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Sourav GHOSH

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AVATARS OF MARGULIS INVARIANTS AND PROPER ACTIONS

by Sourav GHOSH (*)

ABSTRACT. — In this article, we provide a necessary and sufficient criterion for proper actions on $\mathbb{H}^{n,n-1}$ in terms of certain special Anosov representations in $\mathrm{SO}(n, n)$. Moreover, we show that affine Anosov representations of any word hyperbolic group in $\mathrm{SO}(n, n-1) \ltimes \mathbb{R}^{2n-1}$ are infinitesimal versions of such special Anosov representations. Finally, using the above two results we interpret Margulis spacetimes as infinitesimal versions of quotient manifolds of $\mathbb{H}^{n,n-1}$.

In the appendix, we give a description of the appropriate cross-ratios in our setting and their infinitesimal versions.

RÉSUMÉ. — Dans cet article, nous fournissons un critère nécessaire et suffisant pour des actions propres sur l'espace $\mathbb{H}^{n,n-1}$ en termes de certaines représentations Anosov spéciales dans le groupe $\mathrm{SO}(n, n)$. De plus, nous montrons que les représentations Anosov affines de tout groupe hyperbolique dans le groupe $\mathrm{SO}(n, n-1) \ltimes \mathbb{R}^{2n-1}$ sont des versions infinitésimales de ces représentations Anosov spéciales. Enfin, en utilisant les deux résultats ci-dessus, nous interprétons les espaces-temps de Margulis comme des versions infinitésimales de variétés quotients de l'espace $\mathbb{H}^{n,n-1}$.

En annexe, nous donnons une description des birapports appropriés dans notre contexte et de leurs versions infinitésimales.

1. Introduction

The study of tilings gives rise to the study of proper actions. In a celebrated result, Bieberbach classified the symmetries of “crystals”. Any subgroup of $\mathrm{O}(n, \mathbb{R}) \ltimes \mathbb{R}^n$ whose action on \mathbb{R}^n is proper and cocompact is called a *crystallographic* group. Bieberbach [7, 8] (see also [13]) showed that for each fixed n there are only finitely many isomorphism classes of n -dimensional crystallographic groups and they contain a normal subgroup of

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finite index isomorphic to \mathbb{Z}^n . Later, Auslander and Markus [5] constructed examples of subgroups of $\text{Aff}(n, \mathbb{R}) := \text{GL}(n, \mathbb{R}) \ltimes \mathbb{R}^n$ whose action on \mathbb{R}^n are proper and cocompact and which do not have a normal subgroup of finite index isomorphic to \mathbb{Z}^n . The examples they constructed were normal subgroups of finite index isomorphic to $\mathbb{Z}^{\ltimes n}$. Auslander [3, 4] attempted to show that any subgroup of $\text{O}(n, \mathbb{R}) \ltimes \mathbb{R}^n$ whose action on \mathbb{R}^n is proper and cocompact is virtually polycyclic and failed. Later, the statement was rechristened as *Auslander conjecture* by Fried–Goldman [26]. This conjecture is still unsolved for the general case but has been shown to hold true in dimension less than 7 by the work of Fried–Goldman [26] and Abels–Margulis–Soifer [2]. Recently, a generalization of the Auslander conjecture for homogeneous spaces appeared in [70] with partial resolution. In fact, the study of proper action of discrete groups on homogeneous spaces can be traced back to Borel [10]. He studied compact quotients of symmetric spaces. Later, the study of pseudo-Riemannian homogeneous spaces was pioneered by Kulkarni [53], Kobayashi [50], Benoist [6] and Okuda [63]. They found that the signature of a homogeneous space acts as an obstruction for the existence of proper actions of a discrete group on it (please check [51] for a recent survey). Let $g \in \text{GL}(n, \mathbb{R})$, $v \in \mathbb{R}^n$ and $(g, v) \in \text{Aff}(n, \mathbb{R})$. Then we call $\text{GL}(n, \mathbb{R})$ (resp. g) the *linear* part of $\text{Aff}(n, \mathbb{R})$ (resp. (g, v)) and \mathbb{R}^n (resp. v) the *translational* part of $\text{Aff}(n, \mathbb{R})$ (resp. (g, v)). We observe that the map L sending any element to its linear part is a homomorphism and for any subgroup H of $\text{Aff}(n, \mathbb{R})$ we call $L(H)$ the linear part of H .

Meanwhile, Margulis [58, 59] answered a question asked by Milnor [61] related to the Auslander conjecture in the negative. He showed that the conjecture would fail if one drops the cocompactness assumption. He constructed actions of non-abelian free subgroups of $\text{GL}(3, \mathbb{R}) \ltimes \mathbb{R}^3$ which act properly on \mathbb{R}^3 . It was known due to prior works of Kostant–Sullivan [52] and Fried–Goldman [26] that the linear part of the Zariski closure of such a group has to be $\text{SO}(2, 1)$. The quotient manifolds obtained from \mathbb{R}^3 under the proper action of a non-abelian free subgroup of $\text{SO}(2, 1) \ltimes \mathbb{R}^3$ are called *Margulis spacetimes*. Furthermore, Drumm [24] classified the linear holonomy of Margulis spacetimes and constructed [23] nice fundamental domains for a large class of them. It was conjectured that any Margulis space-time admits a fundamental domain of the type constructed by Drumm. This conjecture is known as the *Crooked plane conjecture*. Recently, the Crooked plane conjecture was resolved by Danciger–Guéritaud–Kassel [20]. While constructing Margulis spacetimes, Margulis introduced certain invariants,

which later came to be known as the Margulis invariants, to give a necessary criterion for proper actions. Later, Labourie [54] introduced continuous versions of the Margulis invariants. We call these invariants as Labourie–Margulis invariants. Subsequently, Goldman–Labourie–Margulis [35] used the Labourie–Margulis invariants to give a necessary and sufficient condition for a representation to be a Margulis spacetime. Classification results of similar nature were also independently obtained by Danciger–Guéritaud–Kassel [19] using different techniques. In fact, using those techniques they went on to prove a twenty year old conjecture of Drumm–Goldman. They demonstrated that Margulis spacetimes are tame, i. e. Margulis spacetimes are homeomorphic to the interior of a compact manifold with boundary. Tameness of Margulis spacetimes was also independently proved by Choi–Goldman [17]. Recently, similar results for Margulis spacetimes with parabolics was obtained by Choi–Drumm–Goldman [16]. The construction of Margulis spacetimes and Margulis invariants were generalized by Abels–Margulis–Soifer [1]. For n even they constructed more examples of non-abelian free subgroups of $\mathrm{GL}(2n-1, \mathbb{R}) \ltimes \mathbb{R}^{2n-1}$ whose Zariski closure is $\mathrm{SO}(n, n-1) \ltimes \mathbb{R}^{2n-1}$ and which act properly on \mathbb{R}^{2n-1} . They also proved the nonexistence of any such subgroups for odd n . Suppose Γ is a finitely generated hyperbolic group. We abuse notation and call an injective homomorphism $(\rho, u) : \Gamma \rightarrow \mathrm{SO}(n, n-1) \ltimes \mathbb{R}^{2n-1}$, a Margulis spacetime if $(\rho, u)(\Gamma) \backslash \mathbb{R}^{2n-1}$ is a manifold. The constructions of [1] were further generalized by Smilga in [67, 68]. Generalizing previous works of Kim [49] and Ghosh [30], recently it was proved in [32] that only finitely many well chosen Margulis invariants are enough to determine a conjugacy class of Margulis spacetimes. It might appear from the above discussion that free groups are the only groups which admit proper affine actions but that is not the case. First examples of proper affine actions of right-angled Coxeter groups was obtained by Danciger–Guéritaud–Kassel [21]. It remains to be seen if there are other classes of groups which admit proper actions or not.

The classification results of Goldman–Labourie–Margulis [35] (see also Goldman–Labourie [34]) point towards a description of Margulis spacetimes in terms of Anosov representations. This project is still ongoing with affirmative results by Ghosh [29, 31] and Ghosh–Treib [33]. Anosov representations generalize the notion of convex cocompactness for higher rank Lie groups [38]. They were introduced by Labourie [55] to provide a geometric characterization of certain special representations of surface

groups called *Hitchin representations*. Hitchin representations can be defined as those representations which lie in the components of the representation variety containing Fuchsian representations of a surface group. One interesting thing about Hitchin representations is that their moduli space is topologically trivial and it is the same as the moduli of solutions of the Yang–Mills equations under certain symmetry conditions [42]. Interestingly, Hitchin representations also make their appearance in the study of Margulis spacetimes through the works of Danciger–Zhang [22] and Labourie [56]. Generalizing previous results of Goldman–Margulis [36], Mess [60] and Labourie [54], they showed that Hitchin representations in $\mathrm{SO}(n, n-1)$ do not admit affine deformations which act properly on \mathbb{R}^{2n-1} . The definition of an Anosov representation given by Labourie is dynamical in nature and its dynamics resembles the notion of an Axiom A flow appearing in the dynamical systems literature [11, 65]. Later on, Guichard–Wienhard [40] extended the notion of an Anosov representation to include representations of any finitely generated hyperbolic group into semisimple Lie groups. They also established relations between Anosov representations and proper actions on homogeneous spaces. Subsequently, a more algebraic description of Anosov representations, in terms of uniform gaps in singular value or eigenvalue spectra, and relation between Anosov representations and proper actions on homogeneous spaces appeared in the works of Kapovich–Leeb–Porti [44, 45, 46], Guéritaud–Guichard–Kassel–Wienhard [38], Bochi–Potrie–Sambarino [9] and Kassel–Potrie [48]. In this article, we stick to the dynamical description of an Anosov representation. Further dynamical properties of these representations were proved by Bridgeman–Canary–Labourie–Sambarino in [12]. The appropriate version of Anosov representations of a hyperbolic group into affine Lie groups of the kind $\mathrm{SO}(n, n-1) \ltimes \mathbb{R}^{2n-1}$, was introduced in the works of Ghosh [28, 29] and Ghosh–Treib [33]. They showed that Margulis spacetimes can be characterized by affine Anosov representations. We prove a similar characterization relating proper actions of hyperbolic groups on the homogeneous space $\mathbb{H}^{n, n-1}$ and certain special Anosov representations (this generalization provide an alternate proof of a statement from [22] communicated to us by Danciger and Zhang). Let Γ be a finitely generated hyperbolic group. We denote the inner product on \mathbb{R}^{2n} whose symmetry group is $\mathrm{SO}(n, n)$ by $\langle \cdot | \cdot \rangle$ and the connected component containing identity of $\mathrm{SO}(n, n)$ by G . Let $\mathbb{R}^{n, n} = \mathbb{R}^{n, n-1} \oplus \mathbb{R}e_{2n}$ and G_0 be the group of those elements in G which fix the vector e_{2n} . Let $\mathbb{H}^{n, n-1} := Ge_{2n} \subset \mathbb{R}^{2n}$ and $G_a := G_0 \ltimes \mathbb{R}^{2n-1}$. We observe that $\mathbb{H}^{n, n-1} \cong G/G_0$. We call a subspace of \mathbb{R}^{2n-1} a *null* subspace

