\mu^{* n}-\pi_{\Lambda}\right|=\left|M_{\mu}^{n}\left(\delta_{e}-\pi_{\Lambda}\right)\right| \leqslant 2(1-(1-c) \eta)^{n}
$$

This inequality implies the result.
We can now describe the asymptotic behavior of $\mu_{\varepsilon}$ when the group $\Gamma_{s}$ is finite.
Lemma 5.3. Assume that $\Gamma_{s}$ is finite. Define a new probability measure $\lambda=\pi_{\Gamma_{s}} * \beta$ (it generates a non-elementary subgroup). When $\varepsilon$ tends to 0 , we have $h\left(\mu_{\varepsilon}\right) \sim \varepsilon h(\lambda)$ and $\ell\left(\mu_{\varepsilon}\right) \sim \varepsilon \ell(\lambda)$.
Proof. The random variable $t_{1}$, being geometric of parameter $\varepsilon$, is of the order of $1 / \varepsilon$ with high probability (i.e., for any $\delta>0$, there exists $u>0$ such that $\left.\mathbb{P}\left(t_{1}>u / \varepsilon\right) \geqslant 1-\delta\right)$. Writing $\Sigma_{s}=\left\{a_{1}, \ldots, a_{s}\right\}$ for the support of $\alpha_{\varepsilon}$, we have $\min _{\sigma \in \Sigma_{s}} \alpha_{\varepsilon}(\sigma)=(1-\varepsilon)^{-1} \mu_{\varepsilon}\left(a_{s}\right)$, which is much bigger than $\varepsilon$ by definition of $s$. Lemma 5.2 shows that the measures $\alpha_{\varepsilon}^{* n}$ are close to $\pi_{\Gamma_{s}}$ for $n \geqslant u / \varepsilon$. This implies that $\tilde{\alpha}_{\varepsilon}$ (defined in (5.1)) converges to $\pi_{\Gamma_{s}}$ when $\varepsilon \rightarrow 0$. As $\beta_{\varepsilon}$ converges to $\beta$, this shows that $\lambda_{\varepsilon}$ converges to $\lambda$.

The support of the measure $\lambda$ contains $\Gamma_{s}$ and $a_{s+1}, \ldots, a_{r}$ (as the support of $\beta$ contains $\left\{e, a_{s+1}, \ldots, a_{r}\right\}$ by construction). Hence, $\Gamma_{\lambda}$ contains the non-elementary subgroup $\Gamma_{r}$. It follows that the entropy and the drift are continuous at $\lambda$, by Proposition 2.3 and Theorem 2.9. We get $h\left(\lambda_{\varepsilon}\right) \rightarrow h(\lambda)$ and $\ell\left(\lambda_{\varepsilon}\right) \rightarrow \ell(\lambda)$. With Lemma 5.1, this completes the proof.

We deduce from the lemma that $h\left(\mu_{\varepsilon}\right) / \ell\left(\mu_{\varepsilon}\right)$ tends to $h(\lambda) / \ell(\lambda)$. Hence, the measure $\mu_{\Sigma}=\lambda$ satisfies the conclusion of the theorem, at least in the non-symmetric case. In the symmetric case, where we are looking for a symmetric measure $\mu_{\Sigma}$, the measure $\lambda=\pi_{\Gamma_{s}} * \beta$ is not an answer to the problem. However, $\lambda^{\prime}=\pi_{\Gamma_{s}} * \beta * \pi_{\Gamma_{s}}$ is symmetric, and it clearly
has the same entropy and drift as $\lambda$ (since $\pi_{\Gamma_{s}} * \pi_{\Gamma_{s}}=\pi_{\Gamma_{s}}$ ). Hence, we can take $\mu_{\Sigma}=\lambda^{\prime}$. This completes the proof of Theorem 1.7 when the group $\Gamma_{s}$ is finite.

Example 5.4. Let $\Gamma=\mathbb{Z} / 2 * \mathbb{Z} / 4$, with $\Sigma=\left\{a, b, b^{-1}\right\}$ (where $a$ is the generator of $\mathbb{Z} / 2$ and $b$ the generator of $\mathbb{Z} / 4$ ), with the word distance coming from $\Sigma$. [MM07, Section 5.1] shows that the supremum over measures supported on $\Sigma$ of $h(\mu) / \ell(\mu)$ is the growth $v$ of the group (note that $\Gamma$ is virtually free), and that it is not realized by a measure supported on $\Sigma$. This shows that, in Theorem 1.7, the fact that $\mu_{\Sigma}$ may need a support larger than $\Sigma$ is not an artefact of the proof.

In this example, any symmetric measure on $\Sigma$ is of the form $\mu_{\varepsilon}=(1-\varepsilon) \delta_{a}+\varepsilon \beta$ where $\beta$ is uniform on $\left\{b, b^{-1}\right\}$. The above proof shows that, when $\varepsilon$ tends to $0, h\left(\mu_{\varepsilon}\right) / \ell\left(\mu_{\varepsilon}\right)$ converges to $h(\lambda) / \ell(\lambda)$ where $\lambda=\pi_{\Gamma_{s}} * \beta=\frac{1}{2}\left(\delta_{e}+\delta_{a}\right) * \frac{1}{2}\left(\delta_{b}+\delta_{b^{-1}}\right)$ is the uniform measure on $\left\{b, b^{-1}, a b, a b^{-1}\right\}$.

