TRANSLATES OF S-ARITHMETIC PERIODIC ORBITS AND APPLICATIONS

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ABSTRACT. We prove that certain sequences of periodic orbits of the diagonal group in the space of lattices equidistribute. As an application we obtain new information regarding the sequence of best approximations to certain vectors with algebraic coordinates. In order to prove these results we generalize the seminal work of Eskin Mozes and Shah about the equidistribution of translates of periodic measures from the real case to the S-arithmetic case.

1. Introduction

1.1. Context. Let G be a locally compact second countable topological group and $\Gamma < G$ a closed subgroup (usually assumed to be discrete or even a lattice in G). The space $X = \Gamma \backslash G$ endowed with the quotient topology is then called a homogeneous space. The group G and its various subgroups act on X from the left by $h \cdot \Gamma g = \Gamma g h^{-1}$. In homogeneous dynamics one studies these actions and tries to understand orbits. Given a closed subgroup H < G the most fundamental question is to understand topological properties of H-orbits: Which orbits are closed, dense or bounded? What can be said about orbit closures in general? A fundamental role in the discussion is played by periodic probability measures. Let $\mathcal{P}(X)$ denote the space of Borel probability measures on X equipped with the weak* topology and note that the left action of G on X induces a continuous left action on $\mathcal{P}(X)$ by $(g,\mu) \mapsto g\mu$ where for $f \in C_c(X)$, $\int_X f(x)d(g\mu) = \int_X f(g^{-1}x)d\mu$. An orbit $Hx \subset X$ is called periodic if there exists $\mu \in \mathcal{P}(X)$ which is H-invariant (i.e. $h\mu = \mu$ for all $h \in H$) and such that $\mu(Hx) = 1$. The measure μ is unique in this case and will be referred to as the algebraic or periodic measure on the orbit. It is denoted by μ_{Hx} or $\mu_{\Gamma \setminus \Gamma gH}$ if $x = \Gamma g$.

A significant portion of the literature in homogeneous dynamics is devoted to the following:

Problem 1.1. Understand the weak*-accumulation points (limits) of periodic measures. That is, what can be said about possible limits of sequences $\mu_{H_n x_n}$ where $H_n x_n$ are periodic orbits.

Common phenomena related to Problem 1.1 in particular cases are:

- (1) There are limit measures in $\mathcal{P}(X)$ this will be referred to as non-escape of mass.
- (2) All limit measures are periodic themselves this will be referred to as rigidity.

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(3) Often, when non-escape of mass and rigidity occurs, if one imposes some mild conditions on the sequence H_nx_n (usually of algebraic nature), then all limit measures must be G-invariant. Since there is a unique G-invariant measure on X, this actually implies that the sequence converges to it – a phenomenon which we refer to as equidistribution.

In the general setting described above Problem 1.1 might be too hard to tackle but under various technical assumptions on the players (G, Γ, H_n, x_n) quite a bit is known. A standing assumption will be that $\Gamma < G$ is a lattice. We mention two seminal works in this direction:

- In [MS95] Mozes and Shah solve this problem when G, H_n are real Lie groups and the periodic measure $\mu_{H_n x_n}$ is ergodic under the action of a one-parameter unipotent subgroup of H_n . It is paramount for their consideration that H_n will contain unipotent elements and thus when H_n is a reductive group and in particular, a maximal torus their results do not apply.
- When G, H are reductive algebraic groups defined over \mathbb{Q} , $\Gamma < G$ is an arithmetic lattice in G and the centralizer of H does not have any non-trivial \mathbb{Q} -characters, then in the seminal papers [EMS96, EMS97] Eskin, Mozes, and Shah analyze the possible limit points of μ_{g_nHx} of deformations of a fixed periodic H-orbit (note that $g_nHx = g_nHg_n^{-1}g_nx$ so this falls into the above setting).

In applications, it is often required to extend the above results to the S-arithmetic setting. The situation in the literature regarding [MS95] is satisfactory. In [GO11] Gorodnik and Oh generalize [MS95] to the S-arithmetic setting. In this paper, we are motivated by an application which require an S-arithmetic version of [EMS96, EMS97]. As far as we are aware of, there is no published paper in the literature establishing such a result and so in order to pursue our application we take upon the endeavour to produce such a proof in a setting general enough for our application. We note that there is a paper on arxiv [RZ16] in which an S-arithmetic generalization of [EMS96, EMS97] is given. The formalism there is quite sophisticated and we were not able to understand how to apply the results from that paper nor to verify its correctness. Eventually we preferred to give a self contained treatment even on the expense of compromising brevity.

1.2. **Motivating example.** Let $X = \Gamma \backslash G$ be the space of homothety classes of lattices in \mathbb{R}^d . Here $G = \mathbf{PGL}_d(\mathbb{R})$, $\Gamma = \mathbf{PGL}_d(\mathbb{Z})$. An element $g \in G$ will be represented by a matrix in $\mathbf{GL}_d(\mathbb{R})$ which by abuse of notation we also denote by g and the coset Γg represents the homothety class of the lattice spanned by the rows of g, $g^t\mathbb{Z}^d$. Let G denote the (connected component of the identity of the) group of diagonal matrices.

Let \mathbb{K}/\mathbb{Q} be a totally real number field of degree d. Let $\sigma_j : \mathbb{K} \to \mathbb{R}$ (j = 1, ..., d) be the distinct embeddings of \mathbb{K} into \mathbb{R} ordered in some way. Let $\varphi : \mathbb{K} \to \mathbb{R}^d$ be the map $\varphi = (\sigma_1, ..., \sigma_d)$. For

¹The alert reader will notice that this identification is not equivariant with respect to the left G-action on X and the natural left action on the space of lattices. Under this identification if a coset $x = \Gamma g$ represents a lattice $\Lambda \subset \mathbb{R}^d$, then gx represents the lattice $g^*\Lambda$, where $g^* = (g^{-1})^t$.

an ordered basis $\vec{\alpha} = (\alpha_1, \dots, \alpha_d)$ of \mathbb{K}/\mathbb{Q} let

$$g_{\vec{\alpha}} = \begin{pmatrix} -\varphi(\alpha_1) - \\ \vdots \\ -\varphi(\alpha_d) - \end{pmatrix} \tag{1}$$

and let $x_{\vec{\alpha}} = \Gamma g_{\vec{\alpha}}$. Then, it is well known (see e.g. [LW01]) that the orbit $Ax_{\vec{\alpha}}$ is periodic. The motivating problem which started this paper was to establish the following result which is proved in §7.

Theorem 1.2. Let \mathbb{K}/\mathbb{Q} be a degree d totally real number field with basis $\vec{\alpha} = (\alpha_1, \dots, \alpha_d)$. For any $m \in \mathbb{Z}$ $(m \neq 0, \pm 1)$ and for any sequence $\mathbf{i}_n = (i_{n,1}, \dots, i_{n,d}) \in \mathbb{Z}^d$ consider the sequence $\vec{\alpha}_n \stackrel{\text{def}}{=} (m^{i_{n,1}}\alpha_1, m^{i_{n,2}}\alpha_2, \dots, m^{i_{n,d}}\alpha_d)$. Then, for any choice of elements $s_n \in G$ the sequence of periodic orbits $s_n A x_{\vec{\alpha}_n} = \Gamma g_{\vec{\alpha}_n} A s_n^{-1}$ equidistribute in X as long for any $1 \leq j < r \leq d$, $\lim_n |i_{n,j} - i_{n,r}| = \infty$.

Remark 1.3. Theorem 1.2 attains an appealing form when one chooses the basis $1, \beta, \ldots, \beta^{d-1}$ for $\beta \in \mathbb{K}$ which generates the extension. In this case we write x_{β} for the corresponding lattice. Given a generator $\alpha \in \mathbb{K}$, we may choose in Theorem 1.2 the sequence $\mathbf{i}_n = (0, n, 2n, \ldots, (d-1)n)$ and obtain that the sequence of periodic orbits $Ax_{m^n\alpha}$ equidistribute in X.

Theorem 1.2 falls into the setting of Problem 1.1. The sequence of orbits under consideration is a deformation of a fixed orbit but this deformation is done in an S-arithmetic extension of X and then projected back to X. This is why the results of [EMS96, EMS97] do not apply and one needs an S-arithmetic version of them. In the case d=2 (with $s_n=e$) this equidistribution was established in [AS18] using different techniques. There the authors used mixing (and hence obtain effective equidistribution). For $d \geq 3$ the strategy of [AS18] fails and one needs to appeal to the techniques of [EMS96, EMS97] as we do here.

In [AS18] this equidistribution result was applied to deduce statistical information about the periods in the continued fraction expansions of quadratic numbers of the form $m^n \alpha$ as $n \to \infty$. In §1.3 we apply Theorem 1.2 to obtain new results regarding the statistics of best approximations of vectors like $(m\alpha, m^2\alpha^2, \ldots, m^{d-1}\alpha^{d-1})$ where $\mathbb{Q}(\alpha)$ is a totally real number field of degree d over \mathbb{Q} .

1.3. An application: Best approximations of algebraic vectors. In this subsection we wish to demonstrate the significance of Theorem 1.2 to the theory of Diophantine approximation. Readers whose interest lies solely in the dynamical aspects can skip this subsection altogether. We briefly recall some basic concepts from the theory of Diophantine approximation. Given an irrational vector $v \in \mathbb{R}^{d-1}$ there is a natural correspondence between rational vectors $\frac{1}{q}\mathbf{p}$ (here $q \in \mathbb{Z}, \mathbf{p} \in \mathbb{Z}^{d-1}$ and $\gcd(\mathbf{p}, q) = 1$) close to v and primitive integral vectors $\binom{\mathbf{p}}{q} \in \mathbb{Z}^d$ approximating the line $\mathbb{R} \cdot \binom{v}{1}$. One defines the sequence of best approximations to v to be a sequence of primitive integral vectors $\binom{\mathbf{p}_k}{q_k} \in \mathbb{Z}^d$ defined recursively as follows: $q_1 = 1$ and \mathbf{p}_1 is (one of the potentially finitely many) integral vectors satisfying $||\mathbf{p}_1 - v|| = \min_{\mathbf{p} \in \mathbb{Z}^{d-1}} ||\mathbf{p} - v||$. For k > 1, q_k is the minimal integer q

for which

$$\min \left\{ ||qv - \mathbf{p}|| : \mathbf{p} \in \mathbb{Z}^{d-1} \right\} < ||q_{k-1}v - \mathbf{p}_{k-1}||$$

and \mathbf{p}_k is then chosen so that

$$||q_k v - \mathbf{p}_k|| = \min \left\{ ||q_k v - \mathbf{p}|| : \mathbf{p} \in \mathbb{Z}^{d-1} \right\}.$$

A more geometric way to define (and understand) the vectors $\binom{\mathbf{p}_k}{q_k}$ is as follows: One travels along the ray $\{t \cdot \binom{v}{1}: t \geq 1\}$ and records the closest integral vectors one sees along the way; that is, one records an integral vector at a certain (integral) time if its distance to the ray beats the distances recorded thus far. In this description the distance one uses on \mathbb{R}^d is not the euclidean one but $d(\binom{w_1}{r_1},\binom{w_2}{r_2}) = \max\{||w_1 - w_2||,|r_1 - r_2|\}$ where $w_i \in \mathbb{R}^{d-1}$ and $r_i \in \mathbb{R}$ and the norm on \mathbb{R}^{d-1} is the euclidean norm (in fact, one can choose any norm on \mathbb{R}^{d-1} which is of interest). This way, each irrational vector $v \in \mathbb{R}^{d-1}$ yields a sequence of primitive integral vectors $\binom{\mathbf{p}_k}{q_k}$ which is referred to as the sequence of best approximations to the ray $\mathbb{R} \cdot \binom{v}{1}$ or by slight abuse of language, to v.

The sequence of best approximations is a central object of study in Diophantine approximation. Recently, in [SW], results pertaining to statistical information regarding these sequences in two cases were established: For Lebesgue almost any v some information about the behaviour of $\binom{\mathbf{p}_k}{q_k}$ was established and it was shown that they obey certain universal laws. In addition, for certain vectors v with algebraic coordinates it was shown that the sequence $\binom{\mathbf{p}_k}{q_k}$ obeys statistical laws which are different from the universal ones. As an application of Theorem 1.2 we show that if $\mathbb{Q}(\alpha)$ is a totally real field of degree d over \mathbb{Q} and if $m \neq 0, \pm 1$ is an integer then the statistical laws corresponding to $v_n = (m^n \alpha, m^{2n} \alpha^2, \dots, m^{(d-1)n} \alpha^{d-1})$ approach the universal statistical law as $n \to \infty$. To be more precise we need to define the objects pertaining the random variables which obey these statistical laws.

To each best approximation vector $\begin{pmatrix} \mathbf{p}_k \\ q_k \end{pmatrix}$ we associate a triple

$$(\Lambda_k, w_k, (\mathbf{p}_k)) \in X_{d-1} \times \mathbb{R}^{d-1} \times \hat{\mathbb{Z}}^d$$

where X_{d-1} is the space of lattices of covolume 1 in \mathbb{R}^{d-1} , and $\hat{\mathbb{Z}} = \prod_p \mathbb{Z}_p$ is the pro-finite completion of \mathbb{Z} . This correspondence is defined in the following way: The lattice Λ_k is referred to as the directional lattice of \mathbb{Z}^d in direction of the k'th best approximation $\begin{pmatrix} \mathbf{p}_k \\ q_k \end{pmatrix}$. It is defined by

$$\Lambda_k \stackrel{\text{def}}{=} q_k^{1/(d-1)} \pi_{\mathbb{R}^{d-1}}^{\begin{pmatrix} \mathbf{p}_k \\ q_k \end{pmatrix}} (\mathbb{Z}^n),$$

where $\pi_{\mathbb{R}^{d-1}}^w$ is the linear projection of \mathbb{R}^d onto the hyperplane $\{\binom{u}{0}: u \in \mathbb{R}^{d-1}\}$ with kernel equal to $\mathbb{R} \cdot w$. Thus, Λ_k is the unimodular lattice of the horizontal hyperplane $\mathbb{R}^{d-1} \times \{0\}$ which records how \mathbb{Z}^d looks like from the direction of the k'th best approximation to v.

The vector $w_k \in \mathbb{R}^{d-1}$ is referred to as the displacement vector of the k'th best approximation $\begin{pmatrix} \mathbf{p}_k \\ q_k \end{pmatrix}$. It is defined by

$$w_k \stackrel{\text{def}}{=} q_k^{-1/(d-1)} (q_k v - \mathbf{p}_k).$$

This vector tells us two things: First, its direction $w_k/||w_k||$ on the unit sphere \mathbb{S}^{d-2} is the direction from which the k'th best approximation approaches the line $\mathbb{R} \cdot \binom{v}{1}$ and second, its length $||w_k||$ captures (the appropriate scaling of the) quality of approximation. It turns out that the normalization factor of $q_k^{-1/(d-1)}$ is the correct one for an interesting limit law to occur.

Finally, in the pro-finite coordinate $\hat{\mathbb{Z}}^d$ we simply take the integral vector $\begin{pmatrix} \mathbf{p}_k \\ q_k \end{pmatrix}$ as it is (embedded diagonally in the product). This coordinate captures congruence questions regarding the k'th best approximation.

The whereabouts of the triple $(\Lambda_k, w_k, \begin{pmatrix} \mathbf{p}_k \\ q_k \end{pmatrix})$ in the product space $X_{d-1} \times \mathbb{R}^{d-1} \times \hat{\mathbb{Z}}^d$ and the statistical properties of the sequence of triples as k changes is very interesting from the perspective of Diophantine approximation. In [SW] the following theorem is proved:

Theorem 1.4 ([SW]). Let $d \geq 2$ and let $||\cdot||$ be either the Euclidean norm or the sup norm on \mathbb{R}^{d-1} . For an irrational vector $v \in \mathbb{R}^{d-1}$ let $\binom{p_k(v)}{q_k(v)} = \binom{p_k}{q_k}$ be the sequence of best approximations of v. Then

(1) There exists a probability measure $\nu_{best} \in \mathcal{P}(X_{d-1} \times \mathbb{R}^{d-1} \times \hat{\mathbb{Z}}^d)$ such that for Lebesgue almost any v

$$\lim_{K} \frac{1}{K} \sum_{k=1}^{K} \delta_{(\Lambda_{k}, w_{k}, \begin{pmatrix} p_{k} \\ q_{k} \end{pmatrix})} = \nu_{best}.$$

(2) If $\vec{\alpha} = (\alpha_1, \dots, \alpha_{d-1}, 1)^t$ is a vector whose coordinates span a totally real number field of degree d over \mathbb{Q} then there exists a probability measure $\nu_{best}^{\vec{\alpha}} \in \mathcal{P}(X_{d-1} \times \mathbb{R}^{d-1} \times \hat{\mathbb{Z}}^d)$ such that

$$\lim_{K} \frac{1}{K} \sum_{k=1}^{K} \delta_{(\Lambda_{k}, w_{k}, \begin{pmatrix} p_{k} \\ q_{k} \end{pmatrix})} = \nu_{best}^{\vec{\alpha}}.$$

The following theorem will be proved in §7 as an application of Theorem 1.2 together with a result from [SW].

Theorem 1.5. Let $(\alpha_1, \ldots, \alpha_{d-1}, 1)^t$ be a vector whose coordinates span a totally real number field of degree d over \mathbb{Q} . Let $\mathbf{i}_n = (i_{n,1}, \ldots, i_{n,d-1}) \in \mathbb{Z}^{d-1}$ be such that for any $j \neq r \in \{1, \ldots, d-1\}$ both $i_{n,r} \to_n \infty$ and $|i_{n,r} - i_{n,j}| \to_n \infty$. Let $m \neq 0, \pm 1$ be an integer and and let

$$\vec{\alpha}_n \stackrel{\text{def}}{=} (m^{i_{n,1}}\alpha_1, m^{i_{n,2}}\alpha_2, \dots, m^{i_{n,d-1}}\alpha_{d-1}, 1).$$

Then, in the notation of Theorem 1.4, $\nu_{best}^{\vec{\alpha}_n} \xrightarrow{n} \nu_{best}$.

Remark 1.6. Quite a bit is known about the measures $\nu_{\mathbf{best}}$, $\nu_{\mathbf{best}}^{\vec{\alpha}}$ appearing in Theorem 1.4 and in particular, about how different from each other they are. For example, it is known that

- (i) The projection of $\nu_{\mathbf{best}}$ to X_{d-1} is (not equal to but) in the same measure class as the periodic $\mathbf{SL}_{d-1}(\mathbb{R})$ -invariant measure, while the projection of $\nu_{\mathbf{best}}^{\vec{\alpha}}$ is singular to it.
- (ii) The projection of $\nu_{\mathbf{best}}$ to \mathbb{R}^{d-1} is a compactly supported absolutely continuous measure with respect to Lebesgue measure and if the norm used to define the best approximations is the

- Euclidean one then it is $SO_{d-1}(\mathbb{R})$ -invariant. On the other hand $\nu_{\mathbf{best}}^{\vec{\alpha}}$ is singular to Lebesgue and is not $SO_{d-1}(\mathbb{R})$ -invariant.
- (iii) The projection of $\nu_{\mathbf{best}}$ to $\hat{\mathbb{Z}}^d$ is the normalized restriction of the Haar measure to the compact open set obtained as the closure of $\mathbb{Z}^d_{\text{prim}}$ embedded diagonally in $\hat{\mathbb{Z}}^d$, while the projection of $\nu_{\mathbf{best}}^{\vec{\alpha}}$ is singular to it.