if its orthogonal is a maximal isotropic subspace of $\mathbb{R}^{n,n-1}$. We note that maximal isotropic subspaces of $\mathbb{R}^{n,n-1}$ are $(n-1)$ -dimensional and null subspaces of $\mathbb{R}^{n,n-1}$ are n -dimensional. Let P_0^\pm be the stabilizer inside G_0 of two transverse n -dimensional null subspaces of $\mathbb{R}^{n,n-1}$, P_a^\pm be the stabilizer inside G_a of two transverse n -dimensional null subspaces of $\mathbb{R}^{n,n-1}$ and P^\pm be the stabilizer inside G of two $(n-1)$ -dimensional isotropic subspaces of $\mathbb{R}^{n,n}$ which are oriented and whose orthogonal subspaces are transverse to each other (for more details see Sections 3.1 and 4.1). Moreover, we denote $\text{Hom}(\Gamma, G_0, P_0^\pm)$ (resp. $\text{Hom}(\Gamma, G, P^\pm)$) to be the space of all injective homomorphisms of Γ inside G_0 (resp. G) which are Anosov with respect to P_0^\pm (resp. P^\pm) and we denote $\text{Hom}(\Gamma, G_a, P_a^\pm)$ to be the space of all injective homomorphisms of Γ inside G_a which are affine Anosov with respect to P_a^\pm . We prove the following:

THEOREM 1.1 (see Theorem 3.32). — *We recall that G is the connected component of $\text{SO}(n, n)$ containing identity. Suppose ρ is a representation of Γ in G which is Anosov with respect to the stabilizer of an oriented $(n-1)$ -dimensional isotropic subspace of $\mathbb{R}^{n,n}$, that is, $\rho \in \text{Hom}(\Gamma, G, P^\pm)$. Then the action of $\rho(\Gamma)$ on $\mathbb{H}^{n,n-1}$ is proper if and only if ρ is Anosov in $\text{SL}(2n, \mathbb{R})$ with respect to the stabilizer of an n -dimensional subspace.*

Goldman–Margulis [36] kicked off the study of Margulis spacetimes via deformation of related objects. They showed that the marked Margulis invariant spectrum of an affine representation in dimension three can be interpreted as derivatives of the marked length spectrum of surfaces. Deformation techniques were also used by Danciger–Guéritaud–Kassel [19] to show that three dimensional Margulis spacetimes are the rescaled limits of collapsing AdS spacetimes. In fact, their proof of the Crooked plane conjecture [21] also used deformation techniques but in the context of arc complexes. A general framework to describe transitions between geometries which are sub-geometries of a larger ambient geometry was developed by Cooper–Danciger–Wienhard [18]. Furthermore, in a different context Kassel [47] and Guéritaud–Kassel [39] showed that small deformations of proper actions still give rise to proper actions. In this article, we generalize the derivative interpretation of the Margulis invariants and show that Margulis invariant spectra can be obtained as the derivative of the middle eigenvalue gap spectra of representations in G . Let $h \in G_0$ be *pseudo-hyperbolic* i. e. the unit eigenspace of h is one dimensional and h does not have -1 as an eigenvalue. Let W_\pm^h respectively be the subspaces of \mathbb{R}^{2n-1} on which the action of h is contracting or expanding. There is a consistent way of choosing a direction along the unique eigenspace of h with eigenvalue 1 (for

more details please see Section 4.1). Let v_0^h be the unique eigenvector of h with eigenvalue 1 such that $\langle v_0^h | v_0^h \rangle = 1$ and which is positively oriented with respect to the aforementioned choice of direction. In particular, when $n = 2$, the subspaces W_\pm^h are one dimensional and light-like. The space of light-like vectors is a double cone. The choice in this case is done as follows: we choose non zero vectors $v_\pm^h \in W_\pm^h$ which lie in the upper cone and choose v_0^h in such a way that (v_-^h, v_0^h, v_+^h) is positively oriented with respect to the standard orientation on \mathbb{R}^3 . Then the *Margulis invariant* of an element $g = (h, u) \in G_a$ is defined as

$$\alpha(h, u) := \langle u | v_0^h \rangle.$$

We observe that the dimensions of W_\pm^h are $(n - 1)$. Let $h_t \in G$ be an analytic one parameter family with $h = h_0$. Then, for t small enough we obtain a pair of attracting and repelling subspaces of h_t denoted by $W_\pm^{h_t}$ and whose dimensions are also $(n - 1)$. Moreover, there is a consistent way of choosing a maximal isotropic subspace $V_+^{h_t}$ inside $(W_+^{h_t})^\perp \subset \mathbb{R}^{2n}$. It is the maximal isotropic subspace of $(W_+^{h_t})^\perp$ which is the deformation of $V_+^h = W_+^h \oplus \mathbb{R}(v_0^h + e_{2n}) \subset (W_+^h)^\perp$ (for more details please see Sections 3.1 and 4.1). We observe that $(W_-^h)^\perp \cap V_+^h = \mathbb{R}(v_0^h + e_{2n})$ and that h_t preserves the line $(W_-^{h_t})^\perp \cap V_+^{h_t}$. Let $\lambda(h_t)$ be the eigenvalue of the action of h_t on this line. We prove the following:

LEMMA 1.2. — *Suppose $g \in G_0$ is a pseudo-hyperbolic element and $g_t \in G$ is an analytic one parameter family whose tangent direction at $g = g_0$ is G and $Ge_{2n} = v$. Then*

$$\alpha(g, v) = \left. \frac{d}{dt} \right|_{t=0} \lambda(g_t).$$

We give a similar interpretation of Labourie–Margulis invariant too and use it to show that the notion of an affine Anosov representation detect the infinitesimal versions of certain special Anosov representations in G :

THEOREM 1.3 (see Theorem 4.21). — *We recall that $\mathbb{R}^{n,n} = \mathbb{R}^{n,n-1} \oplus \mathbb{R}e_{2n}$, G is the connected component of $SO(n, n)$ containing identity, G_0 is the subgroup of G which fixes e_{2n} and $G_a = G_0 \ltimes \mathbb{R}^{n,n-1}$.*

Let $\{\rho_t\}_{t \in (-1,1)}$ be an analytic one parameter family of representations of Γ in G with $\rho_0(\Gamma) \subset G_0$. Let U be the tangent vector to $\{\rho_t\}_{t \in (-1,1)}$ at $\rho = \rho_0$ and $u = Ue_{2n}$. Suppose (ρ, u) is affine Anosov in G_a with respect to the stabilizer of an n -dimensional null subspace of $\mathbb{R}^{n,n-1}$, that is, $(\rho, u) \in \text{Hom}(\Gamma, G_a, P_a^\pm)$. Then there exists $\epsilon > 0$ such that for all t with $|t| \in (0, \epsilon)$, ρ_t is Anosov in $SL(2n, \mathbb{R})$ with respect to the stabilizer of an n -dimensional subspace.

Finally, combining these two results and using the characterization of Margulis spacetimes in terms of affine Anosov representations, we show that Margulis spacetimes of dimension $(2n - 1)$ are infinitesimal versions of manifolds obtained from quotients of the homogeneous space $\mathbb{H}^{n,n-1}$:

COROLLARY 1.4. — *Let $\{\rho_t\}_{t \in (-1,1)}$ be an analytic one parameter family of representations of Γ in \mathbf{G} with $\rho = \rho_0$ Anosov in \mathbf{G}_0 with respect to the stabilizer of an n -dimensional null subspace of $\mathbb{R}^{n,n-1}$. Let U be the tangent vector to $\{\rho_t\}_{t \in (-1,1)}$ at $\rho = \rho_0$ and $u = Ue_{2n}$. Suppose (ρ, u) is a Margulis spacetime. Then there exists $\epsilon > 0$ such that for all t with $|t| \in (0, \epsilon)$, $\rho_t(\Gamma)$ acts properly on $\mathbb{H}^{n,n-1}$.*

Moreover, using the work of Abels–Margulis–Soifer [1] and our main result we provide an alternate proof of the following fact which was first obtained by Benoist [6]:

COROLLARY 1.5. — *Suppose n is even. Then there exists a non-abelian free subgroup with finitely many generators inside \mathbf{G} which act properly on $\mathbb{H}^{n,n-1}$.*

In the appendix, we also generalize results by Charette–Drumm [15] and Ghosh [29]. We define affine cross ratios β for any four mutually transverse affine null subspaces in $\mathbb{R}^{n,n-1}$ and show the following:

PROPOSITION 1.6. — *Suppose n is even, $\rho \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$ and $(\rho, u) \in \text{Hom}(\Gamma, \mathbf{G}_a)$. Suppose $\gamma, \eta \in \Gamma$ are two infinite order elements such that the four points $\gamma^\pm, \eta^\pm \in \partial_\infty \Gamma$ are distinct and the sequence $\{\gamma^m \eta^m\}_{m \in \mathbb{N}} \subset \Gamma$ contains a subsequence $\{\gamma^{n_i} \eta^{n_i}\}_{i \in \mathbb{N}}$ consisting only of infinite order elements. Then the following identity holds:*

$$\lim_{i \rightarrow \infty} (\alpha(\gamma^{n_i} \eta^{n_i}) - \alpha(\gamma^{n_i}) - \alpha(\eta^{n_i})) = \beta(\eta^-, \gamma^-, \gamma^+, \eta^+).$$