It remains to treat the case where $\Gamma_{s}$ is virtually cyclic infinite. Such a group surjects onto $\mathbb{Z}$ or $\mathbb{Z} \rtimes \mathbb{Z} / 2$ (the infinite dihedral group), with finite kernel. From the point of view of the random walk, most things happen in the quotient. Hence, it would suffice to understand these two groups (separating in the case of $\mathbb{Z}$ the centered and non-centered cases). We will rather give direct arguments which do not use this reduction and which avoid separating cases. Let $t \leqslant s$ be the smallest index such that $\left\{a_{1}, \cdots, a_{t}\right\}$ generates an infinite group. Let $\eta=\eta(\varepsilon)=\mu_{\varepsilon}\left(a_{t}\right)$, this parameter governs the equidistribution speed on $\Gamma_{s}$ (or, at least, on $\Gamma_{t}$, which has finite index in $\Gamma_{s}$ since these two groups are virtually cyclic infinite). We will find the asymptotics of the entropy and the drift in terms of $\eta / \varepsilon$ (which tends to infinity by definition of $s$ ). We start with the entropy (for which an upper bound suffices). Note that the random walk directed by $\alpha_{\varepsilon}$ does not live on $\Gamma_{s}$, but on a possibly bigger group since we have put in $\alpha_{\varepsilon}$ all the points that fix the set $\left\{g_{0}^{-}, g_{0}^{+}\right\}$(this will be important in the control of the drift below). Let $\tilde{\Gamma}_{s}$ be the group they generate, it is still virtually cyclic (see, for instance, [GdlH90, Théorème 37 page 157]), and it contains $\Gamma_{s}$ as a finite index subgroup.

Lemma 5.5. There exists a constant $C$ such that $h\left(\lambda_{\varepsilon}\right) \leqslant C \log (\eta / \varepsilon)$.
Proof. Let $K$ be the group generated by $\left\{a_{1}, \ldots, a_{t-1}\right\}$. It is finite by definition of $t$. Let $\Sigma^{\prime}$ be the set of points among $a_{t}, \ldots, a_{p}$ which stabilize $\left\{g_{0}^{-}, g_{0}^{+}\right\}$. The group $\tilde{\Gamma}_{s}$ is generated by $K$ and $\Sigma^{\prime}$. Let us consider the associated word pseudo-distance $d^{\prime}$, where we decide that elements in $K$ have 0 length. This pseudo-distance is quasi-isometric to the usual distance, and it satisfies $d^{\prime}(e, x k)=d^{\prime}(e, x)$ for all $x \in \tilde{\Gamma}_{s}$ and all $k \in K$.

Let us first estimate the average distance to the origin for an element given by $\tilde{\alpha}_{\varepsilon}$. We decompose $\alpha_{\varepsilon}$ as the average of a measure supported on $\left\{a_{1}, \ldots, a_{t-1}\right\} \subset K$, and of a measure supported on $\Sigma^{\prime}$ (the contribution of the latter has a mass $m(\varepsilon)$ bounded by ( $p-$ $t+1) \eta \leqslant C \eta$ ). The measure $\alpha_{\varepsilon}^{* n}$ can be obtained by picking at each step one of these two measures (according to their respective weight), and then jumping according to a random element for this measure. When we use the first measure, the $d^{\prime}$-distance to the origin does not change by definition. Hence, the distance to the origin is bounded by the number of
choices of the second measure. We obtain

$$
\begin{aligned}
\mathbb{E}_{\tilde{\alpha}_{\varepsilon}}\left(d^{\prime}(e, g)\right) & \leqslant \sum_{n=0}^{\infty}(1-\varepsilon)^{n} \varepsilon \sum_{i=0}^{n}\binom{n}{i} m(\varepsilon)^{i}(1-m(\varepsilon))^{n-i} \cdot C i \\
& =C m(\varepsilon) \sum_{n=0}^{\infty}(1-\varepsilon)^{n} \varepsilon \sum_{i=1}^{n} n\binom{n-1}{i-1} m(\varepsilon)^{i-1}(1-m(\varepsilon))^{n-i} \\
& =C m(\varepsilon) \sum_{n=0}^{\infty}(1-\varepsilon)^{n} \varepsilon n=C m(\varepsilon)(1-\varepsilon) / \varepsilon \leqslant C \eta / \varepsilon .
\end{aligned}
$$

A measure supported on the integers with first moment $A$ has entropy bounded by $C \log A+C$ (see, for instance, [EK10, Lemma 2]). The proof also applies to virtually cyclic situations (the finite thickening does not change anything). Therefore, we get $H\left(\tilde{\alpha}_{\varepsilon}\right) \leqslant$ $C \log (\eta / \varepsilon)+C$.