These differences highlight the significance of Theorem 1.5.

1.4. Limiting distributions of periodic orbits. We now state our main results Theorems 1.7, 1.8 and fix the notation of the paper. Let **G** be a connected (with respect to the Zariski topology) algebraic group defined over \mathbb{Q} . Let S be a finite set of valuations on \mathbb{Q} which contains the Archimedean one, and denote by S_f the subset of non-Archimedean places in S. Let $\mathbb{Q}_S = \prod_{p \in S} \mathbb{Q}_p$ and $\mathcal{O}_S = \mathbb{Z}[1/p: p \in S_f]$. Denote by

$$\mathbf{G}(\mathbb{Q}_S) = \prod_{p \in S} \mathbf{G}(\mathbb{Q}_p), \quad \Gamma_S = \mathbf{G}(\mathcal{O}_S), \quad \Gamma_\infty = \mathbf{G}(\mathbb{Z}).$$

For an algebraic group \mathbf{F} , we write \mathbf{F}^0 for the connected component of \mathbf{F} in the Zariski topology, and write $\mathbf{F}(\mathbb{R})^0$ for the connected component of $\mathbf{F}(\mathbb{R})$ as a Lie group. In this paper, we prove the following theorem.

Theorem 1.7. Let G be a connected reductive algebraic group defined over \mathbb{Q} without nontrivial \mathbb{Q} -characters. Let H be a connected reductive \mathbb{Q} -subgroup of G without nontrivial \mathbb{Q} -characters, and suppose that H contains a maximal torus of G. Let H be a subgroup of finite index in $H(\mathbb{Q}_S)$ and $\mu_{\Gamma_S \backslash \Gamma_S H}$ the natural invariant probability measure supported on $\Gamma_S \backslash \Gamma_S H$. Let $\{g_i\}_{i \in \mathbb{N}}$ be a sequence in $G(\mathbb{Q}_S)$. Then the sequence $g_i^* \mu_{\Gamma_S \backslash \Gamma_S H}$ has a subsequence converging to an algebraic probability measure μ on $\Gamma_S \backslash G(\mathbb{Q}_S)$.

As a corollary of Theorem 1.7, we deduce the following theorem regarding the equidistribution of certain arithmetically related periodic orbits in a real homogeneous space, which we will need to prove Theorem 1.2.

Theorem 1.8. Let G be a connected \mathbb{R} -split reductive \mathbb{Q} -group and T a maximal \mathbb{R} -split torus defined over \mathbb{Q} in G which is \mathbb{Q} -anisotropic. Let S be a finite set of valuations on \mathbb{Q} containing the archimedean one, $\{g_i\}$ a sequence in $G(\mathbb{R})$ and $\{h_i\}$ a sequence in Γ_S . Then the set of weak* limits of the sequence $\mu_{\Gamma_\infty h_i \mathbf{T}(\mathbb{R})^0 g_i^{-1}}$ is not empty, and contains only algebraic probability measures on $\Gamma_\infty \backslash G(\mathbb{R})$.

Moreover, if for any non-central element $x \in \mathbf{T}(\mathbb{Q})$, the sequence $(g_i x g_i^{-1}, h_i x h_i^{-1})$ diverges in $\mathbf{G}(\mathbb{Q}_S)$, then $\mu_{\Gamma_\infty h_i \mathbf{T}(\mathbb{R})^0 g_i^{-1}}$ converges to some translate of $\mu_{\Gamma_\infty \setminus \mathbf{G}(\mathbb{R})^0}$.

Remark 1.9. In the context of Theorem 1.8, for any $h \in \mathbf{G}(\mathbb{Q})$ and any $g \in \mathbf{G}(\mathbb{R})$, $\Gamma_{\infty} h \mathbf{T}(\mathbb{R})^0 g^{-1}$ is compact in $\Gamma_{\infty} \backslash \mathbf{G}(\mathbb{R})$. The argument for the proof realizes these orbits as projections of certain deformations of a fixed periodic orbit in the S-arithmetic extension of $\Gamma_{\infty} \backslash \mathbf{G}(\mathbb{R})$. In order to apply Theorem 1.7 we need this new lifted orbit to be an orbit of a group containing a maximal torus.

For this to hold, we need \mathbf{T} to be \mathbb{R} -split and this is the reason for the extra conditions appearing in the statement of Theorem 1.8 on top of the ones appearing in the statement of Theorem 1.7.

- 1.5. **About the proof of the main Theorem 1.7.** We follow closely the structure of [EMS96, EMS97] and divide the argument into two parts:
 - (i) Establish non-escape of mass.
 - (ii) Establish algebraicity of limit measures.

We prove non-escape of mass in §3 and then establish the algebraicity of the limit measures in §5. In order to establish part (i) we go along the lines of [EMS97] but prefer to use the more robust formalism of (C, α) -good functions a la Klainbock and Margulis [KM98] which was developed in its S-arithmetic version in [KT07]. To establish part (ii) we follow [EMS96, GO11] and use linearization and Ratner's classification theorem of measures invariant and ergodic under one parameter unipotent flows. Unlike the real case in [EMS96], in which the key observation is that a limit measure is invariant under a one-parameter unipotent subgroup because it has an invariance group having a non-trivial nilpotent line in its Lie algebra, over the p-adics, this is not enough. This is due to the existence of compact unipotent groups. In §4 we deal with this issue and prove that the limit measures in our discussion are invariant under a full-one parameter unipotent flow where the challenge is to prove invariance under large unipotent elements rather than just small ones. This invariance is then used in §5 in conjuction with the S-arithmetic linearization technique and Ratner's measure classification theorem to establish (ii).

Remark 1.10. For more details about the S-arithmetic Ratner's theorem and linearization technique we refer the reader also to [MT94,Rat98,Tom00]. One of the possible extensions of the results in this paper is to establish a theorem about limiting distributions of deformations of adelic periodic orbits. For limiting distributions of adelic torus orbits, one can read e.g. [ELMV09, ELMV11, Kha19, Ven10, Zam10, DS18a, DS18b]. In general, a sequence of deformations of a fixed periodic adelic torus orbit does not necessarily equidistribute in the ambient space, and one of the obstructions is the escape of mass phenomenon which leads to the limiting measures being not probability. One can refer to [AS18, Example 2.11] for an example of this phenomenon for $\mathbf{G} = \mathbf{PGL}_2$ due to Ubis. It is interesting to know whether the escape of mass phenomenon disappears when \mathbf{G} is of higher rank.

2. Growth rates of functions

In this section, we study (C, α) -good functions in local fields. The notion of (C, α) -good functions is introduced in [KM98]. One may read [Kle10, EMS97] for more details. Here we use the definition in [KT07] for local fields. The goal in this section is Proposition 2.2, which shows that certain functions are (C, α) -good. At the end of this section, we prove a weak analogue of the intermediate value theorem for (C, α) -good functions in local fields (Proposition 2.3) which will be important in the discussion in §4.

Let $k = \mathbb{Q}_p$ $(p = \text{prime or } \infty)$ with the *p*-adic valuation $|\cdot|_p$. Let *F* be a finite extension of *k*. Then there exists a unique extension $|\cdot|_F$ of $|\cdot|_p$ on *F* such that for any $x \in F$

$$|x|_F = |N_{F/k}(x)|_p^{\frac{1}{[F:k]}}$$

where $N_{F/k}(x)$ is the norm of x. Equipped with $|\cdot|_F$, F is a complete field. Note that for any extensions $k \subset F \subset L$, we have $N_{F/k} \circ N_{L/F} = N_{L/k}$ and [L:F][F:k] = [L:k]. Hence one can take the limit over all the finite extensions of k and get a well-defined norm $|\cdot|_{\bar{k}}$ on the algebraic closure \bar{k} of k. In the following, we will use the notation $|\cdot|_F$ for the norm on any finite extension F of $k = \mathbb{Q}_p$ obtained in this manner. We will write $\mathcal{O}_F := \{x \in F : |x|_F \leq 1\}$.

Let F be a complete field with a norm $|\cdot|$, X a metric space, μ a Borel measure on X and U a subset of X. Recall [KT07] that a continuous function $f: U \to F$ is called (C, α) -good on U with respect to μ if for any open ball $B \subset U$ centered in supp μ , we have

$$\mu(\lbrace x \in B : |f(x)| < \epsilon \rbrace) \le C \left(\frac{\epsilon}{\|f\|_{\mu,B}}\right)^{\alpha} \mu(B)$$

for any $\epsilon > 0$. Here

$$||f||_{\mu,B} = \sup\{|f(x)| : x \in B \cap \sup(\mu)\}.$$

In the sequel, we will study the case when X is a product of local fields and the measure μ is the Haar measure on X.

Lemma 2.1 (C.f. [EMS97, Lemma 2.4]). Let F be a finite extension of \mathbb{Q}_p . Then there exists $\gamma \in \mathbb{N}$ such that the determinant of the matrix generated by $s_i^k e^{\lambda_j s_i}$ $(i = 1, 2, ..., n^2, j = 1, 2, ..., n, k = 0, 1, 2, ..., n - 1)$

$$\begin{vmatrix} e^{\lambda_1 s_1} & e^{\lambda_1 s_2} & \cdots & e^{\lambda_1 s_{n^2}} \\ s_1 e^{\lambda_1 s_1} & s_2 e^{\lambda_1 s_2} & \cdots & s_{n^2} e^{\lambda_1 s_{n^2}} \\ \vdots & \vdots & \cdots & \vdots \\ s_1^{n-1} e^{\lambda_1 s_1} & s_2^{n-1} e^{\lambda_1 s_2} & \cdots & s_{n^2}^{n-1} e^{\lambda_1 s_{n^2}} \\ e^{\lambda_2 s_1} & e^{\lambda_2 s_2} & \cdots & e^{\lambda_2 s_{n^2}} \\ s_1 e^{\lambda_2 s_1} & s_2 e^{\lambda_2 s_2} & \cdots & s_{n^2} e^{\lambda_2 s_{n^2}} \\ \vdots & \vdots & \vdots & \vdots \\ s_1^{n-1} e^{\lambda_n s_1} & s_2^{n-1} e^{\lambda_n s_2} & \cdots & s_{n^2}^{n-1} e^{\lambda_n s_{n^2}} \end{vmatrix} = \prod_{i < j} (\lambda_i - \lambda_j)^{m_{i,j}} \prod_{i < j} (s_i - s_j) \phi(\boldsymbol{\lambda}, \boldsymbol{s})$$

where $m_{i,j} \in \mathbb{N}$, $\phi(\boldsymbol{\lambda}, \boldsymbol{s})$ is an analytic function in $\boldsymbol{\lambda} = (\lambda_1, \lambda_2, \dots, \lambda_n) \in (p^{\gamma} \mathcal{O}_F)^n$ and $\boldsymbol{s} = (s_1, s_2, \dots, s_{n^2}) \in (p^{\gamma} \mathcal{O}_F)^{n^2}$, and for any $\boldsymbol{\lambda} \in (p^{\gamma} \mathcal{O}_F)^n$, we have $\phi(\boldsymbol{\lambda}, \boldsymbol{s}) \neq 0$.

Proof. This lemma is proved in [EMS97, Lemma 2.4] for $F = \mathbb{C}$ and \mathbb{R} . Here we give a proof for the p-adic case. Indeed, if we expand the determinant and the analytic function $\phi(\boldsymbol{\lambda}, \boldsymbol{s})$ as power series of $s_1, s_2, ..., s_{n^2}$ and $\lambda_1, \lambda_2, ..., \lambda_n$ over \mathbb{R} , the formula proved in [EMS97, Lemma 2.4] holds as an identity of power series with rational coefficients. By examining the ranges of convergence of power series $s_i^k e^{\lambda_j s_i}$ $(1 \le i \le n^2, 1 \le j \le n, 0 \le k \le n-1)$ in the p-adic case, one concludes

that the formula holds for $\lambda_1, \lambda_2, ..., \lambda_n$ and $s_1, s_2, ..., s_{n^2}$ in small neighborhoods of 0 in \mathcal{O}_F . This completes the proof of the lemma.

Proposition 2.2. Let F be a finite field extension of \mathbb{Q}_p and $n \in \mathbb{N}$.

(1) If p is Archimedean, then there exists $\delta_0 > 0$ such that for any $c_i \in \mathbb{C}$, any $\lambda_i \in \mathbb{C}$ $(1 \le i \le n)$ and any interval I in \mathbb{R} of length at most $\delta_0 > 0$, the function $f: I \to \mathbb{C}$ defined by

$$f(s) = \sum_{i=1}^{n} \sum_{l=0}^{n-1} c_{i,l} s^l e^{\lambda_i s}, \quad s \in \mathbb{Z}_p,$$

is (C, α) -good for some $C, \alpha > 0$ depending only on $\delta_0 > 0$ and n. The measure on I is chosen to be the Lebesgue measure.

(2) If p is non-Archimedean, then there exists $k \in \mathbb{N}$ such that for any $c_i \in F$ and any distinct $\lambda_i \in p^k \mathcal{O}_F$ $(1 \le i \le n)$, the function $f : \mathbb{Z}_p \to F$ defined by

$$f(s) = \sum_{i=1}^{n} \sum_{l=0}^{n-1} c_{i,l} s^l e^{\lambda_i s}, \quad s \in \mathbb{Z}_p,$$

is (C, α) -good for some $C, \alpha > 0$ depending only on F and n. Here we choose the measure μ on \mathbb{Z}_p to be the Haar measure on \mathbb{Q}_p with $\mu(\mathbb{Z}_p) = 1$.

Proof. The case for the Archimedean place $p = \infty$ is essentially proved in [EMS97, Corollary 2.10]. Here we give a proof for non-Archimedean places. Denote by

$$d(s_1, s_2, \dots, s_{n^2}) := \begin{bmatrix} e^{\lambda_1 s_1} & e^{\lambda_1 s_2} & \cdots & e^{\lambda_1 s_{n^2}} \\ s_1 e^{\lambda_1 s_1} & s_2 e^{\lambda_1 s_2} & \cdots & s_{n^2} e^{\lambda_1 s_{n^2}} \\ \vdots & \vdots & \cdots & \vdots \\ s_1^{n-1} e^{\lambda_1 s_1} & s_2^{n-1} e^{\lambda_1 s_2} & \cdots & s_{n^2}^{n-1} e^{\lambda_1 s_{n^2}} \\ e^{\lambda_2 s_1} & e^{\lambda_2 s_2} & \cdots & e^{\lambda_2 s_{n^2}} \\ s_1 e^{\lambda_2 s_1} & s_2 e^{\lambda_2 s_2} & \cdots & s_{n^2} e^{\lambda_2 s_{n^2}} \\ \vdots & \vdots & \vdots & \vdots \\ s_1^{n-1} e^{\lambda_n s_1} & s_2^{n-1} e^{\lambda_n s_2} & \cdots & s_{n^2}^{n-1} e^{\lambda_n s_{n^2}} \end{bmatrix}.$$

Let B be an open ball in \mathbb{Z}_p and define

$$U_{B,f,\epsilon} := \{ s \in B : |f(s)|_F < \epsilon \}.$$

We may assume $\mu(U_{B,f,\epsilon}) > 0$. By computing the volume of p-adic open balls, one can inductively find n^2 points $s_1, s_2, \ldots, s_{n^2} \in U_{B,f,\epsilon}$ such that for any $1 \le i \ne j \le n^2$

$$|s_i - s_j|_p \ge \frac{\mu(U_{B,f,\epsilon})}{n^2}.$$

Now by interpolation formula, we have

$$f(s) = \sum_{j=1}^{n^2} f(s_j) \frac{d(s_1, \dots, s_{j-1}, s, s_{j+1}, \dots, s_{n^2})}{d(s_1, s_2, \dots, s_{n^2})},$$

and by Lemma 2.1

$$f(s) = \sum_{j=1}^{n^2} f(s_j) \frac{\prod_{i \neq j} (s - s_i) \phi(\boldsymbol{\lambda}, s_1, \dots, s, \dots, s_{n^2})}{\prod_{i \neq j} (s_j - s_i) \phi(\boldsymbol{\lambda}, s_1, \dots, s_{n^2})}$$

for some analytic function ϕ . Here $\phi(\lambda, s) \neq 0$ for λ in a small neighborhood $(p^l \mathcal{O}_F)^n$ of $\mathbf{0}$. This implies that for any $s \in B$

$$|f(s)|_F \le \max_{1 \le j \le n^2} \left\{ \left| f(s_j) \frac{\prod_{i \ne j} (s - s_i) \phi(\boldsymbol{\lambda}, s_1, \dots, s_{n^2})}{\prod_{i \ne j} (s_j - s_i) \phi(\boldsymbol{\lambda}, s_1, \dots, s_{n^2})} \right|_F \right\}$$

$$\le C\epsilon \frac{\mu(B)^{n^2 - 1}}{(\mu(U_{B, f, \epsilon})/n^2)^{n^2 - 1}}$$

for some constant C > 0 depending only on n and ϕ . Hence

$$\mu(U_{B,f,\epsilon}) \le n^2 C^{\frac{1}{n^2-1}} \left(\frac{\epsilon}{\sup_{s \in B} |f(s)|_F} \right)^{\frac{1}{n^2-1}} \mu(B).$$

We finish the proof of the proposition by letting $\alpha = 1/(n^2 - 1)$ and replacing $n^2 C^{\frac{1}{n^2 - 1}}$ by C.

We conclude this section by proving the following proposition, which is a weak analogue of the intermediate value theorem in the p-adic case.

Proposition 2.3 (Weak intermediate value theorem). Let F be a finite extension of \mathbb{Q}_p with norm $|\cdot|_F$ and $f:\mathbb{Z}_p\to F$ be a continuous (C,α) -good function. Then there exists a constant C'>0 depending only on C, α and p such that for any $n\in\mathbb{N}$ we have

$$\sup_{s \in p^{n+1}\mathbb{Z}_p} |f(s)|_F \ge \sup_{s \in p^n\mathbb{Z}_p} |f(s)|_F / C'.$$

In particular, if f(0) = 0 and $|f(u)|_F \ge \rho$ for some $u \in \mathbb{Z}_p$ and $\rho > 0$, then one can find $v \in \mathbb{Z}_p$ such that $\rho/C' \le |f(v)|_F < \rho$.

Proof. For any $n \in \mathbb{N}$, we denote by $S_n := \sup_{s \in p^n \mathbb{Z}_p} |f(s)|_F$. By the (C, α) -good property of f, for any $\epsilon > 0$

$$\mu(\{s \in p^n \mathbb{Z}_p : |f(s)|_F < \epsilon\}) \le C \left(\frac{\epsilon}{S_n}\right)^{\alpha} \mu(p^n \mathbb{Z}_p).$$

Let

$$\epsilon_n = S_n \left(\frac{\mu(p^{n+1}\mathbb{Z}_p)}{2C\mu(p^n\mathbb{Z}_p)} \right)^{\frac{1}{\alpha}} = \frac{S_n}{(2Cp)^{\frac{1}{\alpha}}}$$

and we have

$$\mu(\lbrace s \in p^n \mathbb{Z}_p : |f(s)|_F < \epsilon_n \rbrace) \le \frac{1}{2} \mu(p^{n+1} \mathbb{Z}_p).$$

This implies that there exists $v \in p^{n+1}\mathbb{Z}_p$ such that

$$|f(v)|_F \ge \epsilon_n = \frac{S_n}{(2Cp)^{\frac{1}{\alpha}}}.$$

We finish the proof of the first statement by taking $C' = (2Cp)^{\frac{1}{\alpha}}$.