Moreover, we also define the linear counterparts θ of these affine cross ratios defined for any four mutually transverse $(n - 1)$ -dimensional isotropic subspaces in $\mathbb{R}^{n,n}$ and show that

PROPOSITION 1.7. — *Suppose n is even, $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ and $\gamma, \eta \in \Gamma$ are two infinite order elements such that the four points $\gamma^\pm, \eta^\pm \in \partial_\infty \Gamma$ are distinct and the sequence $\{\gamma^m \eta^m\}_{m \in \mathbb{N}} \subset \Gamma$ contains a subsequence $\{\gamma^{n_i} \eta^{n_i}\}_{i \in \mathbb{N}}$ consisting only of infinite order elements. Then the following identity holds:*

$$\lim_{i \rightarrow \infty} \frac{\lambda(\gamma^{n_i} \eta^{n_i})^2}{\lambda(\gamma^{n_i})^2 \lambda(\eta^{n_i})^2} = \theta(\eta^-, \gamma^-, \gamma^+, \eta^+)^2.$$

Remark 1.8. — Lastly, we would like to mention that Danciger–Zhang has announced independent work in [22] which has overlap with some of our results. In particular, Theorem 1.1, Lemma 1.2 and Theorem 1.3 of this article when applied to fundamental groups of compact surfaces without boundary, are respectively similar to Lemma 8.2, Theorem 8.8 and Theorem 6.1 of [22]. On the other hand, the results about cross ratios contained in Sections A.1 and A.2 are not obtained by Danciger–Zhang. We would also like to note that, even though Corollary 1.4 has not been stated as a result in [22], for the case of fundamental groups of compact surfaces without boundary, it can also be obtained by jointly applying Theorems 8.8 and 6.1 of [22]. In [22], Danciger–Zhang use Lemma 8.2, Theorem 8.8, Theorem 6.1 and they also use certain properties special to Hitchin representations obtained from the works of Labourie [55] and Fock–Goncharov [25] to generalize Theorem 1.1 of [54] and conclude that representations in $\mathrm{PSL}(2n-1, \mathbb{R}) \ltimes \mathbb{R}^{2n-1}$ whose linear parts are Hitchin do not admit proper affine actions on \mathbb{R}^{2n-1} .

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2. Preliminaries

In this section we introduce certain preliminary notions and results needed to establish our results.

2.1. Anosov representations

In this subsection we define the notion of an Anosov representation and mention some important properties of Anosov representations which will

be used later on. Anosov representations into $\mathrm{SL}(n, \mathbb{R})$ were introduced by Labourie in [55] to show that Hitchin representations satisfy certain nice geometric properties. Later on, Guichard–Wienhard [40] extended the notion of an Anosov representation to representations of any hyperbolic group into a semisimple Lie group. Recently, Kapovich–Leeb–Porti gave a different algebraic characterization of Anosov representations in [44] and [46]. In this article, we use the dynamical definition of an Anosov representation from the work of Labourie [55] and Guichard–Wienhard [40].

We start by defining the Gromov flow space. It plays a very central role in the dynamical definition of an Anosov representation. Let Γ be a finitely generated word hyperbolic group, $\partial_\infty \Gamma$ be its boundary at infinity and let

$$\partial_\infty \Gamma^{(2)} := \{(p_+, p_-) \mid p_\pm \in \partial_\infty \Gamma, p_+ \neq p_-\}.$$

Gromov [37] (see also Champetier [14] and Mineyev [62]) constructed a cocompact, proper action of Γ on $\widetilde{\mathrm{U}\Gamma} := \partial_\infty \Gamma^{(2)} \times \mathbb{R}$, which commutes with the flow:

$$\begin{aligned} \phi_t : \widetilde{\mathrm{U}\Gamma} &\longrightarrow \widetilde{\mathrm{U}\Gamma} \\ p := (p_+, p_-, p_0) &\longmapsto (p_+, p_-, p_0 + t) \end{aligned}$$

and whose restriction on $\partial_\infty \Gamma^{(2)}$ is the diagonal action coming from the natural action of Γ on its boundary $\partial_\infty \Gamma$. Moreover, there exists a metric on $\widetilde{\mathrm{U}\Gamma}$ well defined up to Hölder equivalence such that the Γ action is isometric, the flow ϕ_t acts by Lipschitz homeomorphisms and every orbit of the flow $\{\phi_t\}_{t \in \mathbb{R}}$ gives a quasi-isometric embedding. The resulting quotient space denoted by $\mathrm{U}\Gamma$ is called the *Gromov flow space*. We note that the Gromov flow space is connected and it admits partitions of unity (for more details please see [33]).

The other important ingredients in the definition of Anosov representations are parabolic subgroups (for a detailed exposition on parabolic subgroups please see Section 3.2 of [40]). In this article, to prove certain results we also work with stabilizers of oriented subspaces. These groups strictly speaking are not necessarily parabolic subgroups but they are subgroups of finite index inside parabolic subgroups. We call such subgroups *virtually parabolic* subgroups. Hence, although the theory of Anosov representations due to Guichard–Wienhard [40] does not directly apply to these cases, the original theory due to Labourie [55] does. Let \mathbf{H} be a connected semisimple Lie group. Moreover, let \mathbf{P}_\pm be a pair of virtually parabolic subgroups of \mathbf{H} such that their respective overgroups which are parabolic are opposite and $\mathbf{P}_+ \cap \mathbf{P}_-$ is a finite index subgroup of the intersection of the two opposite parabolic overgroups. We call such a pair of virtually parabolic subgroups

as *opposite*. Let $\mathcal{X} \subset \mathbf{H}/\mathbf{P}_+ \times \mathbf{H}/\mathbf{P}_-$ be the space of all pairs $(h\mathbf{P}_+, h\mathbf{P}_-)$ for $h \in \mathbf{H}$. We consider the left action of \mathbf{H} on $\mathbf{H}/\mathbf{P}_+ \times \mathbf{H}/\mathbf{P}_-$ and observe that the action is transitive on \mathcal{X} and the stabilizer of the point $(\mathbf{P}_+, \mathbf{P}_-) \in \mathcal{X}$ is $\mathbf{P}_+ \cap \mathbf{P}_-$. Hence $\mathbf{H}/(\mathbf{P}_+ \cap \mathbf{P}_-) \cong \mathcal{X}$. Moreover, \mathcal{X} is open in $\mathbf{H}/\mathbf{P}_+ \times \mathbf{H}/\mathbf{P}_-$. Therefore,

$$\mathbf{T}_{(h\mathbf{P}_+, h\mathbf{P}_-)}\mathcal{X} = \mathbf{T}_{h\mathbf{P}_+}\mathbf{H}/\mathbf{P}_+ \oplus \mathbf{T}_{h\mathbf{P}_-}\mathbf{H}/\mathbf{P}_-.$$

DEFINITION 2.1. — Let Γ be a hyperbolic group and let \mathbf{H} be a semisimple Lie group with a pair of opposite virtually parabolic subgroups \mathbf{P}_\pm . Then any representation $\rho : \Gamma \rightarrow \mathbf{H}$ is called \mathbf{P}_\pm -Anosov if and only if

- (1) There exist continuous, injective, $\rho(\Gamma)$ -equivariant limit maps

$$\xi^\pm : \partial_\infty \Gamma \longrightarrow \mathbf{H}/\mathbf{P}_\pm$$

such that $\xi(p) := (\xi^+(p_+), \xi^-(p_-)) \in \mathcal{X}$ for any $p \in \widetilde{\mathbf{U}\Gamma}$.

- (2) There exist positive constants C, c and a continuous collection of $\rho(\Gamma)$ -equivariant Euclidean metrics $\|\cdot\|_p$ on $\mathbf{T}_{\xi(p)}\mathcal{X}$ for $p \in \widetilde{\mathbf{U}\Gamma}$ such that

$$\|v^\pm\|_{\phi_{\pm t}p} \leq Ce^{-ct}\|v^\pm\|_p$$

for all $v^\pm \in \mathbf{T}_{\xi^\pm(p_\pm)}\mathbf{H}/\mathbf{P}_\pm$ and for all $t \geq 0$.

Notation 2.2. — We denote the space of all representations $\rho : \Gamma \rightarrow \mathbf{G}$ by $\text{Hom}(\Gamma, \mathbf{G})$, the space of \mathbf{P}_\pm -Anosov representations in \mathbf{G} by $\text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}_\pm)$ and the space of all α -Hölder maps from $\partial_\infty \Gamma$ to $\mathbf{H}/\mathbf{P}_\pm$ by $\mathcal{C}^\alpha(\partial_\infty \Gamma, \mathbf{H}/\mathbf{P}_\pm)$.

Now we state a few theorems which will be important for us later on.

THEOREM 2.3 (Labourie [55], Guichard–Wienhard [40]). — Suppose $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}_\pm)$. Then there exists an open neighborhood $U \subset \text{Hom}(\Gamma, \mathbf{H})$ with $\rho \in U$ such that $U \subset \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}_\pm)$.

THEOREM 2.4 (Bridgeman–Canary–Labourie–Sambarino [12]). — Suppose $U \subset \text{Hom}(\Gamma, \mathbf{H}, \mathbf{P}_\pm)$ is an open ball and suppose ξ_ρ^\pm are the limit maps of $\rho \in U$. Then there exists $U_0 \subset U$ such that for all $\rho \in U_0$ the limit maps $\xi_\rho^\pm \in \mathcal{C}^\alpha(\partial_\infty \Gamma, \mathbf{H}/\mathbf{P}_\pm)$ for some $\alpha > 0$ and the following map is analytic:

$$\begin{aligned} \xi^\pm : U_0 &\longrightarrow \mathcal{C}^\alpha(\partial_\infty \Gamma, \mathbf{H}/\mathbf{P}_\pm) \\ \rho &\longmapsto \xi_\rho^\pm. \end{aligned}$$

THEOREM 2.5 (Guichard–Wienhard [40]). — Suppose \mathbf{P}_\pm and \mathbf{Q}_\pm respectively are two pairs of opposite virtually parabolic subgroups of \mathbf{H} such that \mathbf{P}_\pm respectively are subgroups of \mathbf{Q}_\pm . Then $\text{Hom}(\Gamma, \mathbf{H}, \mathbf{P}_\pm) \subset \text{Hom}(\Gamma, \mathbf{H}, \mathbf{Q}_\pm)$.