Finally,

$$
H\left(\lambda_{\varepsilon}\right)=H\left(\tilde{\alpha}_{\varepsilon} * \beta_{\varepsilon}\right) \leqslant H\left(\tilde{\alpha}_{\varepsilon}\right)+H\left(\beta_{\varepsilon}\right) \leqslant C \log (\eta / \varepsilon)+C
$$

since the support of $\beta_{\varepsilon}$ is uniformly bounded. As $\eta / \varepsilon \rightarrow \infty$, this gives $H\left(\lambda_{\varepsilon}\right) \leqslant C \log (\eta / \varepsilon)$. Finally, we estimate $h\left(\lambda_{\varepsilon}\right)=\inf _{n>0} H\left(\lambda_{\varepsilon}^{* n}\right) / n \leqslant H\left(\lambda_{\varepsilon}\right)$ to get the conclusion of the lemma.

For the drift, we need to be more precise since we need a lower bound to conclude. We will use a lemma giving lower bounds on the equidistribution speed in virtually cyclic infinite groups, using comparison techniques.

Lemma 5.6. Let $\Lambda$ be a virtually cyclic infinite group. Let $\Sigma_{\Lambda} \subset \Lambda$ be a finite subset generating an infinite subgroup of $\Lambda$. There exists a constant $C$ with the following property. Let $\eta>0$. Let $\mu$ be a probability measure on $\Lambda$ with $\mu(e) \geqslant 1 / 2$ and $\mu(\sigma) \geqslant \eta$ for any $\sigma \in \Sigma_{\Lambda}$. Then, for all $n \geqslant 1$,

$$
\sup _{g \in \Lambda} \mu^{* n}(g) \leqslant C(\eta n)^{-1 / 2}
$$

The interest of the lemma is that $C$ does not depend on the measure $\mu$, and that we obtain an explicit control on $\mu^{* n}$ just in terms of a lower bound on the transition probabilities of $\mu$.
Proof. We use the comparison method. Let $\rho$ be the uniform measure on $e, \Sigma_{\Lambda}$ and $\Sigma_{\Lambda}^{-1}$. The random walk it generates does not have to be transitive (since $\Sigma_{\Lambda}$ does not necessarily generate the whole group $\Lambda$ ), but $\Lambda$ is partitioned into finitely many classes where it is transitive (and isomorphic to the random walk on the group generated by $\Sigma_{\Lambda}$ ). Moreover, it is symmetric, and therefore reversible for the counting measure $m$ on $\Lambda$. The Dirichlet form associated to $\rho$ is by definition

$$
\mathcal{E}_{\rho}(f, f)=\frac{1}{2} \sum_{x, y}|f(x)-f(y)|^{2} \rho\left(x^{-1} y\right)
$$

for any $f: \Lambda \rightarrow \mathbb{C}$. As $\Lambda$ has linear growth, the following Nash inequality holds (see, for instance, [Woe00, Proposition 14.1]).

$$
\|f\|_{L^{2}}^{6} \leqslant C\|f\|_{L^{1}}^{4} \mathcal{E}_{\rho}(f, f)
$$

where all norms are defined with respect to the measure $m$ on $\Lambda$. Let $P_{\mu}$ be the Markov operator associated to $\mu$. It satisfies

$$
\|f\|_{L^{2}}^{2}-\left\|P_{\mu} f\right\|_{L^{2}}^{2}=\langle f, f\rangle-\left\langle P_{\mu} f, P_{\mu} f\right\rangle=\left\langle\left(I-P_{\mu}^{*} P_{\mu}\right) f, f\right\rangle .
$$

The operator $P_{\mu}^{*} P_{\mu}$ is the Markov operator associated to the symmetric probability measure $\nu=\check{\mu} * \mu$, which satisfies $\nu(\sigma) \geqslant \eta / 2$ for $\sigma \in \Sigma_{\Lambda} \cup \Sigma_{\Lambda}^{-1}$ and $\nu(e) \geqslant 1 / 4$ (since $\mu(e) \geqslant 1 / 2$ ). Therefore, $\rho(g) \leqslant C \eta^{-1} \nu(g)$ for all $g$. We deduce

$$
\begin{aligned}
\|f\|_{L^{2}}^{2}-\left\|P_{\mu} f\right\|_{L^{2}}^{2} & =\sum \overline{f(x)}(f(x)-f(y)) \nu\left(x^{-1} y\right)=\frac{1}{2} \sum|f(x)-f(y)|^{2} \nu\left(x^{-1} y\right) \\
& \geqslant \frac{\eta}{2 C} \sum|f(x)-f(y)|^{2} \rho\left(x^{-1} y\right)=\frac{\eta}{C} \mathcal{E}_{\rho}(f, f) .
\end{aligned}
$$

Combining this inequality with Nash inequality, we obtain

$$
\|f\|_{L^{2}}^{6} \leqslant C \eta^{-1}\|f\|_{L^{1}}^{4}\left(\|f\|_{L^{2}}^{2}-\left\|P_{\mu} f\right\|_{L^{2}}^{2}\right) .
$$