For the second statement, by continuity of f, one can find $n \in \mathbb{N}$ such that

$$\sup_{s \in p^n \mathbb{Z}_p} |f(s)|_p \ge \rho \text{ and } \sup_{s \in p^{n+1} \mathbb{Z}_p} |f(s)|_p < \rho.$$

Now we can apply the first statement and obtain

$$\sup_{s \in p^{n+1}\mathbb{Z}_p} |f(s)|_F \ge \sup_{s \in p^n\mathbb{Z}_p} |f(s)|_F / C' \ge \rho / C'.$$

This implies that there exists $v \in p^{n+1}\mathbb{Z}_p$ such that $\rho/C' \leq |f(v)|_F < \rho$.

As a corollary, we deduce the following result.

Corollary 2.4. Let F be a finite extension of \mathbb{Q}_p with norm $|\cdot|_F$. Let $C, \alpha > 0$ and $\{f_i\}_{i \in \mathbb{N}}$ a sequence of continuous (C, α) -good function from \mathbb{Z}_p to F. Suppose $\sup_{s \in \mathbb{Z}_p} |f_i(s)|_F \to \infty$ as $i \to \infty$. Then there is a sequence $\{s_i\}$ in \mathbb{Z}_p such that $s_i \to 0$ and $f_i(s_i) \to \infty$.

Proof. Let C' be the constant as in Proposition 2.3. Since $\lim_{i\to\infty} \sup_{s\in\mathbb{Z}_p} |f_i(s)|_F = \infty$, one can find a sequence $\{k_i\}$ in \mathbb{N} such that

$$k_i \to \infty \text{ and } \lim_{i \to \infty} \frac{1}{C'^{k_i}} \sup_{s \in \mathbb{Z}_p} |f_i(s)|_F = \infty.$$

By Proposition 2.3, for each $i \in \mathbb{N}$, we have

$$\sup_{s \in p^{k_i} \mathbb{Z}_p} |f(s)|_F \ge \sup_{s \in \mathbb{Z}_p} |f(s)|_F / C'^{k_i}.$$

Hence, for each $i \in \mathbb{N}$, there is $s_i \in p^{k_i}\mathbb{Z}_p$ such that

$$|f(s_i)|_F \ge \sup_{s \in \mathbb{Z}_p} |f(s)|_F / (2C'^{k_i}).$$

Then the sequence $\{s_i\}_{i\in\mathbb{N}}$ satisfies our requirement.

3. Nondivergence

In this section, we study the non-escape of mass of S-arithmetic periodic orbits. The arguments in this section will rely on the results in [KM98, KT07, TW03]. Since **G** is a reductive \mathbb{Q} -group without non-trivial \mathbb{Q} -characters, we can embed **G** into \mathbf{SL}_n for some $n \in \mathbb{N}$.

3.1. (C, α) -good maps and Nondivergence. Let X be a metric space with a locally finite measure μ and S a finite set of places of \mathbb{Q} . A map

$$\phi: (X, \mu) \to \mathbf{SL}_n(\mathbb{Q}_S), \quad x \in X \mapsto \phi(x) = (\phi_p(x))_{p \in S} \in \mathbf{SL}_n(\mathbb{Q}_S),$$

where $\phi_p(x) \in \mathbf{SL}_n(\mathbb{Q}_p)$ $(p \in S)$, is called (C, α) -good if there exist constants C and $\alpha > 0$ such that for every place $p \in S$, each entry of ϕ_p is a (C, α) -good function from X to \mathbb{Q}_p .

For any p in the finite set S of places of \mathbb{Q} , let $\text{Lie}(\mathbf{SL}_n(\mathbb{Q}_p))$ denote the Lie algebra of $\mathbf{SL}_n(\mathbb{Q}_p)$, and let \exp_p be the associated exponential map which is defined in a neighborhood of identity in $\text{Lie}(\mathbf{SL}_n(\mathbb{Q}_p))$. We write \exp_S for the product map

$$(\exp_p)_{p \in S} : \prod_{p \in S} \operatorname{Lie}(\mathbf{SL}_n(\mathbb{Q}_p)) \to \prod_{p \in S} \mathbf{SL}_n(\mathbb{Q}_p)$$

and \log_S for the inverse of \exp_S , whenever these maps are defined. For any set W in $\prod_{p \in S} \mathbf{SL}_n(\mathbb{Q}_p)$, we denote by $\log_S(W)$ the image of W under the map \log_S . We write $\mathrm{Lie}(F)$ for the Lie algebra of a (p-adic) Lie group F.

Let $p \in S$ and $\{X_1^p, X_2^p, ..., X_l^p\}$ a basis of the Lie algebra $\operatorname{Lie}(\mathbf{H}(\mathbb{Q}_p))$ in a small neighborhood of 0 where the exponential map \exp_p is defined. By Jordan decomposition, we can write $X_i^p = Y_i^p + Z_i^p$ where Y_i^p is semisimple, Z_i^p is nilpotent and Y_i^p and Z_i^p commute. Then by Proposition 2.2, the map $s \mapsto \exp_p(sX_i^p) = \exp_p(sY_i^p) \exp_p(sZ_i^p)$ ($s \in \mathbb{Z}_p$) is (C, α) -good for some C and $\alpha > 0$. Consequently, by [KT07, Corollary 3.3], the map

$$\Phi_p: (s_1, s_2, ..., s_l) \mapsto \exp_p(s_1 X_1^p) \exp_p(s_2 X_2^p) \cdots \exp_p(s_l X_l^p)$$

is (C, α) -good. Since the map Φ_p is a local diffeomorphism from a neighborhood of 0 in $\text{Lie}(\mathbf{H}(\mathbb{Q}_p))$ to $\mathbf{H}(\mathbb{Q}_p)$, the Lebesgue measure on a neighborhood of 0 in $\text{Lie}(\mathbf{H}(\mathbb{Q}_p))$ maps to a measure on $\mathbf{H}(\mathbb{Q}_p)$ which is absolutely continuous to the Haar measure on $\mathbf{H}(\mathbb{Q}_p)$.

Now we introduce the notation in [KT07]. Let \mathbb{Q}_S^n be the \mathbb{Q}_S -module $\prod_{p\in S} \mathbb{Q}_p^n = \{(x_p)_{p\in S} : x_p \in \mathbb{Q}_p^n\}$. Denote by $\mathfrak{M}(\mathbb{Q}_S, \mathcal{O}_S, n)$ the space of discrete \mathcal{O}_S -modules in \mathbb{Q}_S^n , which can be identified as

$$\{g^{-1}\Delta: g \in \mathbf{GL}_n(\mathbb{Q}_S), \Delta \text{ is a } \mathcal{O}_S\text{-submodule in } \mathcal{O}_S^n\}.$$

Note that any discrete \mathcal{O}_S -module in \mathbb{Q}^n_S is a free \mathcal{O}_S -module. We write $\Omega_{S,n}$ for the space of lattices in \mathbb{Q}^n_S defined by

$$\Omega_{S,n} := \left\{ g^{-1} \mathcal{O}_S^n : g \in \mathbf{GL}_n(\mathbb{Q}_S) \right\} \cong \mathbf{GL}_n(\mathcal{O}_S) \setminus \mathbf{GL}_n(\mathbb{Q}_S).$$

For any $v \in S$, the valuation $|\cdot|_v$ induces a norm $||\cdot||_v$ on \mathbb{Q}_v^n . The metric on \mathbb{Q}_S^n , which is induced by the metrics $||\cdot||_v$ $(v \in S)$ as a product metric, will be denoted by $||\cdot||$. The content c(x) of $x = (x_v)_{v \in S} \in \mathbb{Q}_S^n = \prod_{p \in S} \mathbb{Q}_p^n$ is defined to be the product of $||x_v||_v$ $(v \in S)$, and will be used in §3.2.

Let

$$\mathbf{GL}_n^1(\mathbb{Q}_S) = \{ x = (x_p)_{p \in S} \in \mathbf{GL}_n(\mathbb{Q}_S) : \prod_p |\det(x_p)|_p = 1 \}.$$

Define the space of unimodular lattices in \mathbb{Q}_S^n by

$$\Omega_{S,n}^1 := \{ \Delta \in \Omega_{S,n} : \operatorname{cov}(\Delta) = 1 \} \cong \mathbf{GL}_n^1(\mathcal{O}_S) \setminus \mathbf{GL}_n^1(\mathbb{Q}_S).$$

Here we refer to [KT07] for the definition of the covolume $cov(\Delta)$ of a discrete \mathcal{O}_S -module Δ in \mathbb{Q}^n_S .

As **G** is embedded into \mathbf{SL}_n for some $n \in \mathbb{N}$, we have

$$X_S = \mathbf{G}(\mathcal{O}_S) \backslash \mathbf{G}(\mathbb{Q}_S) \subset \mathbf{GL}_n^1(\mathcal{O}_S) \backslash \mathbf{GL}_n^1(\mathbb{Q}_S)$$

and X_S can be identified with a subspace of \mathcal{O}_S -modules of covolume one in \mathbb{Q}_S^n . Note that the embedding above is a proper map.

In the following, we prove the non-divergence of the orbits $\Gamma \backslash \Gamma Hg^{-1}$ for $g \in \mathbf{G}(\mathbb{Q}_S)$, where H is a subgroup of finite index in $\mathbf{H}(\mathbb{Q}_S)$. We take a neighborhood of identity $U^* = \prod_{p \in S} U_p$ in H $(U_p \subset \mathbf{H}(\mathbb{Q}_p))$ such that the exponential map is defined on U_p $(\forall p \in S)$ and U_p is a subgroup of $\mathbf{H}(\mathbb{Q}_p)$ $(\forall p \in S_f)$. For convenience, in the proof of Theorem 3.1, we denote by μ the Haar measure on $\mathbb{R}^l \times \prod_{p \in S_f} \mathbb{Q}_p^l$ $(l = \dim \mathbf{H})$. The proof is similar to that of [EMS97, Theorem 3.4].

Theorem 3.1. Let G be a reductive algebraic group defined over \mathbb{Q} , which does not admit nontrivial \mathbb{Q} -characters. Let G be a reductive G-subgroup of G admitting no nontrivial G-characters. Let G be a finite set of places of G containing the Archimedean place, G a subgroup of finite index in G and G in G and G in G in G in G in G in G such that for any G is G in G

$$\Gamma_S \backslash \Gamma_S U^* h^{-1} \cap K_0 \neq \emptyset.$$

Then for any $\epsilon > 0$, there exists a compact subset $K \subset X_S$ such that

$$\mu_{\Gamma_S \setminus \Gamma_S H g^{-1}}(K) \ge 1 - \epsilon.$$

Proof. Let Ω be a fundamental domain of $\Gamma_S \setminus \Gamma_S H$ in H, μ_H the Haar measure on H, and let δ_0 be as in Proposition 2.2. Let W be a compact subset in Ω such that

$$\mu_H(W)/\mu_H(\Omega) \ge 1 - \epsilon$$

We can construct a finite open cover \mathcal{U} of W in H, which consists of translates of U^* $\{B_i\}_{i=1}^N = \{U^*h_i\}_{i=1}^N$, and here it is harmless to assume that the diameter of U^* is equal to $\delta_0 > 0$. Note that the number N depends only on Ω and $\epsilon > 0$.

By the discussion in this section, for any B_i $(1 \le i \le N)$, there is an open ball U_i of radius at most $\delta_0 > 0$ in \mathbb{R}^l such that one can find a (C, α) -good map

$$\phi_i: U_i \times \prod_{p \in S_f} \mathbb{Z}_p^l \to B_i$$

where C and α are uniform for all $1 \leq i \leq N$. Note that for each i = 1, 2, ..., N, the map $\phi_i(s)^{-1}$ $(s \in U_i \times \prod_{p \in S_f} \mathbb{Z}_p^l)$, which is the composition of ϕ_i with the inversion in the group H, is also (C, α) -good.

Now we prove the non-escape of mass of $\mu_{\Gamma \backslash \Gamma Hg^{-1}}$. We know that there is a compact subset K_0 in $\Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$ such that $\Gamma_S \backslash \Gamma_S B_i g^{-1}$ intersects K_0 for every $1 \leq i \leq N$. By [KT07, Theorem 8.8, Lemma 9.2 and Lemma 8.6], this implies that there exists $\rho > 0$ depending only on K_0 , n and S

such that for any $1 \leq i \leq N$, and for any submodule $\Delta \subset \mathcal{O}_S^n$

$$\sup_{s \in U_i \times \prod_{p \in S_f} \mathbb{Z}_p^l} \operatorname{cov}(g\phi_i(s)^{-1}\Delta) \ge \rho.$$

Then we apply [KT07, Theorem 9.3] and for any $\epsilon > 0$, there exists a compact subset Q_{ϵ} such that for each i = 1, 2, ..., N

$$\mu(\{s \in U_i \times \prod_{p \in S_f} \mathbb{Z}_p^l : g\phi_i(s)^{-1}\mathcal{O}_S^n \notin Q_\epsilon\}) \le \frac{\epsilon}{N} \mu(U_i \times \prod_{p \in S_f} \mathbb{Z}_p^l).$$

Now we compute

$$\mu_{H}(\{s \in \Omega : \Gamma_{S}sg^{-1} \notin Q_{\epsilon}\})$$

$$\leq \sum_{i=1}^{N} \mu_{H}(\{s \in B_{i} : \Gamma_{S}sg^{-1} \notin Q_{\epsilon}\}) + \mu_{H}(\Omega - W)$$

$$\leq C_{1} \sum_{i=1}^{N} \mu(\{s \in U_{i} \times \prod_{p \in S_{f}} \mathbb{Z}_{p}^{l} : g\phi_{i}(s)^{-1}\mathcal{O}_{S}^{n} \notin Q_{\epsilon}\}) + \epsilon$$

$$\leq \frac{C_{1}\epsilon}{N} \sum_{i=1}^{N} \mu\left(U_{i} \times \prod_{p \in S_{f}} \mathbb{Z}_{p}^{l}\right) + \epsilon \leq C\epsilon\mu_{H}(\Omega) + \epsilon = (C+1)\epsilon$$

where C_1 and C depend only on Ω , $\epsilon > 0$ and the local diffeomorphisms ϕ_i (i = 1, 2, ..., N). Hence we have $\mu_{\Gamma_S \setminus \Gamma_S Hg^{-1}}(Q^c_{\epsilon}) \leq C\epsilon$, for some constant C > 0. This completes the proof of the theorem.

3.2. Translates of a maximal torus visiting a compact subset. Let $\mathbf{H} \subset \mathbf{G}$ be connected reductive algebraic groups defined over \mathbb{Q} without nontrivial \mathbb{Q} -characters and suppose that \mathbf{H} contains a maximal torus of \mathbf{G} . Let H be a subgroup of finite index in $\mathbf{H}(\mathbb{Q}_S)$. Let $U^* = \prod_{p \in S} U_p$ in H ($U_p \subset \mathbf{H}(\mathbb{Q}_p)$) be defined as in §3.1, which is a neighborhood of identity in H such that the exponential map is defined on U_p ($\forall p \in S$) and U_p is a subgroup of $\mathbf{H}(\mathbb{Q}_p)$ ($\forall p \in S_f$). In this subsection, we prove the following theorem.

Theorem 3.2. There exists a compact subset K in $\Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$ such that for any $g \in \mathbf{G}(\mathbb{Q}_S)$

$$\Gamma_S \backslash \Gamma_S U^* q \cap K \neq \emptyset.$$

We know that **H** contains a maximal torus in **G**. It implies that **H** contains a maximal torus **T** of **G** which is defined and anisotropic over \mathbb{Q} [PR72]. In the rest of this section, we will fix this maximal torus **T**. Now we take a small neighborhood of identity Ω in $\mathbf{T}(\mathbb{Q}_S)$ such that $\Omega \subset U^*$. Then it suffices to prove the following

Theorem 3.3. There exists a compact subset $K \subset \Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$ such that for any g in $\mathbf{G}(\mathbb{Q}_S)$

$$\Gamma_S \backslash \Gamma_S \Omega q \cap K \neq \emptyset$$
.

In order to prove Theorem 3.3, we need the following proposition.

Proposition 3.4. There exists a subset Y_p of $\mathbf{G}(\mathbb{Q}_p)$ such that

- (1) $\mathbf{G}(\mathbb{Q}_p) = Y_p \mathbf{T}(\mathbb{Q}_p)$
- (2) For any linear representation of $\mathbf{G}(\mathbb{Q}_p)$ on a vector space V over \mathbb{Q}_p , there exists a constant c > 0 such that

$$\sup_{t \in \Omega} \|yt.v\| \ge c\|v\|, \quad \forall v \in V, \quad \forall y \in Y_p.$$

Proof. The case $p = \infty$ is essentially proved in [EMS97, Proposition 4.4]. Here we give a proof for the p-adic case. Let F be a Galois extension of \mathbb{Q}_p such that \mathbf{T} is F-split. By Iwasawa decomposition, there exists a compact subgroup K_F of $\mathbf{G}(F)$ and a unipotent subgroup U_F of $\mathbf{G}(F)$ such that

$$G(F) = K_F U_F \mathbf{T}(F).$$

Let $\{\sigma_1, \sigma_2, ..., \sigma_n\} = \operatorname{Gal}(F/\mathbb{Q}_p)$ and define

$$S = \{x \in \mathbf{T}(F) : \sigma_1(x)\sigma_2(x)...\sigma_n(x) = e\}.$$

We prove that $ST(\mathbb{Q}_p)$ is a subgroup of finite index in T(F). We know that $T(\mathbb{Q}_p)^n$ is a subgroup of finite index in $T(\mathbb{Q}_p)$. Let

$$E = \{x \in \mathbf{T}(F) : \sigma_1(x)\sigma_2(x)\cdots\sigma_n(x) \in \mathbf{T}(\mathbb{Q}_p)^n\}.$$

Then E is a subgroup of finite index in $\mathbf{T}(F)$. Now for any $x \in E$, there exists $w \in \mathbf{T}(\mathbb{Q}_p)$ such that

$$\sigma_1(x)\sigma_2(x)\cdots\sigma_n(x)=w^n$$

and hence $w^{-1}x \in S$, $x \in S\mathbf{T}(\mathbb{Q}_p)$. This implies that $E = S\mathbf{T}(\mathbb{Q}_p)$ and $S\mathbf{T}(\mathbb{Q}_p)$ is a subgroup of finite index in $\mathbf{T}(F)$.

Let Λ be a finite set of representatives of $\mathbf{T}(F)/(S\mathbf{T}(\mathbb{Q}_p))$ in $\mathbf{T}(F)$. Since Λ normalizes U_F , we have

$$\mathbf{G}(F) = K_F U_F \Lambda S \mathbf{T}(\mathbb{Q}_p) = (K_F \Lambda) U_F S \mathbf{T}(\mathbb{Q}_p) = K U_F S \mathbf{T}(\mathbb{Q}_p)$$

where K is the compact subset $K_F\Lambda$ in $\mathbf{G}(F)$. Let $Y_p = (KU_FS) \cap \mathbf{G}(\mathbb{Q}_p)$. Then $\mathbf{G}(\mathbb{Q}_p) = Y_p\mathbf{T}(\mathbb{Q}_p)$. We prove that the subset Y_p satisfies the second claim in the proposition. Let ρ be a representation of $\mathbf{G}(\mathbb{Q}_p)$ on a space V over \mathbb{Q}_p , and we extend it to a presentation of $\mathbf{G}(F)$ over F. Let $y = kus \in Y_p$ for $k \in K$, $u \in U_F$ and $s \in S$. Then assuming Lemmas 3.5 and 3.6 below, we have for any $v \in V$

$$\sup_{t \in \Omega} \|yt.v\| = \sup_{t \in \Omega} \|kust.v\| \ge c_1 \sup_{t \in \Omega} \|ut.(s.v)\| \ge c_2 \sup_{t \in \Omega} \|(t^{-1}ut).(s.v)\|$$
$$> c_2c_3\|s.v\| > c_2c_3c_4\|v\|.$$

Hence we only need to prove the following lemmas.