Remark 2.6. — Theorems 2.3 and 2.4 as stated above is not proved in [40] and [12] but these results also hold true if we replace parabolic subgroups with virtually parabolic subgroups. This is because the proofs given in [12] (Theorems 6.1, 6.5, 6.6 and Lemma 6.7) only depend on the fact that the space H/P is an analytic manifold and the limit map at the origin satisfies the contraction property (for more details please see Theorem 3.8 of [41] and Theorem 5.18 of [66]). Finally, Theorem 2.5 would also hold true in this setting as we would get the new limit map for free from the old limit map by composing it with the natural projection of H/P onto H/Q and the contraction property would still hold due to its independence from the particular collection of Euclidean norms chosen. Hence, in the remainder of this article we would use the notion of an Anosov representation to include groups of this general nature too.

Recently, Stecker–Treib [69] have given a more algebraic way of characterizing Anosov representations with respect to stabilizers of oriented flags. They use techniques from [44, 45, 46] and [38] to construct domains of discontinuity for oriented flag manifolds.

2.2. Pseudo-orthogonal groups

In this subsection we give a brief description of certain well known properties of pseudo-orthogonal groups which will be helpful later on.

Let \mathbb{R}^k be the space of all column vectors of size k (i. e. matrices of size $k \times 1$) and let I_k be the identity matrix of size $k \times k$. We endow \mathbb{R}^{p+q} with the following quadratic form of signature (p, q) :

$$I_{p,q} := \begin{bmatrix} I_p & 0 \\ 0 & -I_q \end{bmatrix}.$$

We denote \mathbb{R}^{p+q} endowed with the quadratic form $I_{p,q}$ by $\mathbb{R}^{p,q}$ and the group of invertible linear transformations preserving the quadratic form $I_{p,q}$ by $O(p, q)$. We note that any element of $O(p, q)$ either has determinant 1 or -1 and denote the subgroup of $O(p, q)$ whose elements have determinant 1 by $SO(p, q)$ and denote the connected component of $SO(p, q)$ which contains the identity transformation by $SO_0(p, q)$.

Remark 2.7. — Suppose $g \in O(p, q)$ is such that A_g is a $p \times p$ matrix and

$$g = \begin{bmatrix} A_g & B_g \\ C_g & D_g \end{bmatrix}.$$

Then both A_g and D_g are invertible. Let sgn denote the sign of a non-zero real number. We note that $\text{O}(p, q)$ has four connected components characterized as follows:

$$\text{O}_d^a(p, q) := \{g \mid \text{sgn det}(A_g) = a \text{ and } \text{sgn det}(D_g) = d\},$$

where $a, d \in \{\pm\}$. Moreover, using a quick computation we obtain that $\det(g) = \det(A_g) \det(D_g)^{-1}$. Hence, $\text{SO}_0(p, q) = \text{O}_+^+(p, q)$ and $\text{SO}(p, q) = \text{O}_+^+(p, q) \cup \text{O}_-^-(p, q)$. (For more details please see Chapter 9 of [64] and Proposition 7.3 of [27].)

LEMMA 2.8. — *Any element of the group $\text{O}(k, k)$ which fixes $[I_k, I_k]^t$ lies in $\text{SO}_0(k, k)$.*

Proof. — We start by observing that any matrix X of dimension $2k \times 2k$ which satisfy the equation $X[I_k, I_k]^t = [I_k, I_k]^t$ is of the following form for some matrices A and C of dimension $k \times k$,

$$X = \begin{bmatrix} A & (I - A) \\ (I - C) & C \end{bmatrix}$$

Now we characterize such elements of $\text{O}(k, k)$. We observe that $X^t I_{k,k} X = I_{k,k}$ impose the following set of constraints on A and C :

$$\begin{aligned} A^t A - (I - C^t)(I - C) &= I, \\ (I - A)^t(I - A) - C^t C &= -I, \\ A^t(I - A) - (I - C^t)C &= 0. \end{aligned}$$

The first and third equations give us $(A + C) = 2I$ and replacing $(I - C)$ with $(A - I)$ in the first equation we deduce that $A^t + A = 2I$.

Let us denote $(A^t - A)$ by $2M$. We note that M is skew-symmetric. Hence, we obtain that $A = (I - M)$ and $C = (I + M)$. As M is skew-symmetric, using the spectral theorem for skew-symmetric matrices we deduce that $\det(I - M) > 0$ and $\det(I + M) > 0$. Finally, using Remark 2.7 we conclude our result. \square

LEMMA 2.9. — *Suppose $k = \min\{p, q\}$ and $p \neq q$. Then any element of $\text{SO}(p, q)$ which fixes $[I_k, 0, I_k]^t$ lies inside $\text{SO}_0(p, q)$.*

Proof. — Let us denote $(p - q)$ by m . We start by observing that any matrix X of dimension $(p + q) \times (p + q)$ which satisfy the equation $X[I_k, 0, I_k]^t = [I_k, 0, I_k]^t$ is of the following form for some matrices A, C of dimension $k \times k$,

B of dimension $|m| \times |m|$ and column vectors $u, v, w \in \mathbb{R}^k$,

$$X = \begin{bmatrix} A & v & (I - A) \\ u^t & B & -u^t \\ (I - C) & w & C \end{bmatrix}$$

Now we characterize all such elements of $O(p, q)$. We observe that $X^t I_{p,q} X = I_{p,q}$ impose the following set of constraints on A, B, C, u, v and w :

$$\begin{aligned} A^t A + \operatorname{sgn}(m) u u^t - (I - C^t)(I - C) &= I, \\ (I - A)^t (I - A) + \operatorname{sgn}(m) u u^t - C^t C &= -I, \\ A^t (I - A) - \operatorname{sgn}(m) u u^t - (I - C^t) C &= 0, \\ v^t v + \operatorname{sgn}(m) B^t B - w^t w &= \operatorname{sgn}(m) I, \\ A^t v + \operatorname{sgn}(m) u B - (I - C^t) w &= 0, \\ (I - A^t) v - \operatorname{sgn}(m) u B - C^t w &= 0. \end{aligned}$$

The last two equations give us $v = w$ and plugging it in the fourth equation we obtain that $B \in O(|m|)$. Furthermore, from the first and third equations we obtain that $(A + C) = 2I$ and replacing $(I - C)$ with $(A - I)$ in the first equation we deduce that $\operatorname{sgn}(m) u u^t = 2I - A^t - A$. As $v = w$, replacing $(I - C)$ with $(A - I)$ in the fifth equation gives us $\operatorname{sgn}(m) u B + w = 0$. It follows that, any $X \in O(p, q)$ which fix $[I_k, 0, I_k]^t$ are of the following form

$$X = \begin{bmatrix} A & -\operatorname{sgn}(m) u & (I - A) \\ u^t & I & -u^t \\ (A - I) & -\operatorname{sgn}(m) u & (2I - A) \end{bmatrix} \begin{bmatrix} I & 0 & 0 \\ 0 & B & 0 \\ 0 & 0 & I \end{bmatrix}.$$

We observe that

$$\begin{bmatrix} I & 0 & -I \\ 0 & I & 0 \\ 0 & 0 & I \end{bmatrix} \begin{bmatrix} A & -\operatorname{sgn}(m) u & (I - A) \\ u^t & I & -u^t \\ (A - I) & -\operatorname{sgn}(m) u & (2I - A) \end{bmatrix} \begin{bmatrix} I & 0 & I \\ 0 & I & 0 \\ 0 & 0 & I \end{bmatrix} = \begin{bmatrix} I & 0 & 0 \\ * & I & 0 \\ * & * & I \end{bmatrix}.$$

Hence, for $X \in SO(p, q)$ we have $\det(B) = 1$.

Let us denote $(A^t - A)$ by M . We note that M is skew-symmetric. Hence, for $p > q$ we obtain that $\det(2I - A) = \det(I + (M + u u^t)/2)$ and for $p < q$ we obtain that $\det(A) = \det(I - (M - u u^t)/2)$. As M is skew-symmetric we have $\det(I + (M + u u^t)/2) = \det(I - (M - u u^t)/2)$. Moreover, as $I + u u^t/2$ is a symmetric positive definite matrix, using Cholesky factorization (see Corollary 7.2.9 of [43]) we obtain the existence of a real lower triangular matrix T with positive entries on the diagonal such that $I + u u^t/2 = T T^t$. Hence,

$$\det(I + (u u^t + M)/2) = \det(T) \det(I + T^{-1}(M/2)(T^t)^{-1}) \det(T^t).$$

As $T^{-1}M(T^t)^{-1}$ is skew-symmetric, using the spectral theorem for skew-symmetric matrices we deduce that $\det(I + T^{-1}(M/2)(T^t)^{-1}) > 0$. Finally, using Remark 2.7 we conclude our result. \square

2.3. Maximal isotropic spaces

Let V be a subspace of $\mathbb{R}^{p,q}$. Then V is called *isotropic* if and only if for all $v \in V$ we have $v^t I_{p,q} v = 0$. In this subsection we list and demonstrate a few properties of maximal isotropic subspaces. These properties are used many times in the later part of this text. Although we expect these results to be well known, we were unable to find appropriate sources in the literature for them. Keeping in mind their centrality we have decided to include them.

Remark 2.10. — Let V be a maximal isotropic space of $\mathbb{R}^{p,q}$. Then there exists a maximal isotropic space W which is transverse to V i. e. $V^\perp \cap W = \{0\}$. Note that V, W are spaces of dimension $k = \min\{p, q\}$ and the existence of W is guaranteed by the fact that $I_{p,q}V$ is a maximal isotropic space which is transverse to V .

Notation 2.11. — We denote the space spanned by the column vectors of a matrix X by $\text{cspan}(X)$.