The operator $P_{\mu}^{*}$ satisfies the same inequality, for the same reason. Composing these inequalities, we obtain an estimate for the norm of $P_{\mu}^{n}$ from $L^{1}$ to $L^{\infty}$ (this is [VSCC92, Lemma VII.2.6]), of the form

$$
\left\|P_{\mu}^{n}\right\|_{L^{1} \rightarrow L^{\infty}} \leqslant\left(C^{\prime} \eta^{-1} / n\right)^{1 / 2} .
$$

Applying this inequality to the function $\delta_{e}$, we get the desired result.
The previous lemma implies that, if $C^{\prime}$ is large enough, a neighborhood of size $(\eta n)^{1 / 2} / C^{\prime}$ of the identity has probability for $\mu^{* n}$ at most $1 / 2$. Hence, the average distance to the origin is at least of the order of $(\eta n)^{1 / 2}$.

Now, we study the stationary measure for $\beta_{\varepsilon} * \tilde{\alpha}_{\varepsilon}$ on $\partial \Gamma$. We recall that $g_{0}$ is a hyperbolic element in $\Gamma_{s}$, fixed once and for all.

Lemma 5.7. There exists a neighborhood $U$ of $\left\{g_{0}^{-}, g_{0}^{+}\right\}$in $\partial \Gamma$ such that the stationary measure $\nu_{\varepsilon}$ of $\beta_{\varepsilon} * \tilde{\alpha}_{\varepsilon}$ satisfies $\nu_{\varepsilon}(U) \rightarrow 0$.

Proof. Let us first show that, for any neighborhood $U$ of $\left\{g_{0}^{-}, g_{0}^{+}\right\}$, then $\left(\tilde{\alpha}_{\varepsilon} * \delta_{z}\right)\left(U^{c}\right)$ tends to 0 , uniformly in $z \in \partial \Gamma$. This is not surprising since a typical element for $\tilde{\alpha}_{\varepsilon}$ is large in the virtually cyclic group $\tilde{\Gamma}_{s}$, and sends most points into $U$. To make this argument rigorous, we will use Lemma 5.6. The definition (5.1) shows that it suffices to prove that $\left(\alpha_{\varepsilon}^{* n} * \delta_{z}\right)\left(U^{c}\right)$ is small for $n \geqslant u / \varepsilon$.

The subgroup generated by $g_{0}$ has finite index in $\tilde{\Gamma}_{s}$. Hence, any element in $\tilde{\Gamma}_{s}$ can be written as $g_{0}^{k} \gamma_{i}$, for $\gamma_{i}$ in a finite set. Thus, the measure $\alpha_{\varepsilon}^{* n}$ can be written as $\sum c_{n}(k, i) \delta_{g_{0}^{k} \gamma_{i}}$, for some coefficients $c_{n}(k, i)$. Lemma 5.6 (applied to $\Lambda=\tilde{\Gamma}_{s}$ with $\Sigma_{\Lambda}=\left\{a_{1}, \ldots, a_{t}\right\}$ ) ensures that $\sup _{k, i} c_{n}(k, i) \leqslant C /(\eta n)^{1 / 2}$. When $n \geqslant u / \varepsilon$, this quantity tends to 0 since $\varepsilon=o(\eta)$. We have

$$
\left(\alpha_{\varepsilon}^{* n} * \delta_{z}\right)\left(U^{c}\right)=\sum_{k, i} c_{n}(k, i) 1\left(g_{0}^{k} \gamma_{i} z \notin U\right) .
$$

As the element $g_{0}$ is hyperbolic, there exists $C$ such that, for any $w \in \partial \Gamma$,

$$
\left|\left\{k \in \mathbb{Z}: g_{0}^{k} w \notin U\right\}\right| \leqslant C
$$

The uniformity in $w$ follows from the compactness of $\left(\partial \Gamma \backslash\left\{g_{0}^{-}, g_{0}^{+}\right\}\right) /\left\langle g_{0}\right\rangle$. We obtain

$$
\left(\alpha_{\varepsilon}^{* n} * \delta_{z}\right)\left(U^{c}\right) \leqslant\left(\sup _{k, i} c_{n}(k, i)\right) \sum_{i}\left|\left\{k \in \mathbb{Z}: g_{0}^{k} \gamma_{i} z \notin U\right\}\right| \leqslant C \sup _{k, i} c_{n}(k, i) \leqslant C /(\eta n)^{1 / 2} .
$$

This shows that $\left(\alpha_{\varepsilon}^{* n} * \delta_{z}\right)\left(U^{c}\right)$ is small, as desired.
As $\tilde{\alpha}_{\varepsilon} * \delta_{z}\left(U^{c}\right)$ tends to 0 uniformly in $z$, we deduce that $\left(\tilde{\alpha}_{\varepsilon} * \nu_{\varepsilon}\right)\left(U^{c}\right)$ also tends to 0 , and therefore that $\left(\tilde{\alpha}_{\varepsilon} * \nu_{\varepsilon}\right)(U)$ tends to 1 .