Lemma 3.5. There exists a constant $c_3 > 0$ such that for any $v \in V$

$$\sup_{t \in \Omega} \|(t^{-1}ut).v\| \ge c_3\|v\|$$

Proof. Let Φ be the root system of $(\mathbf{G}(F), \mathbf{T}(F))$ and Φ_+ the set of positive roots. Let $\alpha_1, \alpha_2, \ldots, \alpha_n$ be the positive simple roots of $\mathbf{T}(F)$ in Φ_+ . Then there is an isomorphism $f_{\Phi} : \mathbf{T}(F) \to (F_{\times})^n$ which maps $x \in \mathbf{T}(F)$ to

$$f_{\Phi}(x) := (\alpha_1(x), \alpha_2(x), \dots, \alpha_n(x)).$$

Write $u = \exp(\overline{u})$ where

$$\overline{u} = \sum_{\alpha \in \Phi_+} u_{\alpha} \in \sum_{\alpha \in \Phi_+} \mathfrak{g}_{\alpha}.$$

Note that for $t \in \Omega$, $t^{-1}ut.v = \exp(\sum_{\alpha \in \Phi_+} \alpha(t)u_{\alpha})$. Since the positive roots are generated by simple roots, it follows that the set $\{t^{-1}ut: t \in \Omega\}$ is the image of the map f defined by

$$(\alpha_1(t), \alpha_2(t), \dots, \alpha_n(t)) \mapsto \exp(\sum_{\alpha \in \Phi_+} \alpha(t)u_\alpha).$$

We conclude that $\{t^{-1}ut: t \in \Omega\}$ is the image of a polynomial map f defined on $f_{\Phi}(\Omega)$ which maps $s = (s_1, s_2, \dots, s_n) \in f_{\Phi}(\Omega)$ to

$$f(s_1, s_2, \dots, s_n) := \exp(\sum_{\alpha \in \Phi_+} \alpha(s_1, s_2, \dots, s_n) u_\alpha)$$

where $\alpha(s_1, s_2, ..., s_n)$ is a polynomial in $s_1, s_2, ..., s_n$ without constant term.

To finish the proof of the lemma, it is enough to prove that there exists a constant $c_3 > 0$ such that

$$\sup_{s \in f_{\Phi}(\Omega)} \|f(s).v\| \ge c_3 \|v\|.$$

Since f is a polynomial map on each variable s_i with a uniformly bounded degree, by [KT07, Corollary 3.3], f is (C, α) -good for some C and α . We also have f(0) = e, so f(0).v = v. Now choose a bounded region D in F^n containing $f_{\Phi}(\Omega)$, and by the (C, α) -good property, there are constants $c_3 > 0$ depending only on C, α and $f_{\Phi}(\Omega)$ such that

$$\mu_{F^n}(\{s \in D : ||f(s).v|| \le c_3||v||\}) \le \frac{1}{2}\mu_{F^n}(f_{\Phi}(\Omega))$$

where μ_{F^n} denotes the Haar measure on F^n . Hence there is $s \in f_{\Phi}(\Omega)$ such that $||f(s).v|| \ge c_3 ||v||$. This completes the proof of the lemma.

Lemma 3.6. There exists a constant $c_4 > 0$ such that for any $v \in V$

$$||s.v|| \ge c_4 ||v||.$$

Proof. We write the weight space decomposition for V_F over F as

$$V = \sum_{\alpha} V_{\alpha}$$

and for any $v \in V$ we write $v = \sum_{\alpha} v_{\alpha}$. Then $s.v = \sum_{\alpha} \alpha(s)v_{\alpha}$ and

$$\sigma_i(s.v) = \sigma_i(s).v = \sum_{\alpha} \alpha(\sigma_i(s))v_{\alpha}.$$

We know that $\|\cdot\|$ and $\|\sigma_i(\cdot)\|$ are equivalent norms on the vector space V_F over F, and for any $v \in V_F$, there exists a constant $\kappa > 0$ such that

$$\frac{1}{\kappa} \|v\| \le \|\sigma_i(v)\| \le \kappa \|v\|$$

$$\frac{1}{\kappa} \max_{\alpha} \|v_{\alpha}\| \le \|v\| \le \kappa \max_{\alpha} \|v_{\alpha}\|.$$

For any $v \in V$, let $v_{\beta} = \max ||v_{\alpha}||$. Then

$$||s.v|| \ge \frac{1}{\kappa} ||\sigma_i(s).v|| \ge \frac{1}{\kappa^2} ||\beta(\sigma_i(s)).v_\beta||$$

$$||s.v||^n \ge \frac{1}{\kappa^{2n}} ||\beta(\prod_{i=1}^n \sigma_i(s))|| ||v_\beta||^n \ge \frac{1}{\kappa^{3n}} ||v||^n.$$

This implies that $||s.v|| \ge c_4 ||v||$ for some constant $c_4 > 0$.

This completes the proof of the proposition.

Proof of Theorem 3.3. We embed **G** in \mathbf{SL}_n for some n, and hence

$$\mathbf{G}(\mathcal{O}_S)\backslash\mathbf{G}(\mathbb{Q}_S)\subset\mathbf{SL}_n(\mathcal{O}_S)\backslash\mathbf{SL}_n(\mathbb{Q}_S)\subset\mathbf{GL}_n^1(\mathcal{O}_S)\backslash\mathbf{GL}_n^1(\mathbb{Q}_S)$$

where $\mathbf{GL}_n^1(\mathcal{O}_S)\backslash \mathbf{GL}_n^1(\mathbb{Q}_S)$ is the space of unimodular lattices $\Omega_{S,n}^1$ in \mathbb{Q}_S^n . Note that the group action of $\mathbf{G}(\mathbb{Q}_S)$ on $\mathbf{GL}_n^1(\mathcal{O}_S)\backslash \mathbf{GL}_n^1(\mathbb{Q}_S)$ is defined by

$$\Delta \mapsto h^{-1}\Delta, \ \forall h \in \mathbf{G}(\mathbb{Q}_S), \ \forall \Delta \in \Omega^1_{S,n}$$

For any $p \in S$, let Y_p be the subset of $\mathbf{G}(\mathbb{Q}_p)$ in Proposition 3.4 and we have $\mathbf{G}(\mathbb{Q}_p) = \mathbf{T}(\mathbf{Q}_p) Y_p^{-1}$. For $g \in \mathbf{G}(\mathbb{Q}_S)$, we can write $g = (t_p y_p^{-1})_{p \in S}$ where $t_p \in \mathbf{T}(\mathbb{Q}_p)$ and $y_p \in Y_p^{-1}$ and we have

$$\Gamma_S \backslash \Gamma_S \Omega g = \Gamma_S \backslash \Gamma_S (t_p)_{p \in S} \Omega(y_p^{-1})_{p \in S}.$$

Since **T** is \mathbb{Q} -anisotropic, $\mathbf{T}(\mathcal{O}_S)\backslash \mathbf{T}(\mathbb{Q}_S)$ is compact and there exists $s \in \mathbf{T}(\mathbb{Q}_S)$ in a compact fundamental domain of $\mathbf{T}(\mathcal{O}_S)\backslash \mathbf{T}(\mathbb{Q}_S)$ such that $\Gamma_S(t_p)_{p\in S} = \Gamma_S s$. Let $y = (y_p)_{p\in S}$. We have

$$\Gamma_S \backslash \Gamma_S \Omega g = \Gamma_S \backslash \Gamma_S s \Omega y^{-1}.$$

Let ρ_k be the representation of $\mathbf{G}(\mathbb{Q}_S)$ on the k-th exterior product space V_k of \mathbb{Q}_S^n $(1 \le k \le n)$. Then by Proposition 3.4, for any $v \in V_k$

$$\sup_{t \in \Omega} c((yt)s.v) \ge C_0 c(s.v)$$

for some constant $C_0 > 0$. Here c(v) denotes the content of the element v. This implies that for every submodule $\Delta \subset \mathcal{O}_S^n$

$$\sup_{t \in \Omega} \operatorname{cov}((yt)s.\Delta) \ge C_0 \operatorname{cov}(s.\Delta) \ge C_1 \operatorname{cov}(\Delta)$$

for some constant $C_1 > 0$, as s is in a fixed compact subset. Since Ω is an image of a (C, α) -good map, by [KT07, Lemma 8.2], we can apply [KT07, Theorem 9.3] and there is a compact subset K in $\Omega^1_{S,n}$ depending only on n, Ω , S, C and α such that $yts\mathcal{O}^n_S \in K$ for some $t \in \Omega$. This implies that for any $g \in \mathbf{G}(\mathbb{Q}_S)$, $\Gamma_S \backslash \Gamma_S \Omega g \cap K \neq \emptyset$.

We conclude this section by the following theorem, which is a corollary of Theorems 3.1 and 3.2.

Theorem 3.7. Let G be a reductive algebraic group defined over \mathbb{Q} , which does not admit nontrivial \mathbb{Q} -characters. Let G be a reductive \mathbb{Q} -subgroup of G without nontrivial \mathbb{Q} -characters, which contains a maximal torus of G. Let G be a finite set of places of \mathbb{Q} containing the Archimedean place, G a subgroup of finite index in G and G and G and G are G are G and G are G and G are G are G and G are G are G are G are G and G are G are G are G and G are G are G are G are G are G are G and G are G and G are G and G are G are G and G are G and G are G and G are G and G are G are G are G are G and G are G are G are G are G are G and G are G

4. Auxiliary results in group theory

In this section, we prove some auxiliary lemmas for several arguments in §5.

Recall that H is a subgroup of finite index in $\mathbf{H}(\mathbb{Q}_S)$. For any $p \in S_f$, we denote by H_p a small open compact subgroup of $\mathbf{H}(\mathbb{Q}_p)$ contained in H, where the exponential map \exp_p is defined. The subgroup H_p could be considered as a closed subgroup in $\mathbf{SL}_n(\mathbb{Q}_p)$ for some $n \in \mathbb{N}$. Note that all the eigenvalues of elements in H_p stay in a small neighborhood of $1 \in \mathbb{Q}_p$. In the following, for a sequence $\{W_i\}_{i\in\mathbb{N}}$ of subsets in $\mathbf{SL}_n(\mathbb{Q}_p)$, we say $\{W_i\}$ is unbounded if for any compact subset $K \subset \mathbf{SL}_n(\mathbb{Q}_p)$, there are infinitely many W_i such that $W_i \not\subset K$. Otherwise, we say $\{W_i\}$ is bounded.

Lemma 4.1. Let $\{\gamma_i\}$ be a sequence in $\mathbf{SL}_n(\mathbb{Q}_p)$. Suppose $\{\operatorname{Ad} \gamma_i(H_p)\}\ (i \in \mathbb{N})$ is unbounded. Then for any N > 0, after passing to a subsequence, there exists a sequence $h_i \in H_p$ such that

$$h_i \to e_p$$
, $\operatorname{Ad} \gamma_i(h_i) \to u \quad (i \to \infty)$

for some unipotent element u in $\mathbf{SL}_n(\mathbb{Q}_p)$ with $|u|_p \geq N$.

Proof. Since $\{\operatorname{Ad} \gamma_i(H_p)\}$ is unbounded, by passing to a subsequence, there is $x_i = \exp(X_i)$ in H_p with $X_i \in \operatorname{Lie}(\operatorname{\mathbf{SL}}_n(\mathbb{Q}_p))$ such that $\operatorname{Ad} \gamma_i(x_i) \to \infty$. By Proposition 2.2 and [KT07, Lemma 3.1], we know that the function $f_i : \mathbb{Z}_p \to \mathbb{R}$ defined by

$$f_i(s) = \|\gamma_i(\exp(sX_i))\gamma_i^{-1} - e_p\|_p, \quad s \in \mathbb{Z}_p$$

is a (C, α) -good function with some $C, \alpha > 0$ uniformly for any $i \in \mathbb{N}$. Our choices of x_i 's then imply $f_i(1) \to \infty$ as $i \to \infty$. By Corollary 2.4 we can find $s_i \in \mathbb{Z}_p$ so that

$$s_i X_i \to 0, \quad f_i(s_i) \to \infty.$$

Therefore, by replacing X_i with s_iX_i , we may assume $X_i \to 0$. Now for any N > 0, one can find $k \in \mathbb{N}$ such that $C'^k \geq N$ where C' is the constant in Proposition 2.3. Since $f_i(0) = 0$ and $f_i(s_i) \to \infty$ as $i \to \infty$, we apply Proposition 2.3 with $\rho = C'^{k+1}$, and one can find $v_i \in \mathbb{Z}_p$ such that

$$f_i(v_i) \in [C'^k, C'^{k+1}].$$

Let $h_i = \exp(v_i X_i)$. Then by passing to a subsequence, we have

$$h_i \to e_p$$
, Ad $\gamma_i(h_i) \to u$

for some $u \in \mathbf{SL}_n(\mathbb{Q}_p)$ with $|u - e_p|_p \in [C'^k, C'^{k+1}]$. Since $h_i \to e_p$, all the eigenvalues of $\mathrm{Ad}\,\gamma_i(h_i)$ converges to 1. It follows that u is a unipotent element with $|u|_p \geq N$.

Lemma 4.2. A closed subgroup of $\mathbf{SL}_n(\mathbb{Q}_p)$ containing arbitrarily large unipotent elements has a one-parameter unipotent subgroup.

Proof. The exponential map $\exp(X)$ can be defined for nilpotent elements X in the Lie algebra of $\mathbf{SL}_n(\mathbb{Q}_p)$, and the log function $\log(u)$ can be defined for unipotent elements u in $\mathbf{SL}_n(\mathbb{Q}_p)$. Moreover, the sets of nilpotent elements and of unipotent elements are both closed. Now let $\{u_i\}$ be a sequence of unipotent elements in a closed subgroup M with $u_i \to \infty$. We take $X_i = \log u_i$ and we have $X_i \to \infty$. By passing to a subsequence, $|X_i|_p X_i$ converges to a nilpotent element X with $|X|_p = 1$. Now we prove that the one-parameter unipotent subgroup

$$s \in \mathbb{Q}_p \mapsto \exp(sX)$$

is contained in M.

For $s \neq 0$ and i sufficiently large, since $\overline{\mathbb{Z}} = \mathbb{Z}_p$, one can find $m_i \in \mathbb{Z}$ such that

$$|m_i - s|X_i|_p|_p \le \frac{1}{i|X_i|_p}.$$

This implies that

$$m_i X_i = (m_i - s|X_i|_p) X_i + s|X_i|_p X_i \rightarrow sX.$$

Since $\exp(m_i X_i) = u_i^{m_i} \in M$ and M is closed, we have $\exp(sX) \in M$. This completes the proof of the lemma.

The following is an immediate corollary of Lemmas 4.1 and 4.2.

Proposition 4.3. Let $\{\mu_i\}$ be a sequence of probability measures on $\Gamma\backslash \mathbf{G}(\mathbb{Q}_S)$ which weakly converges to a measure μ , and p a non-Archimedean place in S. Suppose that each measure μ_i is invariant under $\mathrm{Ad}(\gamma_i)H_p$ for some $\gamma_i \in \mathbf{G}(\mathbb{Q}_p)$, and $\{\mathrm{Ad}(\gamma_i)H_p\}$ is unbounded. Then μ is invariant under a one-parameter unipotent subgroup.

Now we consider the case when $\{Ad(\gamma_i)H_p\}$ is bounded.

Proposition 4.4. Let $\gamma_i \in \mathbf{G}(\mathbb{Q}_p)$. The sequence $\{\mathrm{Ad}(\gamma_i)H_p\}$ is bounded if and only if $\gamma_i Z_{\mathbf{G}}(\mathbf{H})(\mathbb{Q}_p)$ is bounded in $\mathbf{G}(\mathbb{Q}_p)/Z_{\mathbf{G}}(\mathbf{H}(\mathbb{Q}_p))$.

Proof. The set of semisimple elements in H_p is Zariski dense in $\mathbf{H}(\mathbb{Q}_p)$, so one can pick a finite number of semisimple elements $\alpha_i \in H_p$ such that

$$Z_{\mathbf{G}}(\alpha_1) \cap Z_{\mathbf{G}}(\alpha_2) \cap ... \cap Z_{\mathbf{G}}(\alpha_n) = Z_{\mathbf{G}}(\mathbf{H}(\mathbb{Q}_p)).$$

Consider the action of $\mathbf{G}(\mathbb{Q}_p)$ on itself by conjugation. Under this action, the stabilizer of α_i in $\mathbf{G}(\mathbb{Q}_p)$ equals $Z_{\mathbf{G}}(\alpha_i)$. This induces a map $\beta_i : \mathbf{G}(\mathbb{Q}_p)/Z_{\mathbf{G}}(\alpha_i) \to \mathbf{G}(\mathbb{Q}_p)$ defined by

$$\alpha(gZ_{\mathbf{G}}(\alpha_i)) = g\alpha_i g^{-1}.$$

By [BT65, Bor91] or [MT94, Proposition 1.6], this map is a homeomorphism between $\mathbf{G}(\mathbb{Q}_p)/Z_{\mathbf{G}}(\alpha_i)$ and the conjugacy class of α_i .

Suppose $\{\operatorname{Ad}(\gamma_i)H_p\}$ is bounded. Then the sequence $\{\operatorname{Ad}(\gamma_i)\alpha_k\}$ is bounded $(\forall k=1,2,...,n)$. The homeomorphism β_i then implies that the sequence $\{\gamma_i Z_{\mathbf{G}}(\alpha_k)\}$ is also bounded in $\mathbf{G}(\mathbb{Q}_p)/Z_{\mathbf{G}}(\alpha_k)$. Consequently, this implies that $\gamma_i Z_{\mathbf{G}}(\mathbf{H}(\mathbb{Q}_p))$ is bounded in $\mathbf{G}(\mathbb{Q}_p)/Z_{\mathbf{G}}(\mathbf{H}(\mathbb{Q}_p))$. Conversely, if $\{\gamma_i Z_{\mathbf{G}}(\mathbf{H}(\mathbb{Q}_p))\}$ is bounded in $\mathbf{G}(\mathbb{Q}_p)/Z_{\mathbf{G}}(\mathbf{H}(\mathbb{Q}_p))$, then one can find a compact subset K in $\mathbf{G}(\mathbb{Q}_p)$ such that $\gamma_i \in KZ_{\mathbf{G}}(\mathbf{H}(\mathbb{Q}_p))$. Then $\mathrm{Ad}(\gamma_i)H_p$ is contained in the compact subset KH_pK^{-1} . \square

The following is an analogue of the previous results over \mathbb{R} which can be proved in a similar manner.

Proposition 4.5. Let $\{\gamma_i\}$ be a sequence in $\mathbf{G}(\mathbb{R})$. Consider the operator

$$\operatorname{Ad} \gamma_i : \operatorname{Lie}(Z_{\mathbf{G}}(\mathbf{H})(\mathbb{R})) \to \operatorname{Lie}(\mathbf{G}(\mathbb{R})).$$

Then we have

- (1) If the sequence $\{\|\operatorname{Ad}\gamma_i\|\}$ is unbounded, then by passing to a subsequence, there exists a one-parameter unipotent subgroup $\{u(t):t\in\mathbb{R}\}$ in $\mathbf{G}(\mathbb{R})$ satisfying the following condition: for any $s\in\mathbb{R}$, one can find a sequence $w_i\in Z_{\mathbf{G}}(\mathbf{H})(\mathbb{R})^0$ such that $\gamma_iw_i\gamma_i^{-1}\to u(s)$.
- (2) The sequence $\{\|\operatorname{Ad}(\gamma_i)\|\}$ is bounded if and only if $\gamma_i Z_{\mathbf{G}}(\mathbf{H})(\mathbb{R})$ is bounded in $\mathbf{G}(\mathbb{R})/Z_{\mathbf{G}}(\mathbf{H})(\mathbb{R})$. Here $\|\cdot\|$ denotes the operator norm.