LEMMA 2.12. — Suppose (V, W) is a pair of transverse maximal isotropic subspaces of $\mathbb{R}^{p,q}$, suppose $k = \min\{p, q\}$ and $J = I_{2k-1,1}$. Then

- (1) for $p = q = k$, there exists $g \in \text{SO}_0(p, q)$ such that either of the following holds:

$$\begin{aligned} (V, W) &= (g\text{cspan}([I_k, I_k]^t), g\text{cspan}([I_k, -I_k]^t)), \\ (V, W) &= (gJ\text{cspan}([I_k, I_k]^t), gJ\text{cspan}([I_k, -I_k]^t)), \end{aligned}$$

- (2) and for $p \neq q$, there exists $g \in \text{SO}_0(p, q)$ such that

$$(V, W) = (g\text{cspan}([I_k, 0, I_k]^t), g\text{cspan}([I_k, 0, -I_k]^t)).$$

Proof. — Let $B = \{v_1, \dots, v_k\}$ be a basis of V . Let $B_j = B \setminus \{v_j\}$. Then $\dim(B_j^\perp \cap W) \geq 1$. Choose $w_j \in (B_j^\perp \cap W)$ such that $v_j^t I_{p,q} w_j = 1/2$. We note that the existence of such a w_j is guaranteed by the maximality of V .

Case $p = q = k$. — We observe that $V^\perp \cap W^\perp = \{0\}$ and deduce that

$$h = [v_1 + w_1, \dots, v_k + w_k, v_1 - w_1, \dots, v_k - w_k] \in \text{O}(k, k).$$

Hence, $h[I_k, I_k]^t = 2[v_1, \dots, v_k]$ and $h[I_k, -I_k]^t = 2[w_1, \dots, w_k]$. Now we use Remark 2.7 and observe that upto a sign change in v_1 (consequently in

w_1) we can make sure that $h \in \mathcal{O}_+^+(k, k) \cup \mathcal{O}_-^+(k, k)$. Now, if $h \in \mathcal{SO}_0(k, k)$, then we choose $g = h$. Otherwise, we choose $g = hJ$. We use Remark 2.7 to observe that $g \in \mathcal{SO}_0(k, k)$, $gJ[I_k, I_k]^t = 2[v_1, \dots, v_k]$ and $gJ[I_k, -I_k]^t = 2[w_1, \dots, w_k]$.

Case $p \neq q$. — As both V and W are maximal, any non-zero $u \in V^\perp \cap W^\perp$ satisfy $u^t I_{p,q} u \neq 0$. As $u^t I_{p,q} u = (-u)^t I_{p,q} (-u)$, using continuity we obtain that either $I_{p,q}$ restricts to a positive definite form on $V^\perp \cap W^\perp$ or a negative definite one. In either case, using a Gram-Schmidt process we obtain an orthonormal basis $u_1, \dots, u_{|p-q|}$. Now we use Remark 2.7 and observe that upto a sign change in v_1 (consequently in w_1 too) and in u_1 , $g = [v_1 + w_1, \dots, v_k + w_k, u_1, \dots, u_{|p-q|}, v_1 - w_1, \dots, v_k - w_k] \in \mathcal{SO}_0(p, q)$. Hence, V is spanned by the column vectors of $g[I_k, 0, I_k]$ and W is spanned by the column vectors of $g[I_k, 0, -I_k]$. \square

Notation 2.13. — Suppose $K \in \mathcal{GL}(k, \mathbb{R})$. We consider the following functions:

$$f_{k,0}(K) := \begin{bmatrix} \frac{K^{-1}+K^t}{2} & \frac{K^{-1}-K^t}{2} \\ \frac{K^{-1}-K^t}{2} & \frac{K^{-1}+K^t}{2} \end{bmatrix} \text{ and } f_{k,m}(K) := \begin{bmatrix} \frac{K^{-1}+K^t}{2} & 0 & \frac{K^{-1}-K^t}{2} \\ 0 & I & 0 \\ \frac{K^{-1}-K^t}{2} & 0 & \frac{K^{-1}+K^t}{2} \end{bmatrix}$$

and observe that $f_{k,m}(K) \in \mathcal{SO}(k+m, k) \cap \mathcal{SO}(k, m+k)$ for all $m \geq 0$.

LEMMA 2.14. — Suppose $V, V' \subset \mathbb{R}^{k,k}$ are such that V is the span of the column vectors of $[I_k, I_k]^t$ and V' is the span of the column vectors of $[I_k, I_{k-m,m}]^t$ with odd m . Then there does not exist any $g \in \mathcal{SO}(k, k)$ such that $gV = V'$.

Proof. — We observe that $J = I_{2k-m,m} \in \mathcal{O}(k, k)$. We will prove this result by contradiction. If possible let us assume that $gV = V'$ for some $g \in \mathcal{SO}(k, k)$ then $JgV = V$. Hence, $Jg[I_k, I_k]^t = [I_k, I_k]^t K$ for some invertible $k \times k$ matrix K . We observe that $f_{k,0}(K)Jg[I_k, I_k]^t = [I_k, I_k]^t$. Hence, using Lemma 2.8 we have $f_{k,0}(K)Jg \in \mathcal{SO}_0(k, k)$. As both $f_{k,0}(K)$ and g have determinant 1, we deduce that $\det(J) = 1$, a contradiction. \square

LEMMA 2.15. — Any element of $\mathcal{SO}_0(p, q)$ which preserves a maximal isotropic space also preserves the orientation on it.

Proof. — Let V be a maximal isotropic subspace of $\mathbb{R}^{p,q}$ and let $g \in \mathcal{SO}_0(p, q)$ be such that $gV = V$.

Case $p = q = k$. — We use Lemma 2.12 to obtain that there exists $h \in \mathcal{O}(k, k)$ such that $h^{-1}gh[I_k, I_k]^t = [I_k, I_k]^t K$ for some $k \times k$ matrix K . We observe that $h^{-1}gh \in \mathcal{SO}_0(k, k)$. Hence, $f_{k,0}(K)h^{-1}gh \in \mathcal{SO}(k, k)$.

Also, $f_{k,0}(K)h^{-1}gh[I_k, I_k]^t = [I_k, I_k]^t$. Now we use Lemma 2.8 and obtain that $f_{k,0}(K)h^{-1}gh \in \mathrm{SO}_0(k, k)$. Therefore, $f_{k,0}(K) \in \mathrm{SO}_0(k, k)$ and by Remark 2.7 it follows that $\det(K) > 0$.

Case $p \neq q$. — We denote $|p - q|$ by m and $\min\{p, q\}$ by k . We use Lemma 2.12 to obtain that there exists $h \in \mathrm{O}(p, q)$ such that $h^{-1}gh[I_k, 0, I_k]^t = [I_k, 0, I_k]^t K$ for some $k \times k$ matrix K . We observe that $h^{-1}gh \in \mathrm{SO}_0(p, q)$. Hence, $f_{k,m}(K)h^{-1}gh \in \mathrm{SO}(p, q)$. Also,

$$f_{k,m}(K)h^{-1}gh[I_k, 0, I_k]^t = [I_k, 0, I_k]^t.$$

Now we use Lemma 2.9 and obtain that $f_{k,m}(K)h^{-1}gh \in \mathrm{SO}_0(p, q)$. Therefore, $f_{k,m}(K) \in \mathrm{SO}_0(p, q)$ and by Remark 2.7 it follows that $\det(K) > 0$. \square

LEMMA 2.16. — *Let W be a $(k - 1)$ -dimensional isotropic subspace of $\mathbb{R}^{k,k}$. Then W^\perp contain exactly two different maximal isotropic subspaces. One of them lies in the orbit of $\mathrm{cspan}([I_k, I_k]^t)$ under the action of $\mathrm{SO}_0(n, n)$ and the other one lies in the orbit of $\mathrm{cspan}(J[I_k, I_k]^t)$ under the action of $\mathrm{SO}_0(n, n)$.*

Proof. — As the dimension of W is $(k - 1)$, we can extend W to a maximal isotropic subspace V of $\mathbb{R}^{k,k}$. Clearly, $V = V^\perp \subset W^\perp$. Moreover, let U be a maximal isotropic subspace of $\mathbb{R}^{k,k}$ transverse to V . We choose a basis $B = \{v_1, \dots, v_k\}$ of V such that $W = \mathrm{span}(B \setminus \{v_k\})$. Moreover, suppose $\{u_1, \dots, u_k\}$ is a basis of U such that $u_j \in (B \setminus \{v_j\})^\perp \cap U$ and $u_j^t I_{k,k} v_j = 1$. Note that by construction $u_k \in W^\perp$ and $u_k^t I_{k,k} v_k = 1$. Hence, $W \oplus \mathbb{R}u_k$ is a maximal isotropic subspace of W^\perp different from V .

Also, $W^\perp = W \oplus \mathbb{R}v_k \oplus \mathbb{R}u_k$. Hence, for any $w \in W^\perp$, there exist $a, b \in \mathbb{R}$ such that $w - av_k - bu_k \in W$. Moreover, if $\langle w | w \rangle = 0$, then $ab = 0$ and it follows that either $w \in V$ or $w \in W \oplus \mathbb{R}u_k$.

Finally, for $J = I_{2k-1,1}$ we observe that

$$h = [v_1 + u_1, \dots, v_k + u_k, v_1 - u_1, \dots, v_k - u_k] \in \mathrm{O}(k, k),$$

$h[I_k, I_k]^t = 2[v_1, \dots, v_{k-1}, v_k]$ and $hJ[I_k, I_k]^t = 2[v_1, \dots, v_{k-1}, u_k]$. Moreover, up to a sign change we can choose v_1 (consequently u_1), such that $h \in \mathrm{O}_+^+(k, k) \cup \mathrm{O}_-^+(k, k)$. Finally, we conclude by noting that $h \in \mathrm{O}_\pm^+(k, k)$ if and only if $hJ \in \mathrm{O}_\mp^+(k, k)$. \square

3. Proper actions on $\mathbb{H}^{n,n-1}$

In this section we demonstrate a necessary and sufficient criterion for proper actions of a hyperbolic group Γ on $\mathbb{H}^{n,n-1}$ in terms of certain

special Anosov representations. This statement is inspired from Danciger–Zhang [22] and we provide an alternate proof. The constructions, statements and proofs in this section parallels similar constructions, statements and proofs given for affine spaces in [33].