Let $A=\left\{g_{0}^{-}, g_{0}^{+}\right\}$. We claim that, for all $g$ such that $g A \cap A \neq \emptyset$, then $g A=A$. Indeed, if $g\left(g_{0}^{-}\right) \in A$ for instance, then $g^{-1} g_{0} g$ is a hyperbolic element stabilizing $g_{0}^{-}$. It also stabilizes $g_{0}^{+}$, by [GdIH90, Théorème 30 page 154], i.e., $g_{0} g\left(g_{0}^{+}\right)=g\left(g_{0}^{+}\right)$. Hence, $g\left(g_{0}^{+}\right)$is a fixed point of $g_{0}$, i.e., $g\left(g_{0}^{+}\right) \in A$.

By definition of $\beta_{\varepsilon}$, the finitely many elements of its support do not fix $A$. They even satisfy $g A \cap A=\emptyset$ for all $g$ in this support, by the previous argument. If $U$ is small enough, we get $g U \cap U=\emptyset$, i.e., $g(U) \subset U^{c}$.

Finally,

$$
\nu_{\varepsilon}\left(U^{c}\right)=\left(\beta_{\varepsilon} * \tilde{\alpha}_{\varepsilon} * \nu_{\varepsilon}\right)\left(U^{c}\right) \geqslant\left(\tilde{\alpha}_{\varepsilon} * \nu_{\varepsilon}\right)(U),
$$

which tends to 1 when $\varepsilon$ tends to 0 .
Lemma 5.8. The drift $\ell\left(\lambda_{\varepsilon}\right)$ satisfies $\ell\left(\lambda_{\varepsilon}\right) \geqslant c \cdot(\eta / \varepsilon)^{1 / 2}$.
Proof. Let $\rho_{\varepsilon}$ be a stationary measure for $\lambda_{\varepsilon}$, on the Busemann boundary $\partial_{B} \Gamma$. By Proposition 2.2,

$$
\ell\left(\lambda_{\varepsilon}\right)=\int c_{B}(g, \xi) \mathrm{d} \rho_{\varepsilon}(\xi) \mathrm{d} \lambda_{\varepsilon}(g),
$$

where $c_{B}(g, \xi)=h_{\xi}\left(g^{-1}\right)$ is the Busemann cocycle. As $\lambda_{\varepsilon}=\tilde{\alpha}_{\varepsilon} * \beta_{\varepsilon}$, this gives

$$
\ell\left(\lambda_{\varepsilon}\right)=\int c_{B}(L b, \xi) \mathrm{d} \rho_{\varepsilon}(\xi) \mathrm{d} \tilde{\alpha}_{\varepsilon}(L) \mathrm{d} \beta_{\varepsilon}(b) .
$$

With the cocycle relation (2.2), this becomes

$$
\ell\left(\lambda_{\varepsilon}\right)=\int c_{B}(L, b \xi) \mathrm{d} \rho_{\varepsilon}(\xi) \mathrm{d} \tilde{\alpha}_{\varepsilon}(L) \mathrm{d} \beta_{\varepsilon}(b)+\int c_{B}(b, \xi) \mathrm{d} \rho_{\varepsilon}(\xi) \mathrm{d} \tilde{\alpha}_{\varepsilon}(L) \mathrm{d} \beta_{\varepsilon}(b) .
$$

The second integral is bounded independently of $\varepsilon$ since the support of $\beta_{\varepsilon}$ is finite. In the first integral, $\xi^{\prime}=b \xi$ is distributed according to the measure $\tilde{\rho}_{\varepsilon}:=\beta_{\varepsilon} * \rho_{\varepsilon}$, which is stationary for $\beta_{\varepsilon} * \tilde{\alpha}_{\varepsilon}$. Lemma 5.7 implies that its projection $\left(\pi_{B}\right)_{*} \tilde{\rho}_{\varepsilon}$ on the geometric boundary, which is again stationary for $\beta_{\varepsilon} * \tilde{\alpha}_{\varepsilon}$, gives a small measure to a neighborhood $U$ of $\left\{g_{0}^{-}, g_{0}^{+}\right\}$.

As the limit set of $\tilde{\Gamma}_{s}$ is $\left\{g_{0}^{-}, g_{0}^{+}\right\}$, there exists a constant $C$ such that, for all $\xi \notin \pi_{B}^{-1} U$ and $g \in \tilde{\Gamma}_{s}$, we have $\left|h_{\xi}\left(g^{-1}\right)-d(e, g)\right| \leqslant C$. For $\xi \in \pi_{B}^{-1} U$, we only use the trivial bound $h_{\xi}\left(g^{-1}\right) \geqslant-d(e, g)$, since horofunctions are 1-Lipschitz and vanish at the origin. We get

$$
\begin{aligned}
\ell\left(\lambda_{\varepsilon}\right) & \geqslant \int_{(L, \xi) \in \Gamma \times \pi_{B}^{-1} U^{c}} d(e, L) \mathrm{d} \tilde{\alpha}_{\varepsilon}(L) \mathrm{d} \tilde{\rho}_{\varepsilon}(\xi)-\int_{(L, \xi) \in \Gamma \times \pi_{B}^{-1} U} d(e, L) \mathrm{d} \tilde{\alpha}_{\varepsilon}(L) \mathrm{d} \tilde{\rho}_{\varepsilon}(\xi)-C \\
& =\left(\int d(e, L) \mathrm{d} \tilde{\alpha}_{\varepsilon}(L)\right)\left(\tilde{\rho}_{\varepsilon}\left(\pi_{B}^{-1} U^{c}\right)-\tilde{\rho}_{\varepsilon}\left(\pi_{B}^{-1} U\right)\right)-C .
\end{aligned}
$$