Similarly, we obtain the following

Proposition 4.6. Let $\{\mu_i\}$ be a sequence of probability measures on $\Gamma\backslash \mathbf{G}(\mathbb{Q}_S)$ which weakly converges to a measure μ . Suppose that each measure μ_i is invariant under $\mathrm{Ad}(\gamma_i)Z_{\mathbf{G}}(\mathbf{H})(\mathbb{R})^0$ for some $\gamma_i \in \mathbf{G}(\mathbb{R})$, and the norms of the operators $\mathrm{Ad}(\gamma_i) : \mathrm{Lie}(\mathbf{T}(\mathbb{R})) \to \mathrm{Lie}(\mathbf{G}(\mathbb{R}))$ are unbounded. Then μ is invariant under a one-parameter unipotent subgroup.

5. Measure rigidity: Proof of Theorem 1.7

In this section, we prove Theorem 1.7 using S-arithmetic Ratner's theorem. We will mainly follow the frameworks and the presentations of [EMS96, GO11].

In order to prove Theorem 1.7, we make the following reduction. We know that the sequence $g_i^*\mu_{\Gamma_S\backslash\Gamma_S H}$ converges to a probability measure μ . Then there exist $a_i \in H$, $\gamma_i \in \Gamma_S$ and $g \in \mathbf{G}(\mathbb{Q}_S)$ such that $\pi_S(g) \in \text{supp}(\mu)$ and $\gamma_i a_i g_i^{-1} \to g$. We have

$$\gamma_i^* \mu_{\Gamma_S \backslash \Gamma_S H} = (\gamma_i a_i g_i^{-1})^* (g_i)^* \mu_{\Gamma_S \backslash \Gamma_S H} \to g^* \mu.$$

Therefore, to prove that μ is an algebraic measure, it suffices to study the sequence of probability measures $\gamma_i^* \mu_{\Gamma_S \backslash \Gamma_S H}$ where $\{\gamma_i\}$ is a sequence in Γ_S . Now for each $\gamma_i \in \Gamma_S$, we denote by $\rho_i(x) := \gamma_i x \gamma_i^{-1}$ the conjugation of γ_i in **G**. Then

$$\gamma_i^* \mu_{\Gamma_S \backslash \Gamma_S H} = \mu_{\Gamma_S \backslash \Gamma_S \rho_i(H)}$$

is the Haar measure on the finite-volume orbit $\Gamma_S \backslash \Gamma_S \rho_i(H) \subset \Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$. Without loss of generality, we may assume that $\mu_{\Gamma_S \backslash \Gamma_S \rho_i(H)}$ converges to a probability measure μ .

Let **T** be a maximal torus in **G** which is defined and anisotropic over \mathbb{Q} , and contained in **H**. Suppose that there are infinitely many $i \in \mathbb{N}$ such that $\gamma_i \mathbf{H} \gamma_i^{-1}$ is contained in a connected \mathbb{Q} -subgroup **F** of dimension smaller than dim **G**. By [EMS97, Lemma 5.1] and the condition that $\mathbf{T} \subset \mathbf{F}$, **F** is reductive and has no non-trivial \mathbb{Q} -characters. Since

$$\mathbf{T}(\mathbb{C}) \subset \mathbf{F}(\mathbb{C}) \text{ and } \gamma_i \mathbf{T}(\mathbb{C}) \gamma_i^{-1} \subset \mathbf{F}(\mathbb{C}),$$

and both $\mathbf{T}(\mathbb{C})$ and $\gamma_i \mathbf{T}(\mathbb{C}) \gamma_i^{-1}$ are maximal tori in $\mathbf{F}(\mathbb{C})$, there exists an element $h_i \in \mathbf{F}(\mathbb{C})$ such that $h_i \mathbf{T}(\mathbb{C}) h_i^{-1} = \gamma_i \mathbf{T}(\mathbb{C}) \gamma_i^{-1}$. Hence $h_i^{-1} \gamma_i$ is in the normalizer of $\mathbf{T}(\mathbb{C})$. Let W be a finite set of representatives of Weyl group in $\mathbf{G}(\mathbb{C})$. Then $h_i^{-1} \gamma_i \in \mathbf{T}(\mathbb{C}) W$. This implies that there exists $w_i \in W$ such that $\gamma_i w_i^{-1} \in \mathbf{F}(\mathbb{C})$. By passing to a subsequence, we may assume that $w_i^{-1} = w$ for all $i \in \mathbb{N}$. Now let $\delta_i = \gamma_i \gamma_1^{-1}$ and $\mathbf{H}' = \gamma_1 \mathbf{H} \gamma_1^{-1}$. Then we have

$$\delta_i = \gamma_i w (\gamma_1 w)^{-1} \in \mathbf{F}(\mathbb{C}) \cap \mathbf{G}(\mathcal{O}_S) = \mathbf{F}(\mathcal{O}_S)$$

and $\{\delta_i \mathbf{H}' \delta_i^{-1}\} = \{\gamma_i \mathbf{H} \gamma_i^{-1}\}$ is contained in the \mathbb{Q} -reductive group \mathbf{F} . If we write the conjugation by δ_i as $\rho'_i(x) := \delta_i x \delta_i^{-1}$ and $H' = \gamma_1 H \gamma_1^{-1}$, then

$$\mu_{\Gamma_S \setminus \Gamma_S \rho_i(H)} = \mu_{\gamma_i H \gamma_i^{-1}} = \mu_{\delta_i H' \delta_i^{-1}} = \mu_{\Gamma_S \setminus \Gamma_S \rho_i'(H')}$$

is actually a sequence of probability measures on $\mathbf{F}(\mathcal{O}_S)\backslash\mathbf{F}(\mathbb{Q}_S)$. Therefore, without loss of generality, we may assume that \mathbf{G} is the smallest \mathbb{Q} -subgroup containing the sequence $\{\gamma_i\mathbf{H}\gamma_i^{-1}\}$ and $\gamma_i \in \mathbf{G}(\mathcal{O}_S)$.

Moreover, it is harmless to replace Γ_S by a subgroup of finite index Γ in Γ_S , and consider the homogeneous space $\Gamma \backslash \mathbf{G}(\mathbb{Q}_S)$. Let $Z(\mathbf{G})$ be the center of \mathbf{G} . Then by passing to the quotient space $\mathbf{G}/Z(\mathbf{G})$, we may further assume that \mathbf{G} is semisimple. We will denote by $\pi_S : \mathbf{G}(\mathbb{Q}_S) \to \Gamma \backslash \mathbf{G}(\mathbb{Q}_S)$ be the natural quotient map.

Theorem 5.1. Let G be a connected semisimple algebraic group defined over \mathbb{Q} . Let H be a connected reductive \mathbb{Q} -subgroup without nontrivial \mathbb{Q} -characters, and assume that H contains a maximal torus in G. Let $\rho_i : H \to G(i \in \mathbb{N})$ be a sequence of \mathbb{Q} -homomorphisms defined by $\rho_i(x) = \gamma_i x \gamma_i^{-1}$ ($\gamma_i \in \Gamma$) where Γ is a subgroup of finite index in Γ_S , such that there exists no proper \mathbb{Q} -subgroup M containing infinitely many $\rho_i(H)$. Let H be a subgroup of finite index in $H(\mathbb{Q}_S)$. Suppose that the sequence $\{\mu_{\Gamma \setminus \Gamma \rho_i(H)}\}$ converges to a probability measure μ as $i \to \infty$. Then μ is an algebraic measure on $\Gamma \setminus G(\mathbb{Q}_S)$.

Theorem 1.7 is then an immediate corollary of Theorem 3.7 and Theorem 5.1. The rest of this section will be devoted to the proof of Theorem 5.1.

5.1. **Prerequisites.** Let \mathbf{M} be an algebraic variety defined over \mathbb{Q} , and S a finite set of places of \mathbb{Q} containing the Archimedean place. The Zariski topology on $\mathbf{M}(\mathbb{Q}_S) = \prod_{p \in S} \mathbf{M}(\mathbb{Q}_p)$ is defined to be the product of the Zariski topologies on $\mathbf{M}(\mathbb{Q}_p)$ $(p \in S)$. On the other hand, the topologies on the local fields \mathbb{Q}_p $(p \in S)$ induce a topology on $\mathbf{M}(\mathbb{Q}_S) = \prod_{p \in S} \mathbf{M}(\mathbb{Q}_p)$, which we will refer to as Hausdorff topology on $\mathbf{M}(\mathbb{Q}_S)$.

In the following, a map $f: \mathbf{M}(\mathbb{Q}_S) \to \mathbb{Q}_S$ of the form $f = (f_p)_{p \in S}$, where f_p is a \mathbb{Q}_p -valued function on $\mathbf{M}(\mathbb{Q}_p)$, will be called a \mathbb{Q}_S -valued function on $\mathbf{M}(\mathbb{Q}_S) = \prod_{p \in S} \mathbf{M}(\mathbb{Q}_p)$. The space of \mathbb{Q}_S -valued functions on $\mathbf{M}(\mathbb{Q}_S)$ is a \mathbb{Q}_S -module. A \mathbb{Q}_S -valued function on $\mathbf{M}(\mathbb{Q}_S)$ is called regular if f_p is regular for every $p \in S$.

The following lemma could be verified by [Bor91, Ch.1].

Lemma 5.2. Let $\rho_i : \mathbf{H} \to \mathbf{G}$ $(i \in \mathbb{N})$ be a sequence of \mathbb{Q} -homomorphisms defined by $\rho_i(x) = \gamma_i x \gamma_i^{-1}$, and \mathbf{V} an affine variety on which \mathbf{G} acts morphically. Then for any regular \mathbb{Q}_S -valued function f on $\mathbf{V}(\mathbb{Q}_S)$, the maps $f(\rho_i|_{\mathbf{H}(\mathbb{Q}_S)}w)$ $(w \in \mathbf{V}(\mathbb{Q}_S), i \in \mathbb{N})$ span a finitely generated \mathbb{Q}_S -module in the space of \mathbb{Q}_S -valued functions on $\mathbf{H}(\mathbb{Q}_S)$.

The following is an immediate corollary of [PR94, Theorem 7.9].

Lemma 5.3. Let **H** be as in Theorem 5.1, and S a finite set of places of \mathbb{Q} containing the Archimedean place. Then there exists an open neighborhood of identity W_{S_f} in $\mathbf{H}(\mathbb{Q}_{S_f})$ such that $\mathbf{H}(\mathbb{Q})$ is dense in $\mathbf{H}(\mathbb{R}) \times W_{S_f}$.

5.2. **Invariant measures for unipotent flows and linearization.** In this subsection, we briefly discuss invariant measures for unipotent flows and linearization in the S-arithmetic setting. One can refer to e.g. [MT94, MT96, Rat95, Rat98, Tom00, GO11] for more details. Here we will follow the notation and the formulations in [GO11].

We say that a connected \mathbb{Q} -subgroup \mathbf{P} of \mathbf{G} is in class \mathcal{F} ralative to S if the radical of \mathbf{P} is unipotent and every \mathbb{Q} -simple factor of \mathbf{P} is \mathbb{Q}_v -isotropic for some $v \in S$. A closed subgroup L of $\mathbf{G}(\mathbb{Q}_S)$ is in class \mathcal{H} if there exists a connected \mathbb{Q} -subgroup \mathbf{P} in class \mathcal{F} relative to S such that L has a finite index in $\mathbf{P}(\mathbb{Q}_S)$ and L_u acts ergodically on $\Gamma \setminus \Gamma L$ with respect to the L-invariant probability measure. Here L_u is the subgroup in L generated by unipotent subgroups.

For a closed subgroup L of $\mathbf{G}(\mathbb{Q}_S)$, the Mumford-Tate subgroup, denoted by MT(L), is defined to be the smallest connected \mathbb{Q} -subgroup \mathbf{G} such that

$$\bar{L}^0 \subset \prod_{v \in S} MT(L)(\mathbb{Q}_v)$$

where \bar{L}^0 denotes the identity component of the Zariski closure of L in $\mathbf{G}(\mathbb{Q}_S)$.

Let $X_S = \Gamma \backslash \mathbf{G}(\mathbb{Q}_S)$. Let W be a closed subgroup of $\mathbf{G}(\mathbb{Q}_S)$ generated by one-parameter unipotent subgroups. For each $L \in \mathcal{H}$, define

$$\mathcal{N}(L,W) = \{ g \in \mathbf{G}(\mathbb{Q}_S) : W \subset g^{-1}Lg \}, \quad \mathcal{S}(L,W) = \bigcup_{M \in \mathcal{H}, MT(M) \subsetneq MT(L)} \mathcal{N}(M,W)$$

$$\mathcal{T}_L(W) = \pi_S(\mathcal{N}(L, W) - \mathcal{S}(L, W))$$

where $\pi_S: \mathbf{G}(\mathbb{Q}_S) \to \Gamma \backslash \mathbf{G}(\mathbb{Q}_S)$ denotes the canonical projection. Note that we have

$$\mathcal{N}(L, W) = \{ g \in \mathbf{G}(\mathbb{Q}_S) : W \subset g^{-1}MT(L)(\mathbb{Q}_S)g \},\$$

and for any $P, Q \in \mathcal{H}$ with MT(P) = MT(Q),

$$\mathcal{N}(P,W) = \mathcal{N}(Q,W), \quad \mathcal{S}(P,W) = \mathcal{S}(Q,W), \quad \mathcal{T}_P(W) = \mathcal{T}_Q(W).$$

By [GO11, Lemma 6.10], for any $P, Q \in \mathcal{H}$, either $\mathcal{T}_P(W) \cap \mathcal{T}_Q(W) = \emptyset$ or $\mathcal{T}_P(W) = \mathcal{T}_Q(W)$.

Let \mathcal{F}^* be the collection of Γ -conjugacy classes of Mumford-Tate subgroups of $L \in \mathcal{H}$. For each $[\mathbf{L}] \in \mathcal{F}^*$, choose one subgroup $L \in \mathcal{H}$ with $MT(L) = \mathbf{L}$, and the collection of such L's will be denoted by \mathcal{H}^* .

Let $L \in \mathcal{H}^*$. Let \mathfrak{g} and \mathfrak{l} denote the Lie algebras of \mathbf{G} and MT(L) respectively. For $d = \dim \mathfrak{l}$, consider the representation

$$\wedge^d \operatorname{Ad} : \mathbf{G} \to \mathbf{GL}(\mathbf{V}_L)$$

where $\mathbf{V}_L = \wedge^d \mathfrak{g}$. Let $\mathbf{V}_L(\mathbb{Q}_S) = \prod_{p \in S} \mathbf{V}_L(\mathbb{Q}_p)$ and $p_L \in \wedge^d \mathfrak{l}(\mathbb{Q}) \in \mathbf{V}_L(\mathbb{Q}_S)$. Then we have a map $\eta_L : \mathbf{G}(\mathbb{Q}_S) \to \mathbf{V}_L(\mathbb{Q}_S)$ defined by

$$\eta_L((g_p)_{p \in S}) = (\wedge^d \operatorname{Ad}(g_p)p_L)_{p \in S}.$$

Let $\Gamma_L := \{ \gamma \in \mathbf{G}(\mathcal{O}_S) : \gamma^{-1} MT(L) \gamma = MT(L) \}$ and denote by

$$\Gamma_L^0 := \{ \gamma \in \mathbf{G}(\mathcal{O}_S) : \eta_L(\gamma) = p_L \} = \{ \gamma \in \Gamma_L : \det(\mathrm{Ad}(\gamma)|_{\mathfrak{l}}) = 1 \}.$$

By definition, we have $\eta_L(\Gamma_L) \subset \mathcal{O}_S^{\times} \cdot p_L$ where \mathcal{O}_S^{\times} denotes the group of units in \mathcal{O}_S .

Recall the notion of $S(v_0)$ -small subsets of $\mathbf{V}_L(\mathbb{Q}_S)$. Fix $\delta > 0$ such that for any $w \in S$, if $\alpha \in \mathcal{O}_S^{\times}$ satisfies $\max_{v \in S \setminus \{w\}} |1 - \alpha|_v < \delta$, then $\alpha = 1$. Now let $v_0 \in S$. A subset $C = \prod_{v \in S} C_v \subset \mathbf{V}_L(\mathbb{Q}_S)$ is $S(v_0)$ -small if for any $v \in S \setminus \{v_0\}$ and $\alpha \in K_v^{\times}$, $\alpha C_v \cap C_v \neq \emptyset$ implies that $|1 - \alpha|_v < \delta$. Note that for any $\alpha \in \mathcal{O}_S^{\times}$ and any $S(v_0)$ -small subset C of $\mathbf{V}_L(\mathbb{Q}_S)$, $\alpha C \cap C \neq \emptyset$ implies $\alpha = 1$.

If $-p_L \notin \eta_L(\Gamma_L)$, we set $\overline{\mathbf{V}}_L(\mathbb{Q}_S) := \mathbf{V}_L(\mathbb{Q}_S)$; otherwise, set $\overline{\mathbf{V}}_L(\mathbb{Q}_S) := \mathbf{V}_L(\mathbb{Q}_S)/\{1, -1\}$. Denote by $\bar{\eta}_L$ the composition map of η_L with the quotient map $\mathbf{V}_L(\mathbb{Q}_S) \to \overline{\mathbf{V}}_L(\mathbb{Q}_S)$. Denote by A_L the Zariski closure of $\bar{\eta}_L(\mathcal{N}(L,W))$ in $\overline{\mathbf{V}}_L(\mathbb{Q}_S)$. Then $\bar{\eta}_L^{-1}(A_L) = \mathcal{N}(L,W)$. For D a compact $S(v_0)$ -small subset of A_L for some $v_0 \in S$, we define

$$\mathcal{S}(D) = \{ g \in \bar{\eta}_L^{-1}(D) : \gamma g \in \bar{\eta}_L^{-1}(D) \text{ for some } \gamma \in \Gamma - \Gamma_L \}.$$

We finish this subsection by proving the following lemma.

Lemma 5.4. Let \mathbf{M} be a connected \mathbb{Q} -group in class \mathcal{F} , and $M \in \mathcal{H}$ a subgroup of finite index in $\mathbf{M}(\mathbb{Q}_S)$. Suppose that $\mathbf{M}(\mathbb{R})^0$ is generated by one-parameter unipotent subgroups. Then

$$\overline{\Gamma\backslash\Gamma(\mathbf{M}(\mathbb{R})^0\times\{e_f\})}=\Gamma\backslash\Gamma M.$$

Proof. By [Tom00, Theorem 2], we have $\overline{\Gamma\backslash\Gamma\mathbf{M}(\mathbb{R})^0} = \Gamma\backslash\Gamma F$ for some $F \in \mathcal{H}$ and $\mathbf{M}(\mathbb{R})^0 \subset F$. This implies that $MT(M) \subset MT(F)$. On the other hand, we also have

$$\Gamma \backslash \Gamma F = \overline{\Gamma \backslash \Gamma \mathbf{M}(\mathbb{R})^0} \subset \Gamma \backslash \Gamma M.$$

So by [GO11, Lemma 6.7], $MT(F) \subset MT(M)$. This implies that $MT(L) = MT(M) = \mathbf{M}$. Hence by [GO11, Lemma 6.7], $\Gamma \cap F = \Gamma \cap M$. This completes the proof of the lemma.