3.1. Proto-neutral sections

In this subsection we introduce a few notions in order to state the necessary and sufficient criterion of Theorem 1.1.

We consider \mathbb{R}^{2n} and henceforth for all vectors v, w we denote $v^t I_{n,n} w$ by $\langle v \mid w \rangle$. Let $e_j \in \mathbb{R}^{2n}$ be the column vector whose only non-vanishing entry is the j -th entry and its j -th entry is 1. We consider the embedding of \mathbb{R}^{2n-1} inside $\mathbb{R}^{n,n}$ which is spanned by the vectors e_1, \dots, e_{2n-1} . We note that the quadratic form $I_{n,n}$ induces a form of signature $(n, n-1)$ on \mathbb{R}^{2n-1} . Henceforth, we denote this embedding along with the induced form by $\mathbb{R}^{n,n-1}$, the group $\mathrm{SO}_0(n, n)$ by G and the subgroup of G which fixes e_{2n} by G_0 i. e.

$$G_0 := \{g \in G \mid ge_{2n} = e_{2n}\}.$$

We observe that $ge_{2n} = e_{2n}$ imply $g^t e_{2n} = e_{2n}$. Hence, G_0 is a subgroup of $O(n, n-1)$. Moreover, we use Remark 2.7 to obtain that $G_0 \cong \mathrm{SO}_0(n, n-1)$. We fix the following four subspaces:

$$V_{\pm} := \mathrm{span}\{e_j \pm e_{n+j} \mid j = 1, \dots, n\},$$

$$W_{\pm} := \mathrm{span}\{e_j \pm e_{n+j} \mid j = 1, \dots, n-1\}.$$

We observe that V_{\pm} are a transverse pair of maximal isotropic subspaces of $\mathbb{R}^{n,n}$ and W_{\pm} are a transverse pair of maximal isotropic subspaces of $\mathbb{R}^{n,n-1}$. We define \vec{V}_{\pm} to be the space V_{\pm} along with the orientation coming from the ordered basis $(e_1 \pm e_{n+1}, \dots, e_n \pm e_{2n})$ and \vec{W}_{\pm} to be the space W_{\pm} along with the orientation coming from $(e_1 \pm e_{n+1}, \dots, e_{n-1} \pm e_{2n-1})$. Furthermore, we define

$$P^{\pm} := \mathrm{Stab}_G(\vec{W}_{\pm}) \text{ and } P_0^{\pm} := \mathrm{Stab}_{G_0}(\vec{W}_{\pm}).$$

Remark 3.1. — We use Lemma 2.15 to deduce that $P_0^{\pm} = \mathrm{Stab}_{G_0}(W_{\pm})$.

We observe that $P_0^{\pm} = P^{\pm} \cap G_0$ and $G_0/P_0^{\pm} \subset G/P^{\pm}$. Let $L_0 := P_0^+ \cap P_0^-$ and $L := P^+ \cap P^-$.

Remark 3.2. — We observe that $W_+ \oplus W_- \oplus \mathbb{R}e_n = \mathbb{R}^{2n-1}$. As elements of L_0 preserve orientations on W_{\pm} , we obtain that they preserve orientations on $\mathbb{R}e_n$ too. Also, for any $g \in G_0$ we have $\langle ge_n \mid ge_n \rangle = 1$. It follows that the elements of L_0 preserve the vector e_n .

We use the above Lemma and obtain a well defined map:

$$\begin{aligned}\nu : \mathbf{G}_0/\mathbf{L}_0 &\longrightarrow \mathbb{R}^{2n+1} \\ [g] &\longmapsto ge_n.\end{aligned}$$

The map ν is called the *neutral section*.

Remark 3.3. — Any element in $\mathbf{G}_0/\mathbf{L}_0$ can be interpreted as a tuple of transverse null subspaces. The neutral section keeps track of the directional intersection of a pair of transverse null subspaces (for more details about its usage please see Remark 4.1). In particular, for $n = 2$, let (V, W) be a pair of transverse null subspaces of \mathbb{R}^3 . Then V, W are tangent planes to the light-cone, they touch the light-cone at two distinct lines. In fact, $V \cap W$ is also a line. Let $v \in V$, $w \in W$ be such that v, w lie in the upper light-cone and let $u \in V \cap W$ be such that $[v, u, w]$ give the standard orientation on \mathbb{R}^3 . Then we have

$$\langle u \mid u \rangle > 0 \text{ and } \nu(V, W) = \frac{u}{\sqrt{\langle u \mid u \rangle}}.$$

Henceforth, we fix an Euclidean norm $\|\cdot\|$ on \mathbb{R}^{2n} and choose the following two vectors

$$v_+ := \frac{e_n + e_{2n}}{2} \text{ and } v_- := \frac{e_n - e_{2n}}{2}.$$

LEMMA 3.4. — Suppose $\mathbf{L} = \mathbf{P}^+ \cap \mathbf{P}^-$. Then the elements of \mathbf{L} preserve the orientations on the isotropic lines $\mathbb{R}v_\pm$.

Proof. — We observe that $V_\pm = W_\pm \oplus \mathbb{R}v_\pm$. As elements of \mathbf{L} preserve the oriented space \vec{W}_\pm , we deduce that they preserve orientations on $\mathbb{R}v_+ \oplus \mathbb{R}v_-$. Hence, the action of \mathbf{L} on $\mathbb{R}v_+ \oplus \mathbb{R}v_-$ is isomorphic to the action of $\mathbf{SO}(1, 1)$ on $\mathbb{R}^{1,1}$. As the action of $\mathbf{SO}(1, 1)$ on $\mathbb{R}^{1,1}$ preserve the individual isotropic lines, we deduce that the elements of \mathbf{L} preserve the lines $\mathbb{R}v_\pm$. Hence, the action of \mathbf{L} preserves V_\pm . Moreover, V_\pm are maximal isotropic subspaces of $\mathbb{R}^{n,n}$. Hence, by Lemma 2.15 we obtain that the action of \mathbf{L} preserves \vec{V}_\pm . In fact, the action of \mathbf{L} also preserves \vec{W}_\pm . It follows that the elements of \mathbf{L} preserve the orientations on the isotropic lines $\mathbb{R}v_\pm$. \square

We use the above Lemma and obtain two well defined maps:

$$\begin{aligned}\nu_\pm : \mathbf{G}/\mathbf{L} &\longrightarrow \mathbb{R}^{2n} \\ [g] &\longmapsto gv_\pm / \|gv_\pm\|.\end{aligned}$$

We also observe that for all $g, h \in \mathbf{G}$, $\|\nu_\pm([g])\| = 1$ and

$$h\nu_\pm([g]) = \frac{\|hgv_\pm\|}{\|gv_\pm\|} \nu_\pm([hg]) = \|h\nu_\pm([g])\| \nu_\pm(h[g]).$$

We call the maps ν^\pm as *proto-neutral sections* as they play a role similar to the one played by the neutral section ν in [33].

Remark 3.5. — Any element in \mathbf{G}/\mathbf{L} can be interpreted as a tuple of oriented $(n-1)$ -dimensional isotropic subspaces of $\mathbb{R}^{n,n}$ whose orthogonal spaces are transverse to each other. The intersection of these two orthogonal spaces is a space of dimension two which contains exactly two isotropic lines. The proto-neutral sections provide a consistent way of distinguishing one isotropic line from the other of the two isotropic lines (for more details about its usage please see Remark 4.10).

LEMMA 3.6. — Suppose $x, y \in \{\pm\}$ and $g \in \mathbf{G}$ is such that $g \in \mathbf{P}^x$ then

$$W_x + \mathbb{R}\nu_y([g]) = W_x + \mathbb{R}v_y.$$

Proof. — Suppose $g \in \mathbf{P}^x$. Hence, $gW_x = W_x$ and $(W_x + \mathbb{R}\nu_y([g]))$ is in the \mathbf{G} -orbit of $(W_x + \mathbb{R}v_y)$. As $W_x \subset (W_x + \mathbb{R}v_y)$ and $(W_x + \mathbb{R}v_y)$ is maximally isotropic, we have $(W_x + \mathbb{R}v_y) \subset W_x^\perp$. It follows that

$$g(W_x + \mathbb{R}v_y) \subset gW_x^\perp = (gW_x)^\perp = W_x^\perp.$$

As $g(W_x + \mathbb{R}v_y)$ is in the \mathbf{G} -orbit of $(W_x + \mathbb{R}v_y)$ and both of them are maximal isotropic subspaces of W_x^\perp , we use Lemma 2.16 and conclude that $g(W_x + \mathbb{R}v_y) = (W_x + \mathbb{R}v_y)$ and our result follows. \square

LEMMA 3.7. — Suppose $x, y \in \{\pm\}$ and $g, h \in \mathbf{G}$ satisfy $g^{-1}h \in \mathbf{P}^x$. Then

$$\begin{aligned} \langle \nu_y([g]) \mid \nu_y([h]) \rangle &= 0, \\ \langle \nu_y([g]) \mid \nu_{-y}([h]) \rangle &\neq 0. \end{aligned}$$

Proof. — As $g^{-1}h \in \mathbf{P}^x$, using Lemma 3.6 we deduce that both $g^{-1}hv_y$ and v_y lie in the same maximal isotropic subspace $W_x \oplus \mathbb{R}v_y$. Hence, $\langle v_y \mid g^{-1}hv_y \rangle = 0$. Moreover, $g^{-1}hv_y \in W_x \oplus \mathbb{R}v_y$ and $g^{-1}hv_y \notin W_x$. Hence, $g^{-1}hv_y + av_y \in W_x$ for some $a \neq 0$. It follows that $\langle v_{-y} \mid g^{-1}hv_y \rangle = a\langle v_{-y} \mid v_y \rangle \neq 0$. Therefore,

$$\begin{aligned} \langle \nu_y([g]) \mid \nu_y([h]) \rangle &= \frac{\langle gv_y \mid hv_y \rangle}{\|gv_y\| \|hv_y\|} = \frac{\langle v_y \mid g^{-1}hv_y \rangle}{\|gv_y\| \|hv_y\|} = 0, \\ \langle \nu_y([g]) \mid \nu_{-y}([h]) \rangle &= \frac{\langle gv_y \mid hv_{-y} \rangle}{\|gv_y\| \|hv_{-y}\|} = \frac{\langle v_y \mid g^{-1}hv_{-y} \rangle}{\|gv_y\| \|hv_{-y}\|} \neq 0, \end{aligned}$$

and our result follows. \square

Remark 3.8. — We note that $I_{n,n} \in \mathbf{O}_+^+(n, n)$ for even n and $I_{n,n} \in \mathbf{O}_-^+(n, n)$ for odd n . Hence, for even n we deduce that $\nu_\pm([I_{n,n}]) = \nu_\mp([I])$,

$$\nu_\pm([gI_{n,n}]) = \frac{g\nu_\pm([I_{n,n}])}{\|g\nu_\pm([I_{n,n}])\|} = \frac{g\nu_\mp([I])}{\|g\nu_\mp([I])\|} = \nu_\mp([g]).$$

Moreover, $-I \in \mathbf{O}_+^+(n, n)$ for even n and we have $\nu_\pm([-g]) = -\nu_\pm([g])$.