For small enough $\varepsilon$, we have $\tilde{\rho}_{\varepsilon}\left(\pi_{B}^{-1} U\right) \leqslant 1 / 4$ (and therefore $\tilde{\rho}_{\varepsilon}\left(\pi_{B}^{-1} U^{c}\right) \geqslant 3 / 4$ ). Moreover, Lemma 5.6 ensures that the average distance to the origin for the measure $\tilde{\alpha}_{\varepsilon}$ is at least $c \cdot(\eta / \varepsilon)^{1 / 2}$. Hence, the previous formula completes the proof.

Combining Lemmas 5.5 and 5.8 , we get

$$
h\left(\lambda_{\varepsilon}\right) / \ell\left(\lambda_{\varepsilon}\right) \leqslant C \log (\eta / \varepsilon) /(\eta / \varepsilon)^{1 / 2} .
$$

This tends to 0 since $\eta / \varepsilon$ tends to infinity. We deduce from Lemma 5.1 that $h\left(\mu_{\varepsilon}\right) / \ell\left(\mu_{\varepsilon}\right)$ tends to 0 . This is a contradiction since we were assuming that it converges to the maximum $M$, which is positive.

This concludes the proof of Theorem 1.7.
The study of the case where $\Gamma_{s}$ is virtually cyclic infinite gives in particular the following result.

Theorem 5.9. Let $(\Gamma, d)$ be a metric hyperbolic group. Let $\Sigma$ be a finite subset of $\Gamma$ which generates a non-elementary group. Let $\mu_{i}$ be a sequence of measures on $\Sigma$, with $h\left(\mu_{i}\right)>$ 0 , converging to a probability measure $\mu$ such that $\Gamma_{\mu}$ is infinite virtually cyclic. Then $h\left(\mu_{i}\right) / \ell\left(\mu_{i}\right) \rightarrow 0$.

Note that the precise value of $\ell\left(\mu_{i}\right)$ depends on the choice of the distance, but if two distances are equivalent then the associated drifts vary within the same constants. Hence, the convergence $h\left(\mu_{i}\right) / \ell\left(\mu_{i}\right) \rightarrow 0$ does not depend on the distance.

We recover results of Le Prince [LP07]: In any metric hyperbolic group, there exist admissible probability measures with $h / \ell<v$. The construction of Le Prince is rather similar to the examples given by Theorem 5.9.

Example 5.10. We can use the above proof to also find an example where $h\left(\mu_{\varepsilon}\right) / \ell\left(\mu_{\varepsilon}\right) \rightarrow 0$ although $\mu_{\varepsilon}$ tends to a measure $\mu$ for which $\Gamma_{\mu}$ is finite and nontrivial. Consider $\Gamma=$ $\mathbb{Z} / 2 \times \mathbb{F}_{2}=\{0,1\} \times\langle a, b\rangle$, endowed with the probability measure $\mu_{\varepsilon}$ given by

$$
\mu_{\varepsilon}(0, e)=\mu_{\varepsilon}(1, e)=1 / 2-\varepsilon-\varepsilon^{2}, \quad \mu_{\varepsilon}(0, a)=\mu_{\varepsilon}\left(0, a^{-1}\right)=\varepsilon, \quad \mu_{\varepsilon}(0, b)=\mu_{\varepsilon}\left(0, b^{-1}\right)=\varepsilon^{2} .
$$

The measure $\mu_{\varepsilon}$ converges to $\mu=\left(\delta_{(0, e)}+\delta_{(1, e)}\right) / 2$. With the above notations, $\Gamma_{\mu}=\mathbb{Z} / 2 \times\{e\}$ but $\Gamma_{s}=\mathbb{Z} / 2 \times\langle a\rangle$ is virtually cyclic infinite (so that $h\left(\mu_{\varepsilon}\right) / \ell\left(\mu_{\varepsilon}\right) \rightarrow 0$ ) and $\Gamma_{r}=\Gamma$.

## 6. Examples for non-symmetric measures

In this section, we describe the additional difficulties that arise if one tries to prove Theorem 1.3 for non-symmetric measures. The main problem is that the random walk lives on the subsemigroup $\Gamma_{\mu}^{+}$, which is not a subgroup any more. While many cases can be handled with the tools we have described in this article, one case can not be treated in this way: when the subsemigroup $\Gamma_{\mu}^{+}$has no nice geometric properties (it is not quasi-convex, it is not a subgroup), but $\Gamma_{\mu}=\Gamma$.