5.3. An Auxiliary Proposition. Let \mathbf{H} be as in Theorem 5.1. Let $\rho_i : \mathbf{H} \to \mathbf{G}(i \in \mathbb{N})$ be a sequence of \mathbb{Q} -homomorphisms defined by $\rho_i(x) = \gamma_i x \gamma_i^{-1}$ ($\gamma_i \in \Gamma$) where Γ is a subgroup of finite index in Γ_S , such that there exists no proper \mathbb{Q} -subgroup \mathbf{M} containing infinitely many $\rho_i(\mathbf{H})$. Now for the group H, suppose that the sequence of subsets $\{\rho_i(H)\}$ in $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$ is bounded, and the norms $\|\operatorname{Ad} \gamma_i\|$ of the operators $\operatorname{Ad} \gamma_i : \operatorname{Lie}(\mathbf{T}(\mathbb{R})) \to \operatorname{Lie}(\mathbf{G}(\mathbb{R}))$ are bounded. Then by Proposition 4.4 and Proposition 4.5, for each $p \in S$, we can find a bounded sequence $\gamma_{i,p} \in \mathbf{G}(\mathbb{Q}_p)$ and a sequence $t_{i,p} \in Z_{\mathbf{G}}(\mathbf{H})(\mathbb{Q}_p)$ such that $\gamma_i = \gamma_{i,p}t_{i,p}$. Therefore,

$$\pi_S(\rho_i(H)) = \Gamma \backslash \Gamma H \gamma_i^{-1} = \Gamma \backslash \Gamma(t_{i,p}^{-1})_{p \in S} H(\gamma_{i,p}^{-1})_{p \in S}.$$

Since $\mu_{\Gamma \setminus \Gamma \rho_i(H)}$ converges to a probability measure, the sequence $\{\Gamma(t_{i,p})_{p \in S}\}$ is bounded in the space $\Gamma \setminus \Gamma Z_{\mathbf{G}}(\mathbf{H})(\mathbb{Q}_S)$. Let

$$\Gamma(t_{i,p}^{-1})_{p \in S} = \Gamma(a_{i,p})_{p \in S}$$

where $(a_{i,p})_{p\in S}$ is a bounded sequence in $Z_{\mathbf{G}}(\mathbf{H})(\mathbb{Q}_S)$. Then

$$\pi_S(\rho_i(H)) = \Gamma \backslash \Gamma H(a_{i,p} \gamma_{i,p}^{-1}).$$

By passing to a subsequence, the limiting measure of $\mu_{\Gamma \setminus \Gamma \rho_i(H)}$ is a translate of the periodic measure $\mu_{\Gamma \setminus \Gamma H}$. Theorem 5.1 then follows in this case.

Now we assume that either the sequence $\{\rho_i(H)\}$ is unbounded in $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$ or the norms $\|\operatorname{Ad}\gamma_i\|$ of the operators $\operatorname{Ad}\gamma_i: \operatorname{Lie}(\mathbf{T}(\mathbb{R})) \to \operatorname{Lie}(\mathbf{G}(\mathbb{R}))$ are unbounded. Then by Proposition 4.3 and Proposition 4.6, after passing to a subsequence, the limiting probability measure μ of $\mu_{\Gamma \setminus \Gamma \rho_i(H)}$ is invariant under a one parameter unipotent subgroup W.

We continue the notation and the discussion in the previous subsection. According to [GO11, Theorem 6.11], there is some subgroup $L \in \mathcal{H}^*$ such that $\mu(\mathcal{T}_L(W)) > 0$ and $\mu(\mathcal{S}(L,W)) = 0$. Then one can find a compact subset $C^* \subset \mathcal{T}_L(W)$ such that $\mu(C^*) = \alpha > 0$. Since $\bar{\eta}_L$ never takes the zero vector in every place $p \in S$, we may take the compact subset C^* so that the subset

$$C_L := \bar{\eta}_L(C^*) \subset A_L$$

is contained in a subset $\prod_{v \in S} C_v$ of $\mathbf{V}_L(\mathbb{Q}_S)$ which is $S(v_0)$ -small for every $v_0 \in S$.

Let Ω be a fundamental domain of $\Gamma \backslash \Gamma H$ in H, and μ_H the Haar measure on H. Let Ω_c be a compact subset in Ω such that

$$\mu_H(\Omega - \Omega_c) \le \epsilon_0 \mu_H(\Omega)$$

where $\epsilon_0 < \alpha/2$. Let $\{B_k\}_{k=1}^N$ be a cover of Ω_c with open boxes of diameter δ_0 , where δ_0 is as in Proposition 2.2. Moreover, we assume that for any B_k $(1 \le k \le N)$, one can find an element $g_k \in \mathbf{G}(\mathbb{Q}_S)$ and an open box U_k of radius at most $\delta_0 > 0$ around 0 in $\mathrm{Lie}(\mathbf{H}(\mathbb{Q}_S)) = \prod_{p \in S} \mathrm{Lie}(\mathbf{H}(\mathbb{Q}_p))$ such that $\exp_S(U_k) = B_k g_k^{-1}$. The orbit $\Gamma \backslash \Gamma \rho_i(H)$ is then equal to

$$\Gamma \backslash \Gamma \rho_i(H) = \bigcup_{k=1}^N \Gamma \backslash \Gamma \rho_i(B_k) \cup \Gamma \backslash \Gamma \rho_i(\Omega - \Omega_c).$$

Note that N depends only on ϵ_0 and Ω_c . By Proposition 2.2, [KT07, Corollary 3.3] and the discussion in §3, for each ρ_i and B_k , one can find a bijective (C, α) -good map

$$\phi_i: U_k \to \rho_i(B_k)$$

where C and α are uniform for all $i \in \mathbb{N}$ and $1 \le k \le N$.

Write A_L as $\prod_{p \in S} V(f_1^{(p)}, f_2^{(p)}, \dots, f_n^{(p)})$, where $f_i^{(p)}$ $(1 \leq i \leq n)$ are polynomials on $\mathbf{V}_L(\mathbb{Q}_p)$ $(p \in S)$ and $V(f_1^{(p)}, f_2^{(p)}, \dots, f_n^{(p)})$ is the zero set of $f_i^{(p)}$ (i = 1, 2, ..., n). We define

$$W_p(R,\alpha) = \{ x \in \mathbf{V}_L(\mathbb{Q}_p) : ||x||_p \le R, \quad |g_i^{(p)}(x)|_p \le \alpha \}.$$

Define $R_0 := \sup_{x \in C_L} d(x, 0)$. Then $\prod_{p \in S} W_p(R_0, \alpha)$ is a compact neighborhood of C_L in $\mathbf{V}_L(\mathbb{Q}_S)$.

Proposition 5.5. Let $C_L = \bar{\eta}_L(C^*) \subset A_L$ and ϕ_i defined as above. Fix B_k for some $1 \leq k \leq N$ $(\phi_i(U_k) = \rho_i(B_k))$. Then there exists a closed subset S in $\pi_S(S(L, W))$ satisfying the following property: for a given compact subset K in $X_S - S$, one of the following holds:

(1) For any $\epsilon > 0$ there exists a neighborhood Ψ of C_L such that

$$\mu(\lbrace x \in U_k : \pi_S(\phi_i(x)) \in \mathcal{K} \cap \pi_S(\overline{\eta}_L^{-1}(\Psi))\rbrace) \le \epsilon \mu(U_k)$$

for infinitely many $i \in \mathbb{N}$.

(2) There exists R > 0 such that for any $\alpha > 0$ and for all large enough $i \in \mathbb{N}$, one can find $w_i \in \Gamma \bar{p}_L$ with

$$\rho_i(B_k)w_i \subset \prod_{p \in S} W_p(R, \alpha).$$

Proof. Suppose that the second claim does not hold. Let $\epsilon > 0$, and take R > 0 sufficiently large such that $R_0/R \le \epsilon$. Then we can find $\alpha > 0$ such that for infinitely many $i \in \mathbb{N}$

$$\rho_i(B_k)w \not\subset \prod_{p\in S} W_p(R,\alpha)$$

for any $w \in \Gamma \bar{p}_L$. In the sequel, we fix such an R.

For any $p \in S$ define

$$D(p) = \prod_{v \neq p} C_v \times W_p(R, \alpha).$$

Then D(p) is a compact S(p)-small subset which contains $\prod_{v \in S} C_v$. Let $S := \bigcup_{p \in S} S(D(p))$. For each $p \in S$ we take a small neighborhood of D(p)

$$\Psi(p) = \prod_{v \in S} \Psi_v(p)$$

as in [GO11, Proposition 6.16], and define

$$\Psi = (\prod_{p \in S} W_p(R_0, \epsilon \alpha)) \cap \bigcap_{p \in S} \Psi(p).$$

We show that the subset Ψ satisfies the first claim in the proposition.

Fix a compact \mathcal{K} in $X_S - \mathcal{S}$. For $p \in S$, denote by

$$\mathcal{F}_p = \{ w \in \Gamma \bar{p}_L : \rho_i(B_k) w \not\subset \prod_{v \in S, v \neq p} \mathbf{V}_L(\mathbb{Q}_v) \times W_p(R, \alpha) \}.$$

By our assumption, $\bigcup_{p\in S} \mathcal{F}_p = \Gamma \bar{p}_L$ and we have

$$\mu(\{x \in U_k : \pi_S(\phi_i(x)) \in \mathcal{K} \cap \pi_S(\overline{\eta}_L^{-1}(\Psi))\})$$

$$=\mu(\{x \in U_k : \pi_S(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)\Gamma\bar{p}_L \cap \Psi \neq \emptyset\})$$

$$\leq \sum_{p \in S} \mu(\{x \in U_k : \pi_S(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)\mathcal{F}_p \cap \Psi \neq \emptyset\}).$$

In the sequal, we will estimate

$$\mu(\lbrace x \in U_k : \pi_S(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)\mathcal{F}_p \cap \Psi \neq \emptyset \rbrace)$$

for each $p \in S$.

For convenience, for any $v \in S$, we write π_v for the natural projection $\text{Lie}(\mathbf{H}(\mathbb{Q}_S)) \to \text{Lie}(\mathbf{H}(\mathbb{Q}_v))$, and π^v for the natural projection from $\text{Lie}(\mathbf{H}(\mathbb{Q}_S))$ to $\text{Lie}(\mathbf{H}(\mathbb{Q}_{S\setminus\{v\}}))$. We denote by μ_v the Haar measure on $\text{Lie}(\mathbf{H}(\mathbb{Q}_v))$, and μ_S the Haar measure on $\text{Lie}(\mathbf{H}(\mathbb{Q}_S))$.

Case 1: $p \in S_f$. Define

$$\mathcal{S}_p = \{ y \in \text{Lie}(\mathbf{H}(\mathbb{Q}_{S \setminus \{p\}})) : \exists x \in U_k \text{ such that } \pi_S(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)\mathcal{F}_p \cap \Psi \neq \emptyset, \ \pi^p(x) = y \}.$$

For $y \in \mathcal{S}_p$ and $w \in \mathcal{F}_p$, let

$$I_w(y) = \{ s \in \text{Lie}(\mathbf{H}(\mathbb{Q}_p)) : \exists x \in U_k \text{ such that } \pi_S(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)w \in \Psi, \pi^p(x) = y, \pi_p(x) = s \}.$$

For any $s \in I_w(y)$, we denote by $B_{s,w}(y)$ the maximal ball in $U_k \cap \text{Lie}(\mathbf{H}(\mathbb{Q}_p))$ containing s such that

- (1) $\phi_i(B_{s,w}(y))w \subset \overline{\Psi_p(p)}$.
- (2) There exists $s' \in B_{s,w}(y)$ with $\phi_i(s')w \in \overline{\Psi_p(p)} \setminus \Psi_p(p)$.

By definition of \mathcal{F}_p , such $B_{s,w}(y)$ exists and by [GO11, Proposition 6.16]

$$\{B_{s,w}(y): w \in \mathcal{F}_p, s \in I_w(y)\}$$

is a collection of disjoint balls in $U_k \cap \text{Lie}(\mathbf{H}(\mathbb{Q}_p))$, which we rewrite as $\{B_i(w,y)\}_{i\in\mathbb{N}}$. Now since $\phi_i|_{U_k}$ is (C,α) -good, for $y\in\mathcal{S}_p$ and $w\in\mathcal{F}_p$

$$\mu_p(I_w(y))$$

$$\leq \sum_{i=1}^{\infty} \mu_p(\{s' \in B_i(w, y) : \exists x \in U_k \text{ such that } \phi_i(x)w \in \Psi, \pi_p(x) = s', \pi^p(x) = y\})$$

$$\leq \sum_{i=1}^{\infty} \epsilon \mu_p(B_i(w, y)).$$

Then we have

$$\mu_S(\{x \in U_k : \pi(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)\mathcal{F}_p \cap \Psi \neq \emptyset\})$$

$$\leq \sum_{w \in \mathcal{F}_p} \mu_S(\{x \in U_k : \pi(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)w \in \Psi\})$$

$$= \sum_{w \in \mathcal{F}_p} \int_{y \in \mathcal{S}_p} \mu_p(I_w(y)) dy \leq \int_{y \in \mathcal{S}_p} \sum_{w \in \mathcal{F}_p} \sum_{i=1}^{\infty} \epsilon \mu_p(B_i(w, y)) dy$$

$$\leq \int_{y \in \mathcal{S}_p} \epsilon \mu_p(U_k \cap \text{Lie}(\mathbf{H}(\mathbb{Q}_p))) dy \leq \epsilon \mu_S(U_k).$$

Case 2: $p = \infty$. Define

 $\mathcal{S}_{\infty} = \{ y \in \text{Lie}(\mathbf{H}(\mathbb{Q}_{S \setminus \{\infty\}})) : \exists x \in U_k, \pi_S(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)\mathcal{F}_{\infty} \cap \Psi \neq \emptyset \text{ such that } \pi^{\infty}(x) = y \}.$ For $y \in \mathcal{S}_{\infty}$ let

$$I(y) = \{ s \in \text{Lie}(\mathbf{H}(\mathbb{R})) : \exists x \in U_k, \ \pi(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x) \mathcal{F}_{\infty} \cap \Psi \neq \emptyset, \pi_{\infty}(x) = s, \pi^{\infty}(x) = y \}.$$

Now fix $y \in \mathcal{S}_{\infty}$ and choose $s_0 \in I(y)$. By [GO11, Proposition 6.16], one can find $w_0 \in \Gamma \bar{p}_L$ such that

$$\pi_S(\phi_i(s_0, y)) \in \mathcal{K}, \ \phi_i(s_0, y)w_0 \in \Psi$$

and for other $w \in \Gamma \bar{p}_L - \{w_0\}$, we have $\phi_i(s_0, y)w \notin \Psi(\infty)$. We split I(y) into two subsets

$$I(y, \mathcal{F}_{\infty} - \{w_0\}) := \{s \in \text{Lie}(\mathbf{H}(\mathbb{R})) : (s, y) \in U_k, \ \pi_S(\phi_i(s, y)) \in \mathcal{K}, \ \phi_i(s, y)(\mathcal{F}_{\infty} - \{w_0\}) \cap \Psi \neq \emptyset\},\$$

and

$$I(y, w_0) := \{ s \in \text{Lie}(\mathbf{H}(\mathbb{R})) : (s, y) \in U_k, \ \pi_S(\phi_i(s, y)) \in \mathcal{K}, \ \phi_i(s, y) w_0 \in \Psi \}.$$

We first consider $I(y, \mathcal{F}_{\infty} - \{w_0\})$. Let $\vec{v} \in \text{Lie}(\mathbf{H}(\mathbb{R}))$ and $||\vec{v}|| = 1$. Define

$$I_{s_0}(\vec{v}) := \{ x \in [0, \delta_0] : s_0 + x\vec{v} \in U_k \cap \text{Lie}(\mathbf{H}(\mathbb{R})) \}.$$

For $w \in \Gamma \bar{p}_L - \{w_0\}$, define

$$I_{s_0,w}(\vec{v}) := \{ x \in I_{s_0}(\vec{v}) : \pi_S(\phi_i(s_0 + x\vec{v}, y)) \in \mathcal{K}, \ \phi_i(s_0 + x\vec{v}, y)w \in \Psi \}.$$

For any $x \in I_{s_0,w}(\vec{v})$, we denote by $B_{x,w}(\vec{v})$ the maximal interval in $I_{s_0}(\vec{v})$ containing x such that

- (1) $\phi_i(s_0 + B_{x,w}(\vec{v})\vec{v}, y)w \subset \overline{\Psi(\infty)}$.
- (2) There exists $x' \in B_{x,w}(\vec{v})$ with $\phi_i(s_0 + x'\vec{v}, y)w \in \overline{\Psi(\infty)} \setminus \Psi(\infty)$.

By definition of \mathcal{F}_{∞} and our choice of s_0 and w_0 , such an interval $B_{x,w}(\vec{v})$ exists, and by [GO11, Proposition 6.16]

$$\{B_{x,w}(\vec{v}): w \in \mathcal{F}_{\infty} - \{w_0\}, x \in I_{s_0,w}(\vec{v})\}$$

is a collection of intervals which cover $I_{s_0}(\vec{v})$ at most twice. We rewrite this collection as $\{B_i(w, \vec{v})\}_{i \in \mathbb{N}}$. Now since $\phi_i|_{U_k}$ is (C, α) -good, for any $y \in \mathcal{S}_{\infty}$ and $\vec{v} \in S_1$

$$\sum_{w \in \mathcal{F}_{\infty} - \{w_{0}\}} \mu(I_{s_{0}, w}(\vec{v})) \leq \sum_{w \in \mathcal{F}_{\infty} - \{w\}} \sum_{i=1}^{\infty} \mu(\{x' \in B_{i}(w, \vec{v}) : \phi_{i}(s_{0} + x'\vec{v}, y)w \in \Psi\})$$

$$\leq \sum_{w \in \mathcal{F}_{\infty} - \{w\}} \sum_{i=1}^{\infty} \epsilon \mu(B_{i}(w, \vec{v})) \leq 2\epsilon \mu(I_{s_{0}}(\vec{v})) \leq 2\epsilon \delta_{0}.$$

We have

$$\mu(I(y, \mathcal{F} - \{w_0\})) \leq \sum_{w \in \mathcal{F}_{\infty} - \{w_0\}} \int_{S^1} \int_{r \in [0, \delta_0]} \chi_{I_{s_0, w}(\vec{v})}(r) r^{l-1} dr d\vec{v}$$

$$\leq \int_{S^1} \sum_{w \in \mathcal{F}_{\infty} - \{w_0\}} \mu(I_{s_0, w}(\vec{v})) \delta_0^{l-1} d\vec{v} \leq 2\epsilon \delta_0^l = 2\epsilon \mu(U_k \cap \text{Lie}(\mathbf{H}(\mathbb{R})))$$

where $l = \dim \mathbf{H}$.