3.2. Proto-Labourie–Margulis invariants

In this subsection we define the notion of diffused eigenvalues and call them proto-Labourie–Margulis invariants due to their relation with diffused Margulis invariants introduced by Labourie.

Suppose Γ is a word hyperbolic group and $\mathbf{U}\Gamma$ is its Gromov flow space. We denote the period of the periodic orbit in $\mathbf{U}\Gamma$ corresponding to an infinite order element $\gamma \in \Gamma$ by $l(\gamma)$ and the flow invariant probability measure on this periodic orbit by μ_γ . Now for any $\rho \in \mathbf{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$, any infinite order element $\gamma \in \Gamma$ and $t \in \mathbb{R}$ we have,

$$\nu_\pm(\xi_\rho(\gamma^+, \gamma^-, t)) = \nu_\pm(\xi_\rho(\gamma^+, \gamma^-, 0)).$$

Moreover, we denote $(\gamma^+, \gamma^-, 0)$ by p_γ and denote $\nu_\pm \circ \xi_\rho$ by ν_ρ^\pm . We observe that $\nu_\rho^\pm(p_\gamma)$ is an eigenvector of $\rho(\gamma)$. Indeed, we have

$$\nu_\rho^\pm(p_\gamma) = \nu_\pm(\rho(\gamma)\xi_\rho(p_\gamma)) = \frac{\rho(\gamma)\nu_\rho^\pm(p_\gamma)}{\|\rho(\gamma)\nu_\rho^\pm(p_\gamma)\|}.$$

As $\langle \rho(\gamma)\nu_\rho^+(p_\gamma) \mid \rho(\gamma)\nu_\rho^-(p_\gamma) \rangle = \langle \nu_\rho^+(p_\gamma) \mid \nu_\rho^-(p_\gamma) \rangle$ it follows that

$$\|\rho(\gamma)\nu_\rho^+(p_\gamma)\| \|\rho(\gamma)\nu_\rho^-(p_\gamma)\| = 1.$$

Notation 3.9. — We denote the eigenvalues of $\rho(\gamma)$ corresponding to the eigenvectors $\nu_\rho^\pm(p_\gamma)$ by $\lambda_\rho^\pm(\gamma)$. We observe that

$$\lambda_\rho^\pm(\gamma) = \|\rho(\gamma)\nu_\rho^\pm(p_\gamma)\| \text{ and } \lambda_\rho^+(\gamma)\lambda_\rho^-(\gamma) = 1.$$

Suppose $\{\rho_t\}_{t \in (-1, 1)} \subset \mathbf{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ is an analytic one parameter family. We now present a partition of unity type argument to construct a family of Euclidean norms indexed by points in the Gromov flow space $\mathbf{U}\Gamma$ and \mathbf{P}^\pm -Anosov representations in \mathbf{G} . Recall that $\pi : \widetilde{\mathbf{U}\Gamma} \rightarrow \mathbf{U}\Gamma$ is the standard projection map.

Remark 3.10. — As $\mathbf{U}\Gamma$ is compact, there exist $\{V_i\}_{i=1}^k$ such that $V_i \subset \widetilde{\mathbf{U}\Gamma}$ are small open balls and $\bigcup_{i=1}^k \pi(V_i) = \mathbf{U}\Gamma$. Hence $\bigcup_{\gamma \in \Gamma} \bigcup_{i=1}^k \gamma V_i = \widetilde{\mathbf{U}\Gamma}$. We know from Section 8.2 of [33] that there exist maps

$$\{f_i : \mathbf{U}\Gamma \rightarrow \mathbb{R}^+\}_{i=1}^k$$

with $\text{Supp}(f_i) \subset \pi(V_i)$ such that the functions f_i are Hölder continuous and differentiable along flow lines with $\sum_{i=1}^k f_i = 1$. We use this to construct a collection

$$\left\{ \|\cdot\|_p^t \mid (\rho_t, p) \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm) \times \widetilde{\mathbf{U}\Gamma} \right\}$$

of Euclidean norms on \mathbb{R}^{2n} indexed by $\text{Hom}(\Gamma, \mathbf{G}) \times \widetilde{\mathbf{U}\Gamma}$ such that:

- (1) it is Hölder continuous in the variable $p \in \widetilde{\mathbf{U}\Gamma}$,
- (2) it is smooth along the flow lines of $\{\phi_s\}_{s \in \mathbb{R}}$,
- (3) it is analytic along the variable ρ_t ,
- (4) it is equivariant i. e. $\|\rho_t(\gamma)v\|_{\gamma p}^t = \|v\|_p^t$ for all $v \in \mathbb{R}^{2n}$ and $\gamma \in \Gamma$.

We start by considering the fixed Euclidean norm $\|\cdot\|$ on \mathbb{R}^{2n} mentioned in Section 3.1. We observe that for any $p \in \Gamma V_i$ there exists a unique $\gamma_{p,i}$ such that $\gamma_{p,i}p \in V_i$. Note that in such a situation $\gamma_{\eta p,i}\eta = \gamma_{p,i}$. We define

$$\|v\|_{p,i}^t := \|\rho_t(\gamma_{p,i})v\|$$

for all $v \in \mathbb{R}^{2n}$ and for any $p \in \widetilde{\mathbf{U}\Gamma}$ we define:

$$\|v\|_p^t := \sum_{i=1}^k f_i(\pi(p)) \|v\|_{p,i}^t.$$

We check that this collection of norms are equivariant. Indeed, as

$$\begin{aligned} \|\rho_t(\gamma)v\|_{\gamma p}^t &= \sum_{i=1}^k f_i(\pi(\gamma p)) \|\rho_t(\gamma)v\|_{\gamma p,i}^t = \sum_{i=1}^k f_i(\pi(p)) \|\rho_t(\gamma_{\gamma p,i})\rho_t(\gamma)v\| \\ &= \sum_{i=1}^k f_i(\pi(p)) \|\rho_t(\gamma_{\gamma p,i}\gamma)v\| = \sum_{i=1}^k f_i(\pi(p)) \|\rho_t(\gamma_{p,i})v\| \\ &= \sum_{i=1}^k f_i(\pi(p)) \|v\|_{p,i}^t = \|v\|_p^t. \end{aligned}$$

Moreover, it follows from our construction that this collection of norms satisfies all the first three conditions listed above. More details about properties of these kinds of constructions can be found in [41] and [66] (see also [28, 29] and [12]).

Notation 3.11. — Henceforth, to simplify our notations we denote $\nu_\pm \circ \xi$ (resp. $\nu_\pm \circ \xi_t$) by ν^\pm (resp. ν_t^\pm) and observe that by definition it is invariant under the flow ϕ .

Suppose $\xi : \widetilde{\mathbf{U}\Gamma} \rightarrow (\mathbf{G}/\mathbf{L})$ are the limit maps corresponding to the representations $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$. We recall the proto-neutral sections from

Section 3.1 and using the above collection of norms we define:

$$\begin{aligned}\sigma^\pm : \widetilde{\mathrm{U}\Gamma} &\longrightarrow \mathbb{R}^{2n} \\ p &\longmapsto \nu^\pm(p) / \|\nu^\pm(p)\|_p.\end{aligned}$$

We call these maps *proto-neutralised sections*. We observe that for all $p \in \widetilde{\mathrm{U}\Gamma}$, $\langle \sigma^+(p) \mid \sigma^-(p) \rangle > 0$, and for $\gamma \in \Gamma$ we obtain that,

$$\sigma^\pm(\gamma p) = \frac{\nu^\pm(\gamma p)}{\|\nu^\pm(\gamma p)\|_{\gamma p}} = \frac{\rho(\gamma)\nu^\pm(p)}{\|\rho(\gamma)\nu^\pm(p)\|_{\gamma p}} = \frac{\rho(\gamma)\nu^\pm(p)}{\|\nu^\pm(p)\|_p} = \rho(\gamma)\sigma^\pm(p).$$

Also, we deduce that

$$\sigma^\pm(\phi_s p) = \frac{\nu^\pm(\phi_s p)}{\|\nu^\pm(\phi_s p)\|_{\phi_s p}} = \frac{\nu^\pm(p)}{\|\nu^\pm(p)\|_{\phi_s p}} = \frac{\|\nu^\pm(p)\|_p \sigma^\pm(p)}{\|\nu^\pm(p)\|_{\phi_s p}}.$$

Remark 3.12. — We recall that for $p_\gamma = (\gamma^+, \gamma^-, 0)$ we have $\phi_{l(\gamma)} p_\gamma = \gamma p_\gamma$ and hence $\|\nu^\pm(p_\gamma)\|_{\phi_{l(\gamma)} p_\gamma} = \lambda^\mp(\gamma) \|\nu^\pm(p_\gamma)\|_{p_\gamma}$. It follows that

$$\sigma^\pm(\gamma p_\gamma) = \lambda^\pm(\gamma) \sigma^\pm(p_\gamma).$$

Notation 3.13. — We consider the proto-neutralised sections and define

$$\nabla_\phi \sigma^\pm(p) := \left. \frac{\partial}{\partial s} \right|_{s=0} \sigma^\pm(\phi_s p) = - \left. \frac{\partial}{\partial s} \right|_{s=0} \log \|\nu^\pm(p)\|_{\phi_s p} \sigma^\pm(p).$$

Furthermore, for all $p \in \widetilde{\mathrm{U}\Gamma}$ we define:

$$f^\pm(p) := \frac{\langle \nabla_\phi \sigma^\pm(p) \mid \sigma^\mp(p) \rangle}{\langle \sigma^+(p) \mid \sigma^-(p) \rangle} = - \left. \frac{\partial}{\partial s} \right|_{s=0} \log \|\nu^\pm(p)\|_{\phi_s p}.$$

Remark 3.14. — As the action of Γ commutes with the flow and the proto-neutralised sections are Γ -equivariant, we obtain that the functions f^\pm are Γ -invariant i. e. $f^\pm(\gamma p) = f^\pm(p)$ for all $\gamma \in \Gamma$ and $p \in \widetilde{\mathrm{U}\Gamma}$. Hence, f^\pm induce functions on $\mathrm{U}\Gamma$. We abuse notation and denote these functions by f^\pm too.