Let us first show that the growth properties of such a subsemigroup can be more complicated than what happens for subgroups. If $\Lambda$ is a subgroup of $\Gamma$, either $\left|B_{n} \cap \Lambda\right| \asymp e^{n v}$, or $\left|B_{n} \cap \Lambda\right|=o\left(e^{n v}\right)$ (the first case happens if and only if $\Lambda$ has finite index in $\Gamma$, see the discussion at the beginning of Paragraph 4.3). Unfortunately, the behavior of semigroups can be more complicated.

Proposition 6.1. In $\mathbb{F}_{2}$, there exists a subsemigroup $\Lambda^{+}$such that $\liminf \left|B_{n} \cap \Lambda^{+}\right| /\left|B_{n}\right|=0$ and $\lim \sup \left|B_{n} \cap \Lambda^{+}\right| /\left|B_{n}\right|>0$.

Proof. Let $\mathbb{S}_{a, a}^{n}$ denote the geodesic words in $\mathbb{F}_{2}=\langle a, b\rangle$ of length $n$ which start and end with $a$. Let $n_{j}$ be a sequence tending very quickly to infinity. Let $\Lambda^{+}$be the subsemigroup generated by $\bigcup \mathbb{S}_{a, a}^{n_{j}}$. Then $\left|B_{n_{j}} \cap \Lambda^{+}\right| \geqslant c\left|B_{n_{j}}\right|$. We claim that

$$
\left|B_{n_{j}-1} \cap \Lambda^{+}\right| /\left|B_{n_{j}-1}\right| \rightarrow 0
$$

Indeed, the subsemigroup $\Lambda_{j-1}^{+}$generated by $\bigcup_{k<j} \mathbb{S}_{a, a}^{n_{k}}$ has a growth rate which is $<e^{n v}$, since some subwords such as $b^{n_{j-1}}$ are forbidden in this subsemigroup. Hence, if $n_{j}$ is large enough with respect to $n_{j-1}$, we have $\left|\mathbb{S}^{n_{j}-1} \cap \Lambda^{+}\right|=\left|\mathbb{S}^{n_{j}-1} \cap \Lambda_{j-1}^{+}\right|=o\left(e^{\left(n_{j}-1\right) v}\right)$.

In this example, most points in $\mathbb{S}^{n_{j}} \cap \Lambda^{+}$are introduced by $\mathbb{S}_{a, a}^{n_{j}}$. This shows that $\Lambda^{+}$is far from being quasi-convex. In particular, techniques based only on non-quasi-convexity and sub- or super-multiplicativity will never show that $\left|B_{n} \cap \Lambda^{+}\right|=o\left(\left|B_{n}\right|\right)$ for subsemigroups.

Now, we give an example of a well-behaved measure (apart from the fact that it is not symmetric, not admissible and not finitely supported) for which $h=\ell v$. The construction is done in free products. The idea is to forbid simplifications, so that we have an explicit control on the random walk at time $n$. To enforce this behavior, we will work in a free product $\Gamma_{1} * \Gamma_{2}$, and consider a probability measure supported on elements of the form $g_{1} g_{2}$ with $g_{i} \in \Gamma_{i} \backslash\{e\}$. The next statement applies to some non virtually free hyperbolic groups, for instance the free product of two surface groups. It also applies to some non-hyperbolic groups, more precisely to all finitely generated groups without torsion and with infinitely many ends, by Stallings' theorem. It would be of interest to extend it to all groups with infinitely many ends. For this, we would need to also handle amalgamated free products and HNN extensions.

Proposition 6.2. Let $\Gamma_{1}$ and $\Gamma_{2}$ be two nontrivial groups, generated respectively by finite symmetric sets $S_{1}$ and $S_{2}$. Let $\Gamma=\Gamma_{1} * \Gamma_{2}$ with the generating set $S=S_{1} \cup S_{2}$ and the corresponding word distance. There exists on $\Gamma$ a (nonsymmetric, nonadmissible) probability measure $\mu$, with an exponential moment and nonzero entropy, satisfying $h(\mu)=\ell(\mu) v$.

Proof. For $i=1,2$, let $\Gamma_{i}^{*}=\Gamma_{i} \backslash\{e\}$. We claim that

$$
\begin{equation*}
\sum_{g_{1} \in \Gamma_{1}^{*}, g_{2} \in \Gamma_{2}^{*}} e^{-v\left|g_{1} g_{2}\right|}=1 \tag{6.1}
\end{equation*}
$$

where $v$ is the growth rate of $\Gamma$.
Let $F_{i}(z)$ be the growth series of $\Gamma_{i}$, i.e., $F_{i}(z)=\sum_{g \in \Gamma_{i}} z^{|g|}$. The spheres $\mathbb{S}_{i}^{n} \in \Gamma_{i}$ satisfy $\mathbb{S}_{i}^{n+m} \subset \mathbb{S}_{i}^{n} \cdot \mathbb{S}_{i}^{m}$. Hence, the sequence $\log \left|\mathbb{S}_{i}^{n}\right|$ is subadditive. This implies that $\log \left|\mathbb{S}_{i}^{n}\right| / n$ converges to its infimum $v_{i}$, and moreover that $\left|\mathbb{S}_{i}^{n}\right| \geqslant e^{n v_{i}}$. We deduce that the radius of convergence of $F_{i}$ is $e^{-v_{i}}$, and moreover $F_{i}\left(e^{-v_{i}}\right)=+\infty$.