For $I(y, w_0)$, by our assumption on \mathcal{F}_{∞} , one can find $s_1 \in U_k$ such that $\phi_i(s_1, y)w_0\bar{p}_L \notin \Psi(\infty)$. We can then repeat the arguments above by replacing s_0 and $\mathcal{F}_{\infty} - \{w_0\}$ with s_1 and $\{w_0\}$, and obtain $\mu(I(y, w_0)) \leq 2\epsilon \mu(B_k)$. Therefore, we have

$$\mu(\{x \in U_k : \pi(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)\mathcal{F}_{\infty} \cap \Psi \neq \emptyset\})$$

$$\leq \mu(\{x \in U_k : \pi(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)(\mathcal{F}_{\infty} - \{w_0\}) \cap \Psi \neq \emptyset\})$$

$$+ \mu(\{x \in U_k : \pi(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)w_0 \in \Psi\})$$

$$\leq \int_{y \in \mathcal{S}_{\infty}} \mu_{\infty}(I(y, \mathcal{F} - \{w_0\}))dy + \int_{y \in \mathcal{S}_{\infty}} \mu_{\infty}(I(y, w_0))dy$$

$$\leq \int_{y \in \mathcal{S}_{\infty}} 4\epsilon \mu_{\infty}(U_k \cap \text{Lie}(\mathbf{H}(\mathbb{R})))dy \leq 4\epsilon \mu_S(U_k).$$

Combining Case 1 and Case 2, we have

$$\mu(\{x \in U_k : \pi(\phi_i(x)) \in \mathcal{K} \cap \pi(\overline{\eta}_L^{-1}(\Psi))\})$$

$$\leq \sum_{p \in S} \mu(\{x \in U_k : \pi(\phi_i(x)) \in \mathcal{K}, \ \phi_i(x)\mathcal{F}_p \cap \Psi \neq \emptyset\})$$

$$\leq \sum_{p \in S_f} \epsilon \mu(U_k) + 4\epsilon \mu(U_k) \leq 4\epsilon |S| \mu(U_k).$$

We complete the proof of the proposition by replacing $4\epsilon |S|$ with ϵ .

5.4. **Proof of Theorem 5.1.** We will follow the arguments in [EMS96, §4] and [EMS98]. Let μ be the limiting measure of $\mu_{\Gamma \setminus \Gamma \rho_i(H)}$. We continue the assumption that either the sequence $\{\rho_i(H)\}$ is unbounded in $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$ or the norms of the operators $\operatorname{Ad} \gamma_i : \operatorname{Lie}(\mathbf{T}(\mathbb{R})) \to \operatorname{Lie}(\mathbf{G}(\mathbb{R}))$ are unbound. Then μ is invariant under a one-parameter unipotent subgroup W. Let L be a subgroup in \mathcal{H}^* such that $\mu(\mathcal{T}_L(W)) > 0$ and $\mu(\mathcal{S}(L,W)) = 0$. Then one can find a compact subset $C \subset \mathcal{T}_L(W)$ such that $\mu(C) = \alpha > 0$, and $C_L := \bar{\eta}_L(C) \subset A_L$ is contained in a subset $\prod_{v \in S} C_v$ of $\mathbf{V}_L(\mathbb{Q}_S)$ which is $S(v_0)$ -small for every $v_0 \in S$.

Let K be a compact neighborhood of C. The positivity of the measure $\mu(C)$ implies that the first claim in Proposition 5.5 does not hold for some B_k $(1 \le k \le N)$. In other words, there exist B_k , $\alpha_i \to 0$ and $w_i \in \Gamma$ such that

$$\rho_i(B_k)w_i\bar{p}_L\subset\Phi_i$$

where $\Phi_i = \prod_{p \in S} W_p(R, \alpha_i)$ is a decreasing sequence of relatively compact subsets in \bar{V}_L with $\bigcap_{i=1}^{\infty} \Phi_i \subset A_L$.

Let S be the collection of all \mathbb{Q}_{S} -valued functions on B_k of the form

$$\omega = (\omega_p)_{p \in S} \mapsto (f_p(\rho_i(\omega_p)wp_L))_{p \in S}, \quad \omega \in B_k$$

where $w \in \Gamma$, $i \in \mathbb{N}$ and f_p is a linear functional on $\mathbf{V}_L(\mathbb{Q}_p)$ for each $p \in S$. By Lemma 5.2, \mathcal{S} spans a finitely generated \mathbb{Q}_S -module in the space of functions on B_k . Therefore, by [EMS97, Lemma 4.1], there exists a finite set $\Sigma \subset B_k$ such that for any \mathbb{Q}_S -valued function $\{\phi_p\}_{p \in S} \in \mathcal{S}$, if $\{\phi_p : p \in S\}$ vanishes on Σ , then $\{\phi_p\}_{p \in S} = 0$ on B_k . Note that since ϕ_p $(p \in S)$ are regular and B_k are Zariski dense in $\mathbf{H}(\mathbb{Q}_S)$, this indeed implies $\{\phi_p\}_{p \in S} = 0$ on $\mathbf{H}(\mathbb{Q}_S)$. By Lemma 5.3, $\mathbf{H}(\mathbb{Q}) \cap B_k$ is dense in B_k . So we may assume that $\Sigma \subset \mathbf{H}(\mathbb{Q}_S) \cap B_k$.

Now for any $s \in \Sigma$ and $i \in \mathbb{N}$, $\rho_i(s)w_ip_L \subset \Phi_1$. One can find $k \in \mathbb{N}$ such that

$$\{\rho_i(s): i \in \mathbb{N}, s \in \Sigma\} \subset \mathbf{G}(\frac{1}{k}\mathcal{O}_S).$$

Since $\mathbf{G}(\frac{1}{k}\mathcal{O}_S)w_ip_L$ is a discrete subset of $\bar{\mathbf{V}}_L(\mathbb{Q}_S)$ and Φ_1 is bounded, by passing to a subsequence, we have $\rho_i(s)w_ip_L = \rho_1(s)w_1p_L$ for all $i \in \mathbb{N}$ and $s \in \Sigma$. Therefore, by our choice of Σ , for all $\omega \in \mathbf{H}(\mathbb{Q}_S)$, we have $\rho_i(\omega)w_ip_L = \rho_1(\omega)w_1p_L$. Now set $\omega = e_S$, we have $w_ip_L = w_1p_L$. Thus

$$\rho_i(\omega)w_1p_L = \rho_1(\omega)w_1p_L$$

for all $i \in \mathbb{N}$ and $\omega \in \mathbf{H}(\mathbb{Q}_S)$. Consequently, $w_1^{-1}\rho_1(\omega)^{-1}\rho_i(\omega)w_1 \in N(MT(L))$ for all $i \in \mathbb{N}$ and $\omega \in \mathbf{H}(\mathbb{Q})$, where N(MT(L)) is the normalizer of MT(L) in \mathbf{G} .

Without loss of generality, in the following, one may assume $w_1 = e_S$. Therefore

$$\rho_1(\omega)^{-1}\rho_i(\omega) \in N(MT(L)), \quad \omega \in \mathbf{H}(\mathbb{Q}).$$

We prove that N(MT(L)) equals **G**. Let **F** be the smallest algebraic subgroup of **G** containing $\rho_1(\omega)^{-1}\rho_i(\omega)$ ($\omega \in \mathbf{H}(\mathbb{Q})$). Then **F** is defined over \mathbb{Q} . Let $\omega_1, \omega_2 \in \mathbf{H}(\mathbb{Q})$. We compute

$$\rho_1(\omega_1)^{-1}(\rho_1(\omega_2)^{-1}\rho_i(\omega_2))\rho_1(\omega_1) = (\rho_1(\omega_2\omega_1)^{-1}\rho_i(\omega_2\omega_1))(\rho_1(\omega_1)^{-1}\rho_i(\omega_1))^{-1} \in \mathbf{F}(\mathbb{Q}).$$

This implies that $\rho_1(\mathbf{H}) \subset N(\mathbf{F})$ where $N(\mathbf{F})$ is the normalizer of \mathbf{F} , and $\rho_1(\mathbf{H})\mathbf{F}$ is a \mathbb{Q} -subgroup. Since $\rho_i(\mathbf{H}) \subset \rho_1(\mathbf{H})\mathbf{F}$ for all $i \in \mathbb{N}$ and \mathbf{G} is the smallest \mathbb{Q} -group containing $\rho_i(\mathbf{H})$ $(i \in \mathbb{N})$, we have $\mathbf{G} = \rho_1(\mathbf{H})\mathbf{F}$. Hence \mathbf{F} is normal in \mathbf{G} .

By the definition of A_L , we know that $\rho_i(\omega) \in \mathcal{N}(L, W)$ ($\omega \in \mathbf{H}(\mathbb{Q}_S)$). We have

$$\mathbf{G}(\mathbb{Q}_S) = \mathbf{F}(\mathbb{Q}_S) \rho_1(\mathbf{H}(\mathbb{Q}_S)) \subset \mathbf{F}(\mathbb{Q}_S) \mathcal{N}(L, W) = \mathcal{N}(L, W).$$

Let P be the closed subgroup generated by gWg^{-1} $(g \in \mathbf{G}(\mathbb{Q}_S))$. Then we have

$$\overline{\Gamma \backslash \Gamma P} = \Gamma \backslash \Gamma L.$$

For any $\omega \in \mathbf{H}(\mathbb{Q})$, $\rho_1(\omega)$ normalizes P and hence

$$P \subset \rho_1(\omega)L\rho_1(\omega)^{-1}, \quad \overline{\Gamma\backslash\Gamma P} \subset \Gamma\backslash\Gamma(\rho_1(\omega)L\rho_1(\omega)^{-1}).$$

This implies that $\rho_1(\omega)MT(L)\rho_1(\omega)^{-1}=MT(L)$ and hence MT(L) is normal in **G**.

Since $\mu(\pi_S(S(L, W))) = 0$, by [GO11, Theorem 6.11], every W-ergodic component of μ is L-invariant, and hence so is μ . Since **G** is \mathbb{Q} -semisimple, MT(L) is an almost direct product of

 \mathbb{Q} -factors of \mathbf{G} . Since L is a subgroup of finite index in $MT(L)(\mathbb{Q}_S)$, we may choose L_0 a subgroup of finite index in $MT(L)(\mathbb{Q}_S)$ such that L_0 is normal in $\mathbf{G}(\mathbb{Q}_S)$ and $L_0 \subset L$. Hence μ is L_0 -invariant.

Now consider the quotient space $\overline{X} = (\Gamma/(\Gamma \cap L_0)) \setminus (\mathbf{G}(\mathbb{Q}_S)/L_0)$, and denote by $\overline{\mu}_{\Gamma \setminus \Gamma \rho_i(H)}$ the pushforward of the measure $\mu_{\Gamma \setminus \Gamma \rho_i(H)}$ via the map $\Gamma \setminus \mathbf{G}(\mathbb{Q}_S) \to \overline{X}$. Since **H** is connected, we have

$$\rho_i(H)/(L_0 \cap \rho_i(H)) \subset \rho_i(\mathbf{H}(\mathbb{Q}_S))/(\rho_i(\mathbf{H}(\mathbb{Q}_S)) \cap L_0) \subset \mathbf{G}(\mathbb{Q}_S)/MT(L)(\mathbb{Q}_S)$$

and the measure $\overline{\mu}_{\Gamma \backslash \Gamma \rho_i(H)}$ is indeed supported on the space $(\Gamma/(\Gamma \cap L_0))\backslash (\mathbf{G}(\mathbb{Q}_S)/MT(L)(\mathbb{Q}_S))$. We complete the proof of Theorem 5.1 by induction on dim \mathbf{G} .

6. Correspondence between $\Gamma_{\infty}\backslash \mathbf{G}(\mathbb{R})$ and $\Gamma_{S}\backslash \mathbf{G}(\mathbb{Q}_{S})$: Proof of Theorem 1.8

In this section, we prove Theorem 6.4 which is a sharpening of Theorem 1.7 in the sense that it gives explicit algebraic conditions on the deforming sequence that ensures equidistribution in the ambient space. This is done in a special case where \mathbf{G} an \mathbb{R} -split semisimple group defined over \mathbb{Q} and \mathbf{H} a maximal \mathbb{R} -split and \mathbb{Q} -anisotropic torus in \mathbf{G} . As a corollary of Theorem 6.4, we will deduce Theorem 1.8.

6.1. Theorem 1.7 for an \mathbb{R} -split semisimple group G and H a maximal \mathbb{R} -split and \mathbb{Q} -anisotropic torus. In the rest of this section, we will assume that G is an \mathbb{R} -split semisimple group defined over \mathbb{Q} and T is a maximal \mathbb{R} -split and \mathbb{Q} -anisotropic torus in G.

Lemma 6.1. Let G be an \mathbb{R} -split semisimple group and S a maximal \mathbb{R} -split and \mathbb{Q} -anisotropic torus in G. Let F be a connected \mathbb{Q} -subgroup of G containing S. Then F is reductive and has no non-trivial \mathbb{Q} -characters.

Proof. By [EMS97, Lemma 5.1], **F** is reductive. Let χ be a \mathbb{Q} -character of **F**. Then $\chi = 1$ on **S**, as **S** is \mathbb{Q} -anisotropic. Moreover, $\chi = 1$ on any one parameter unipotent subgroup. Since **S** is maximal, **F** is generated by **S** and unipotent subgroups. Hence $\chi = 1$ on **F**.

Lemma 6.2. Let G be an \mathbb{R} -split semisimple group and S a maximal \mathbb{R} -split and \mathbb{Q} -anisotropic torus in G. Let F be a connected algebraic subgroup of G containing S. Let W be a set of representatives of the Weyl group of $S(\mathbb{R})$ in $G(\mathbb{R})$. Suppose that for some $\gamma \in G(\mathbb{R})$, $\gamma S(\mathbb{R}) \gamma^{-1} \subset F(\mathbb{R})$. Then there is $w \in W$ such that $\gamma w \in F(\mathbb{R})$.

Proof. By [EMS97, Lemma 5.1], \mathbf{F} is a reductive subgroup of \mathbf{G} . Now both $\mathbf{S}(\mathbb{R})$ and $\gamma \mathbf{S}(\mathbb{R}) \gamma^{-1}$ are maximal \mathbb{R} -split tori in $\mathbf{F}(\mathbb{R})$, so there is $x \in \mathbf{F}(\mathbb{R})$ such that $\gamma \mathbf{S}(\mathbb{R}) \gamma^{-1} = x \mathbf{S}(\mathbb{R}) x^{-1}$. This implies that $\gamma^{-1}x$ is in the normalizer of $\mathbf{S}(\mathbb{R})$ in $\mathbf{G}(\mathbb{R})$. Since the Weyl group of $\mathbf{S}(\mathbb{R})$ is the normalizer of $\mathbf{S}(\mathbb{R})$ modulo $\mathbf{S}(\mathbb{R})$, one can find $w \in W$ and $y \in \mathbf{S}(\mathbb{R})$ so that $\gamma^{-1}x = wy$. Hence $\gamma w = xy^{-1} \in \mathbf{F}(\mathbb{R})$. This completes the proof of the lemma.

Let Λ be a subgroup of finite index in $\mathbf{T}(\mathbb{Z}) \cap \mathbf{T}(\mathbb{R})^0$. We think of Λ and $\mathbf{T}(\mathbb{Z}) \cap \mathbf{T}(\mathbb{R})^0$ as subgroups of $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$ by diagonal embedding. Since the closure of $\mathbf{T}(\mathbb{Z}) \cap \mathbf{T}(\mathbb{R})^0$ is compact in $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$, we can choose Λ so that it is contained in a small neighborhood of identity in

 $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$ where the exponential map of p-adic Lie groups are defined for $p \in S_f$. We denote by F the closure of Λ in $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$ and H the group generated by $\Gamma_S \cap \mathbf{T}(\mathbb{Q}_S)$ and $\mathbf{T}(\mathbb{R}) \times F$. By strong approximation, H is a subgroup of finite index in $\mathbf{T}(\mathbb{Q}_S)$ and

$$\Gamma_S \backslash \Gamma_S H = \Gamma_S \backslash \Gamma_S (\mathbf{T}(\mathbb{R}) \times F).$$

By Theorem 1.7, one can conclude that for any $\{g_i\} \subset \mathbf{G}(\mathbb{Q}_S)$, the sequence $(g_i)^*\mu_{\Gamma_S\backslash\Gamma_S H}$ has a subsequence converging to an algebraic probability measure μ on $\Gamma_S\backslash\mathbf{G}(\mathbb{Q}_S)$. In order to prove Theorem 1.8, we have to review the arguments in §5 and analyze the limiting algebraic measure μ .

Lemma 6.3. The closure $\overline{\Gamma_S \backslash \Gamma_S \mathbf{T}(\mathbb{R})^0}$ of $\Gamma_S \backslash \Gamma_S \mathbf{T}(\mathbb{R})^0$ in $\Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$ is equal to $\Gamma_S \backslash \Gamma_S (\mathbf{T}(\mathbb{R})^0 \times F)$.

Proof. By diagonal embedding, we know that Λ is a lattice in $(\mathbf{T}(\mathbb{R})^0 \times F)$ and $\Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F)$ is closed in $\Gamma_S \setminus \mathbf{G}(\mathbb{Q}_S)$. Now let Ω be a fundamental domain of $\Lambda \setminus \mathbf{T}(\mathbb{R})^0$. Then for any $x \in \mathbf{T}(\mathbb{R})^0$, one can find $x' \in \Omega$ and $\gamma \in \Lambda$ such that $x = \gamma^{-1}x'$, and

$$\Gamma_S(x, e_f) = \Gamma_S(\gamma^{-1}x', e_f) = \Gamma_S\gamma(\gamma^{-1}x', e_f) = \Gamma_S(x', \gamma).$$

Note that x' and γ vary in $\mathbf{T}(\mathbb{R})^0$ and Λ respectively, and Λ is dense in F. So we have

$$\overline{\Gamma_S(\mathbf{T}(\mathbb{R})^0, e_f)} = \Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F).$$

This completes the proof of the lemma.

At the end of §5, we prove that if μ is the limiting measure of $(g_i)^*\mu_{\Gamma_S\backslash\Gamma_S H}$, then μ is a translate of the M-invariant probability measure $\mu_{\Gamma_S\backslash\Gamma_S M}$ for some subgroup M in $\mathbf{G}(\mathbb{Q}_S)$. Moreover, there is a \mathbb{Q} -subgroup \mathbf{M} of \mathbf{G} containing \mathbf{T} which does not have nontrivial \mathbb{Q} -characters such that $M \subset \mathbf{M}(\mathbb{Q}_S)$ is of finite index. Due to the condition that \mathbf{T} is \mathbb{R} -split, by Lemmas 5.4 and 6.3, and reviewing the induction argument in §5, we have $\mathbf{M}(\mathbb{R})^0 \subset M$ and

$$\Gamma_S \backslash \Gamma_S M = \overline{\Gamma_S \backslash \Gamma_S \mathbf{M}(\mathbb{R})^0}.$$

Furthermore, suppose that for any non-central element $x \in \mathbf{T}(\mathbb{Q})$ the sequence $g_i x g_i^{-1}$ diverges in $\mathbf{G}(\mathbb{Q}_S)$. We prove that $\mathbf{M} = \mathbf{G}$. By the discussion in §5, it is enough to prove that \mathbf{G} is the smallest group containing the sequence $\gamma_i \mathbf{T} \gamma_i^{-1}$ $(i \in \mathbb{N})$ where γ_i is defined at the beginning of §5. Note that since for any non-central element $x \in \mathbf{T}(\mathbb{Q})$ the sequence $g_i x g_i^{-1}$ diverges in $\mathbf{G}(\mathbb{Q}_S)$, the same is true for $\gamma_i x \gamma_i^{-1}$.