Let $F(z)$ be the growth series of $\Gamma$. As in the proof of Proposition 4.8, it is given by

$$
F(z)=\frac{F_{1}(z) F_{2}(z)}{1-\left(F_{1}(z)-1\right)\left(F_{2}(z)-1\right)}
$$

Assume for instance $v_{1} \geqslant v_{2}$. As $F_{1}\left(e^{-v_{1}}\right)=+\infty$, the function $\left(F_{1}(z)-1\right)\left(F_{2}(z)-1\right)$ takes the value 1 when $z$ increases to $e^{-v_{1}}$, at a point which is precisely the radius of convergence $e^{-v}$ of $F$. This shows that $\left(F_{1}\left(e^{-v}\right)-1\right)\left(F_{2}\left(e^{-v}\right)-1\right)=1$. This is precisely the equality (6.1).

We define a probability measure $\mu$ on $\Gamma$ as follows: for $\left(g_{1}, g_{2}\right) \in \Gamma_{1}^{*} \times \Gamma_{2}^{*}$, let

$$
\mu\left(g_{1} g_{2}\right)=e^{-v\left|g_{1} g_{2}\right|}
$$

Since there is only one way to generate the word $g_{1}^{1} g_{2}^{1} \cdots g_{1}^{n} g_{2}^{n}$ using $\mu$, we have

$$
\mu^{* n}\left(g_{1}^{1} g_{2}^{1} \cdots g_{1}^{n} g_{2}^{n}\right)=e^{-v \sum_{i}\left|g_{1}^{i} g_{2}^{i}\right|}
$$

Denoting by $X_{n}$ the position of the random walk at time $n$, it follows that $-\log \mu^{* n}\left(X_{n}\right)=$ $v\left|X_{n}\right|$. Dividing by $n$ and letting $n$ tend to infinity, this gives $h(\mu)=\ell(\mu) v$.

If one is interested in measures with finite support, one can only get the following approximation result. It has the same flavor as Theorem 1.4, but it is both stronger since it also applies to some non-hyperbolic groups, and weaker since the measures it produces are not admissible nor symmetric.

Proposition 6.3. Let $\Gamma_{1}$ and $\Gamma_{2}$ be two nontrivial groups, generated respectively by finite symmetric sets $S_{1}$ and $S_{2}$. Let $\Gamma=\Gamma_{1} * \Gamma_{2}$ with the generating set $S=S_{1} \cup S_{2}$ and the corresponding word distance. Then

$$
\sup \{h(\mu) / \ell(\mu): \mu \text { finitely supported probability measure in } \Gamma, \ell(\mu)>0\}=v
$$

Proof. Any element in $\Gamma$ can be canonically decomposed as a word in elements of $\Gamma_{1}$ and $\Gamma_{2}$. Let $\mathbb{S}_{i, j}^{p}$ be the set of elements of length $p$ that start with an element in $\Gamma_{i}$ and end with an element in $\Gamma_{j}$. We have the decomposition

$$
\mathbb{S}^{p}=\mathbb{S}_{1,1}^{p} \cup \mathbb{S}_{1,2}^{p} \cup \mathbb{S}_{2,1}^{p} \cup \mathbb{S}_{2,2}^{p}
$$

One term in this decomposition has cardinality at least $\left|\mathbb{S}^{p}\right| / 4$. Hence, there exist $i, j$ such that $\lim \sup \log \left|\mathbb{S}_{i, j}^{p}\right| / p=v$. Multiplying by fixed elements at the beginning and at the end to go from $\Gamma_{1}$ to $\Gamma_{i}$, and from $\Gamma_{j}$ to $\Gamma_{2}$, we get

$$
\begin{equation*}
\limsup \log \left|\mathbb{S}_{1,2}^{p}\right| / p=v \tag{6.2}
\end{equation*}
$$

Let $\mu_{p}$ be the uniform probability measure on $\mathbb{S}_{1,2}^{p}$. By construction, there are no simplifications when one iterates $\mu_{p}$. Hence, $\mu_{p}^{* n}$ is the uniform probability measure on $\left(\mathbb{S}_{1,2}^{p}\right)^{* n}$, whose cardinality is $\left|\mathbb{S}_{1,2}^{p}\right|^{n}$. We get $H\left(\mu_{p}^{* n}\right)=n \log \left|\mathbb{S}_{1,2}^{p}\right|$ and $L\left(\mu_{p}^{* n}\right)=n p$. Therefore, $h\left(\mu_{p}\right)=\log \left|\mathbb{S}_{1,2}^{p}\right|$ and $\ell\left(\mu_{p}\right)=p$, giving

$$
h\left(\mu_{p}\right) / \ell\left(\mu_{p}\right)=\log \left|\mathbb{S}_{1,2}^{p}\right| / p
$$

Together with (6.2), this proves the proposition.

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