Suppose on the contrary that there is a subgroup **H** in **G** such that **H** contains infinitely many $\gamma_i \mathbf{T} \gamma_i^{-1}$. By [EMS97, Lemma 5.1] and Lemma 6.1, **H** is reductive and has no non-trivial \mathbb{Q} -characters. Since

$$\mathbf{T}(\mathbb{R})^0 \subset \mathbf{H}(\mathbb{R}) \text{ and } \gamma_i \mathbf{T}(\mathbb{R})^0 \gamma_i^{-1} \subset \mathbf{H}(\mathbb{R}),$$

by Lemma 6.2, there exists a finite subset $W \subset \mathbf{G}(\mathbb{R})$ such that for every γ_i , one can find $w_i \in W$ with $\gamma_i w_i \in \mathbf{H}(\mathbb{R})$. By passing to a subsequence, we may assume that $w_i = w$ for all $i \in \mathbb{N}$. Now let $\delta_i = \gamma_i \gamma_1^{-1}$ and $\mathbf{T}' = \gamma_1 \mathbf{T} \gamma_1^{-1}$. Then we have

$$\delta_i = \gamma_i w(\gamma_1 w)^{-1} \in \mathbf{H}(\mathbb{R}) \cap \mathbf{G}(\mathcal{O}_S) = \mathbf{H}(\mathcal{O}_S)$$

and $\{\delta_i \mathbf{T}' \delta_i^{-1}\} = \{\gamma_i \mathbf{T} \gamma_i^{-1}\}$ is contained in the \mathbb{Q} -reductive group \mathbf{H} . Since \mathbf{G} is \mathbb{R} -split and \mathbf{T}' is a maximal \mathbb{R} -split torus, there is a non-central element x in \mathbf{T}' such that x is in the center of \mathbf{H} . This implies that $\delta_i x \delta_i^{-1} = x$ for any i, which contradicts the assumption that $\gamma_i x \gamma_i^{-1}$ diverges in $\mathbf{G}(\mathbb{Q}_S)$.

Hence we proved the following

Theorem 6.4. Let $\{g_i\}$ be a sequence in $\mathbf{G}(\mathbb{Q}_S)$. Let $\mu_{\Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F)}$ be the periodic probability measure supported on $\Gamma_S \backslash \Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F)$. Then the sequence $(g_i)^* \mu_{\Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F)}$ has a subsequence converging to a periodic probability measure μ on $\Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$.

Moreover, assume that the measure μ is a translate of the M-invariant probability measure $\mu_{\Gamma_S\backslash\Gamma_S M}$ for some subgroup M in $\mathbf{G}(\mathbb{Q}_S)$. Then there is a \mathbb{Q} -subgroup M of G containing T which does not have nontrivial \mathbb{Q} -characters such that $\mathbf{M}(\mathbb{R})^0 \subset M$, $M \subset \mathbf{M}(\mathbb{Q}_S)$ is of finite index and $\Gamma_S\backslash\Gamma_S M = \overline{\Gamma_S\backslash\Gamma_S \mathbf{M}(\mathbb{R})^0}$. If for any non-central element $x \in \mathbf{T}(\mathbb{Q})$ the sequence $g_ixg_i^{-1}$ diverges in $\mathbf{G}(\mathbb{Q}_S)$, then $\mathbf{M} = \mathbf{G}$.

6.2. **Proof of Theorem 1.8.** Let S be a finite set of places of \mathbb{Q} containing the Archimedean one. We denote by

$$K^S := \prod_{p \notin S} \mathbf{G}(\mathbb{Z}_p), \ K_S := \prod_{p \in S_f} \mathbf{G}(\mathbb{Z}_p) \text{ and } K_f := \prod_{p \neq \infty} \mathbf{G}(\mathbb{Z}_p).$$

We write e_p , e_S and e_f for the identity in $\mathbf{G}(\mathbb{Q}_p)$, $\prod_{p\in S}\mathbf{G}(\mathbb{Q}_p)$ and $\prod_{p\in S_f}\mathbf{G}(\mathbb{Q}_p)$ respectively. If necessary, for an S-arithmetic element $g\in \prod_{p\in S}\mathbf{G}(\mathbb{Q}_p)$, we will index g to indicate its coordinates in the set of places S, that is, $g=(g_p)_{p\in S}$ where g_p is the coordinate of g at the place p.

It is well-known that $\Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$ is a finite disjoint union of closed $\mathbf{G}(\mathbb{R})K_S$ -orbits

$$\Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S) = \bigcup_i \Gamma_S \backslash \Gamma_S x_i \mathbf{G}(\mathbb{R}) K_S$$

for some $x_1 = e_S, x_2, \dots, x_n \in \mathbf{G}(\mathbb{Q}_S)$. One can define a natural projection map from the closed orbit $\Gamma_S \backslash \mathbf{G}(\mathbb{R}) K_S$ to $X_\infty = \Gamma_\infty \backslash G_\infty$ by sending

$$\mathbf{G}(\mathcal{O}_S)g = \mathbf{G}(\mathcal{O}_S)(g_{\infty}, g_f) \in \mathbf{G}(\mathcal{O}_S) \backslash \mathbf{G}(\mathcal{O}_S)\mathbf{G}(\mathbb{R})K_S$$

where $g_{\infty} \in \mathbf{G}(\mathbb{R})$ and $g_f \in K_S$, to the point $\Gamma_{\infty}g_{\infty}$ in $\Gamma_{\infty}\backslash\mathbf{G}(\mathbb{R})$. We denote it by

$$\varpi_{\infty}(\Gamma_S g) = \Gamma_{\infty} g_{\infty}.$$

Note that the map ϖ_{∞} is well-defined on $\Gamma_S \backslash \mathbf{G}(\mathbb{R}) K_S$, and intertwines with the actions of $\mathbf{G}(\mathbb{R})$ on the spaces $\Gamma_S \backslash \mathbf{G}(\mathbb{R}) K_S$ and $\Gamma_{\infty} \backslash \mathbf{G}(\mathbb{R})$. For a probability measure μ defined on $\Gamma_S \backslash \mathbf{G}(\mathbb{R}) K_S$, we denote by $\varpi_{\infty}^*(\mu)$ the pushforward of μ on $\Gamma_{\infty} \backslash G_{\infty}$.

Since $\mathbf{T}(\mathbb{Z})$ is a subgroup of the compact group $\prod_{p \in S_f} \mathbf{G}(\mathbb{Z}_p)$, one can choose a subgroup Λ of finite index in $\mathbf{T}(\mathbb{Z}) \cap \mathbf{T}(\mathbb{R})^0$ such that Λ is contained in a small neighborhood of identity in $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$ where the exponential map of p-adic Lie groups are defined for $p \in S_f$. Let F be the closure of Λ in $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$.

Lemma 6.5. For any $h \in \Gamma_S$ and $g \in \mathbf{G}(\mathbb{R})$,

$$\varpi_{\infty}(\Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F)(g^{-1}, h^{-1})) = \Gamma_{\infty}h\mathbf{T}(\mathbb{R})^0g^{-1}.$$

Proof. By Lemma 6.3, we know that

$$\overline{\Gamma_S \mathbf{T}(\mathbb{R})^0} = \Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F).$$

Now for any $g \in \mathbf{G}(\mathbb{R})$ and $h \in \Gamma_S$

$$\Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F)(g^{-1}, h^{-1})) = \overline{\Gamma_S(\mathbf{T}(\mathbb{R})^0 g^{-1} \times h^{-1})} = \overline{\Gamma_S(h\mathbf{T}(\mathbb{R})^0 g^{-1} \times e_f)}.$$

By continuity of ϖ_{∞} , we have

$$\varpi_{\infty}(\Gamma_{S}(\mathbf{T}(\mathbb{R})^{0} \times F)(g^{-1}, h^{-1})) = \varpi_{\infty}(\overline{\Gamma_{S}(h\mathbf{T}(\mathbb{R})^{0}g^{-1} \times \{e_{f}\})} = \Gamma_{\infty}h\mathbf{T}(\mathbb{R})^{0}g^{-1}.$$

This completes the proof of the lemma.

Lemma 6.6. Let μ be a probability measure on $\Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$ supported on $\Gamma_S \backslash \mathbf{G}(\mathbb{R})K_S$. Suppose that for some subgroup H_0 in $\mathbf{G}(\mathbb{R})$, μ is H_0 -invariant and $\varpi_{\infty}^*(\mu)$ is supported on an H_0 -orbit xH_0 in $\Gamma_{\infty} \backslash \mathbf{G}(\mathbb{R})$. Then $\varpi_{\infty}^*(\mu) = \mu_{xH_0}$. In particular, for any $g \in \mathbf{G}(\mathbb{R})$ and $h \in \Gamma_S$

$$\varpi_{\infty}^*((g,h)^*\mu_{\Gamma_S(\mathbf{T}(\mathbb{R})^0\times F)})=\mu_{\Gamma_{\infty}h\mathbf{T}(\mathbb{R})^0g^{-1}}.$$

Proof. Since μ is H_0 -invariant, the probability measure $\varpi_{\infty}^*(\mu)$ is also H_0 -invariant and hence equal to the Haar measure μ_{xH_0} on xH_0 . The second claim follows from Lemma 6.5.

The next lemma is straightforward.

Lemma 6.7. Suppose that a sequence of measures μ_i on $\Gamma_S \backslash \mathbf{G}(\mathbb{Q}_S)$ converges weakly to ν . Then $\varpi_{\infty}^*(\mu_i)$ converges weakly to $\varpi_{\infty}^*(\nu)$.

Deduction of Theorem 1.8 from Theorems 1.7 and 6.4. By Theorems 1.7 and 6.4, we know that after passing to a subsequence, $(g_i, h_i)^* \mu_{\Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F)}$ converges to some algebraic measure μ . Here the measure μ is supported on $\Gamma_S \backslash \Gamma_S Mg$ for some $g \in \mathbf{G}(\mathbb{Q}_S)$ and M a subgroup of $\mathbf{G}(\mathbb{Q}_S)$. Moreover, there is a \mathbb{Q} -subgroup \mathbf{M} of \mathbf{G} which does not have nontrivial \mathbb{Q} -characters such that $\mathbf{M}(\mathbb{R})^0 \subset M$ and $\Gamma_S \backslash \Gamma_S M = \overline{\Gamma_S \backslash \Gamma_S \mathbf{M}(\mathbb{R})^0}$. Since

$$\Gamma_S \backslash \Gamma_S(\mathbf{T}(\mathbb{R})^0 \times F)(g_i^{-1}, h_i^{-1}) = \overline{\Gamma_S \backslash \Gamma_S(h_i \mathbf{T}(\mathbb{R})^0 g_i^{-1} \times \{e_f\})} \subset \Gamma_S \backslash \Gamma_S \mathbf{G}(\mathbb{R}) K_S,$$

the support $\Gamma_S \backslash \Gamma_S Mg$ of the limiting measure μ is also contained in $\Gamma_S \backslash \Gamma_S \mathbf{G}(\mathbb{R}) K_S$. Hence there exist $\gamma \in \Gamma_S$ and $x = (x_\infty, x_f) \in \mathbf{G}(\mathbb{R}) \times \prod_{p \in S_f} \mathbf{G}(\mathbb{Z}_p)$ such that $g = \gamma x$. We have

$$\Gamma_S \backslash \Gamma_S Mg = \overline{\Gamma_S(\mathbf{M}(\mathbb{R})^0 \times \{e_f\})(\gamma x)} = \overline{\Gamma_S(\gamma^{-1}\mathbf{M}(\mathbb{R})^0 \gamma \times \{e_f\})x}$$

and

$$\varpi_{\infty}(\Gamma_S \backslash \Gamma_S Mg) = \Gamma_{\infty}(\gamma^{-1} \mathbf{M}(\mathbb{R})^0 \gamma) x_{\infty}.$$

Since $\Gamma_{\infty}(\gamma^{-1}\mathbf{M}(\mathbb{R})^{0}\gamma)x_{\infty}$ is closed and of finite volume, by Lemmas 6.6 and 6.7, we conclude that $\mu_{\Gamma_{\infty}h_{i}\mathbf{T}(\mathbb{R})^{0}g_{i}^{-1}}$ converges to the $x_{\infty}^{-1}(\gamma^{-1}\mathbf{M}(\mathbb{R})^{0}\gamma)x_{\infty}$ -invariant probability measure $\mu_{\Gamma_{\infty}(\gamma^{-1}\mathbf{M}(\mathbb{R})^{0}\gamma)x_{\infty}}$. This finishes the proof of Theorem 1.8.

7. Proof of Theorems 1.2,

We end the paper by showing how Theorem 1.2 follows as an application of Theorem 1.8 and then proving Theorem 1.5.

Proof of Theorem 1.2. Let m, $\vec{\alpha}$, \mathbf{i}_n , $\vec{\alpha}_n$, $g_{\vec{\alpha}_n}$ and s_n be as in the statement. We use the notation of §1.4. Let S_f be the set of primes dividing m and $S = S_f \cup \{\infty\}$. Let $\mathbf{G} = \mathbf{PGL}_d$, $\Gamma_{\infty} = \mathbf{G}(\mathbb{Z})$, and $X = \Gamma_{\infty} \setminus \mathbf{G}(\mathbb{R})$. Let A be the connected component of the identity (in the hausdorf topology) of the group of diagonal matrices in $\mathbf{G}(\mathbb{R})$. We need to prove that the sequence of periodic orbits

$$Y_n \stackrel{\text{def}}{=} \Gamma_{\infty} g_{\vec{\alpha}_n} A s_n^{-1}$$

equidistribute in X. The proof is an application of Theorem 1.8. In order to apply Theorem 1.8 we need to find a \mathbb{Q} -anisotropic maximal torus defined over \mathbb{Q} , $\mathbf{T} < \mathbf{G}$ and elements $h_n \in \Gamma_S$, $g_n \in \mathbf{G}(\mathbb{R})$ such that $Y_n = \Gamma_\infty h_n \mathbf{T}(\mathbb{R})^0 g_n$ and such that for any non-central $x \in \mathbf{T}(\mathbb{Q})$ the sequence $(g_n x g_n^{-1}, h_n x h_n^{-1}) \in \mathbf{G}(\mathbb{Q}_S)$ diverges.

Since $\Gamma_{\infty}g_{\vec{\alpha}}A$ is periodic, the intersection $\Gamma_{\infty} \cap g_{\vec{\alpha}}Ag_{\vec{\alpha}}^{-1}$ is a co-compact lattice in $g_{\vec{\alpha}}Ag_{\vec{\alpha}}^{-1}$. We denote its Zariski closure by $\mathbf{T} < \mathbf{G}$. Thus, \mathbf{T} is a \mathbb{Q} -anisotropic \mathbb{Q} -torus and $A = g_{\vec{\alpha}}^{-1}\mathbf{T}(\mathbb{R})^0g_{\vec{\alpha}}$ and

$$Y_n = \Gamma_{\infty} g_{\vec{\alpha}_n} A s_n^{-1} = \Gamma_{\infty} \operatorname{diag} \left(m^{i_1^{(n)}}, \dots, m^{i_d^{(n)}} \right) g_{\vec{\alpha}} A s_n^{-1} = \Gamma_{\infty} \operatorname{diag} \left(m^{i_1^{(n)}}, \dots, m^{i_d^{(n)}} \right) \mathbf{T}(\mathbb{R})^0 g_{\vec{\alpha}} s_n^{-1}.$$

So if we denote $h_n = \operatorname{diag}\left(m^{i_1^{(n)}}, \ldots, m^{i_d^{(n)}}\right) \in \Gamma_S$ and $g_n = g_{\vec{\alpha}} s_n^{-1} \in \mathbf{G}(\mathbb{R})$, then the equidistribution sought follows from Theorem 1.8 once we show that for any non-central $x \in \mathbf{T}(\mathbb{Q})$ we have that $h_n x h_n^{-1} \to \infty$ in $\prod_{p \in S_f} \mathbf{G}(\mathbb{Q}_p)$. Let $x \in \mathbf{T}(\mathbb{Q})$ be a non-central element and by abuse of notation, identify x with a rational matrix in $\mathbf{GL}_d(\mathbb{Q})$. We may assume by passing to a subsequence if necessary that the order between the coordinates of the vector \mathbf{i} is fixed. For simplicity we assume that for all n, $i_1^{(n)} < i_2^{(n)} < \cdots < i_d^{(n)}$. In this case for $1 \le j < r \le d$ such that $x_{jr} \ne 0$, the (j,r)'th coordinate of the conjugation $h_n x h_n^{-1}$ equals $m^{(i_j^{(n)} - i_r^{(n)})} x_{jr}$ which goes to ∞ in \mathbb{Q}_p for any $p \in S_f$ by our assumption that $|i_j^{(n)} - i_r^{(n)}| \to \infty$ and our assumption on the order. We conclude that the proof will be concluded once we show that x cannot be lower triangular. If x is lower triangular then since $g_{\vec{\alpha}}^{-1}\mathbf{T}(\mathbb{R})g_{\vec{\alpha}}$ is the group of diagonal matrices in $\mathbf{G}(\mathbb{R})$ we get a (rational) diagonal matrix a such that

$$xg_{\vec{\alpha}} = g_{\vec{\alpha}}a.$$

Recalling the definition of $g_{\vec{\alpha}}$ in (1) and comparing the first rows of the matrices in both sides of this equation we see (since x is assumed to be lower triangular) that

$$(x_{11}\sigma_1(\alpha_1), x_{11}\sigma_2(\alpha_1), \dots, x_{11}\sigma_d(\alpha_1)) = (a_{11}\sigma_1(\alpha_1), a_{22}\sigma_2(\alpha_1), \dots, a_{dd}\sigma_d(\alpha_1)).$$

As $\sigma_j(\alpha_1) \neq 0$ we get that $a_{jj} = x_{11}$ which means that a is a scalar matrix. This in turn implies that x itself (which is conjugate to a) is a scalar matrix – a contradiction to our assumption that x is non-central.

Proof of Theorem 1.5. Let

$$\vec{\alpha}_n \stackrel{\text{def}}{=} (m^{i_{n,1}}\alpha_1, m^{i_{n,2}}\alpha_2, \dots, m^{i_{n,d-1}}\alpha_{d-1}, 1)$$

be as in the statement. Let $\mathbb{K} \subset \mathbb{R}$ be the totally real number field generated by the coordinates of the vectors $\vec{\alpha}_n$. Let $\varphi = (\sigma_1, \dots, \sigma_{d-1}, \mathrm{id}_{\mathbb{K}})$ be an ordering of the different embeddings of \mathbb{K} into \mathbb{R} so that the last embedding is chosen to be the identity. Using similar notation as in (1), consider the homothety class $x_{\vec{\alpha}_n}$ of the lattice in \mathbb{R}^d spanned by the rows of the matrix

$$g_{\vec{\alpha}_n} = \begin{pmatrix} - & m^{i_{n,1}} \varphi(\alpha_1) & - \\ & \vdots & \\ - & m^{i_{n,d-1}} \varphi(\alpha_{d-1}) & - \\ - & \varphi(1) & - \end{pmatrix}.$$

Theorem 1.2 tells us that for any choice of real matrices s_n , the periodic orbits $s_n A x_{\vec{\alpha}_n}$ equidistribute in the space of homothety classes of lattices in \mathbb{R}^d with respect to the uniform measure on that space. We can now apply [SW, Theorem 14.2], which says that in this case,

$$\nu_{\mathbf{best}}^{\vec{\alpha}_n} \longrightarrow \nu_{\mathbf{best}}.$$
(2)

In fact, [SW, Theorem 14.2] only needs as an input the fact that the orbits $s_n A x_{\vec{\alpha}_n}$ equidistribute for a specific sequence of invertible real matrices s_n but since we know this for any such choice, there is no need to dwell on the particular formula of s_n that appears there.

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