DEFINITION 3.15. — *Let $f, g : \mathrm{U}\Gamma \rightarrow \mathbb{R}$ be two Hölder continuous functions. Then f is said to be Livšic cohomologous to g if there exists a function $h : \mathrm{U}\Gamma \rightarrow \mathbb{R}$ which is differentiable along the flow ϕ_t and satisfies the following property:*

$$f - g = \left. \frac{\partial}{\partial t} \right|_{t=0} h \circ \phi_t.$$

DEFINITION 3.16. — *The Livšic cohomology classes $[f_t^+]$ (resp. $[f_t^-]$), of the functions on $\mathrm{U}\Gamma$ mentioned in Remark 3.14, are called left (resp. right) proto-Labourie–Margulis invariants of ρ_t .*

PROPOSITION 3.17. — Suppose μ_γ is the flow invariant probability measure supported on the periodic orbit of $\mathbf{U}\Gamma$ corresponding to $\gamma \in \Gamma$ with period $l(\gamma)$ and $f^\pm : \mathbf{U}\Gamma \rightarrow \mathbb{R}$ is defined as above. Then

$$\int f^\pm d\mu_\gamma = \frac{\log \lambda^\pm(\gamma)}{l(\gamma)}.$$

Proof. — Suppose p_γ belongs to the periodic orbit corresponding to γ . We deduce that

$$\|\nu^\pm(p_\gamma)\|_{\phi_{l(\gamma)p_\gamma}} = \|\nu^\pm(p_\gamma)\|_{\gamma p_\gamma} = \|\rho(\gamma)^{-1}\nu^\pm(p_\gamma)\|_{p_\gamma}.$$

Hence, we obtain that

$$\int_0^{l(\gamma)} f^\pm(\phi_s p_\gamma) ds = \log \|\nu^\pm(p_\gamma)\|_{p_\gamma} - \log \|\nu^\pm(p_\gamma)\|_{\phi_{l(\gamma)p_\gamma}} = \log \lambda^\pm(\gamma),$$

and our result follows. \square

Remark 3.18. — As $\lambda^+(\gamma)\lambda^-(\gamma) = 1$, we obtain that

$$\int \frac{1}{2}(f^+ + f^-)d\mu_\gamma = 0 \text{ and } \int \frac{1}{2}(f^+ - f^-)d\mu_\gamma = \frac{\log \lambda^+(\gamma)}{l(\gamma)}$$

for all appropriate $\gamma \in \Gamma$. Hence, by Livšic's Theorem [57] we obtain that the functions $(f^+ + f^-)/2$ are Livšic cohomologous to zero and the functions $(f^+ - f^-)/2$ are Livšic cohomologous respectively to f^+ .

3.3. An equivalent criterion

In this subsection we prove an equivalent criterion for proper actions of hyperbolic groups Γ on $\mathbb{H}^{n,n-1}$ in terms of proto-Labourie–Margulis invariants.

LEMMA 3.19. — Suppose $p \in \widetilde{\mathbf{U}\Gamma}$ and σ^\pm are as in Subsection 3.2. Then there exists $h(p) \in \mathbf{G}$ such that $\nu^\pm(p) = \nu_\pm([h(p)])$, $h(\phi_s p) = h(p)$ and

$$\frac{\sigma^\pm(p)}{\sqrt{2\langle \sigma^+(p) \mid \sigma^-(p) \rangle}} = \frac{1}{2} \left(\sqrt{\frac{\|h(p)v_-\|_p}{\|h(p)v_+\|_p}} \right)^{\pm 1} h(p)v_\pm.$$

Moreover, there exists some $g(p) \in \mathbf{G}$ such that

$$\frac{(\sigma^+(p) - \sigma^-(p))}{\sqrt{2\langle \sigma^+(p) \mid \sigma^-(p) \rangle}} = g(p)e_{2n}.$$

Proof. — We observe that $\nu^\pm(p) = \nu_\pm([h(p)]) = (h(p)v_\pm)/\|h(p)v_\pm\|$ for some $h(p) \in \mathbf{G}$. We can choose $h(p)$ such that it does not depend on the flow. Hence, we obtain that

$$\sigma^\pm(p) = \frac{h(p)v_\pm}{\|h(p)v_\pm\|_p}, \quad \langle \sigma^+(p) \mid \sigma^-(p) \rangle = \frac{2}{\|h(p)v_+\|_p \|h(p)v_-\|_p},$$

and for $a = \sqrt{\|h(p)v_-\|_p^t / \|h(p)v_+\|_p^t}$ we deduce that

$$\frac{\sigma^\pm(p)}{\sqrt{2\langle \sigma^+(p) \mid \sigma^-(p) \rangle}} = \frac{1}{2} a^{\pm 1} h(p)v_\pm.$$

Now we can choose $h_a \in \mathbf{G}$ such that $h_a v_\pm = a^{\pm 1} v_\pm$ and consider $g(p) = h(p)h_a$ to conclude our result. \square

Remark 3.20. — As $\text{Stab}_{\mathbf{G}}(e_{2n}) = \mathbf{G}_0$, we obtain that $\mathbb{H}^{n,n-1} := \mathbf{G}e_{2n} \cong \mathbf{G}/\mathbf{G}_0$. We define

$$\sigma := \frac{(\sigma^+ - \sigma^-)}{\sqrt{2\langle \sigma^+ \mid \sigma^- \rangle}}$$

and observe that $\sigma : \widetilde{\mathbf{U}\Gamma} \rightarrow \mathbb{H}^{n,n-1}$. As σ^\pm are Γ -equivariant, we deduce that $\sigma(\gamma p) = \rho(\gamma)\sigma(p)$ for all $p \in \widetilde{\mathbf{U}\Gamma}$.

Suppose $s \in \mathbb{R}$, $\gamma \in \Gamma$ and $(p, x) \in \widetilde{\mathbf{U}\Gamma} \times \mathbb{H}^{n,n-1}$. We observe that \mathbb{R} acts on $\widetilde{\mathbf{U}\Gamma} \times \mathbb{H}^{n,n-1}$ by sending (p, x) to $(\phi_s p, x)$ and Γ acts on $\widetilde{\mathbf{U}\Gamma} \times \mathbb{H}^{n,n-1}$ by sending (p, x) to $(\gamma p, \rho(\gamma)x)$.

LEMMA 3.21. — *The Γ action on $(\widetilde{\mathbf{U}\Gamma} \times \mathbb{H}^{n,n-1})/\mathbb{R}$ is proper if and only if the \mathbb{R} action on $\Gamma \backslash (\widetilde{\mathbf{U}\Gamma} \times \mathbb{H}^{n,n-1})$ is proper.*

Proof. — As the action of Γ and the action of \mathbb{R} on $\widetilde{\mathbf{U}\Gamma}$ commute with each other, we see that

$$\begin{aligned} \gamma \phi_t(p, x) &= \gamma(\phi_t p, x) = (\gamma \phi_t p, \rho(\gamma)x) \\ &= (\phi_t \gamma p, \rho(\gamma)x) = \phi_t(\gamma p, \rho(\gamma)x) = \phi_t \gamma(p, x). \end{aligned}$$

Now we use Lemma 5.2 of [35] (See also Lemma 3.1 of [6]) to conclude our result. \square

LEMMA 3.22 (Goldman–Labourie [34]). — *Suppose $f : \mathbf{U}\Gamma \rightarrow \mathbb{R}$ is a Hölder continuous function such that*

$$\int f d\mu > 0$$

for all ϕ invariant measures μ on $\mathbf{U}\Gamma$. Then f is Livšic cohomologous to a positive function.

Proof. — The proof follows verbatim the proof given in the proof of Lemma 3 of [34] once we replace all the appearances of a manifold by the compact metric space $\mathbf{U}\Gamma$. \square

PROPOSITION 3.23. — *Suppose $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ and the action of $\rho(\Gamma)$ on $\mathbb{H}^{n,n-1}$ is proper. Then the proto-Labourie–Margulis invariants of ρ are non-vanishing.*

Proof. — Let us assume on contrary that the action of $\rho(\Gamma)$ on $\mathbb{H}^{n,n-1}$ is proper but the proto-Labourie–Margulis invariants f^\pm of ρ are not non-vanishing. Then neither f^+ nor f^- is Livšic cohomologous to a strictly positive function. Hence using Lemma 3.22 and the fact that $(f^+ + f^-)/2$ is Livšic cohomologous to 0, we get that there exists a flow invariant probability measure μ such that

$$\int \frac{1}{2}(f^+ - f^-)d\mu = 0.$$

Suppose $T > 0$. We consider

$$f_T(p) := \frac{1}{T} \int_0^T \frac{1}{2}(f^+ - f^-)(\phi_s p) ds$$

and observe that $\int f_T d\mu = 0$. Hence for all $T > 0$ there exists $p_T \in \mathbf{U}\Gamma$ such that $f_T(p_T) = 0$. It follows that

$$\log \left(\frac{\|\nu^+(p_T)\|_{p_T}}{\|\nu^+(p_T)\|_{\phi_T p_T}} \right) - \log \left(\frac{\|\nu^-(p_T)\|_{p_T}}{\|\nu^-(p_T)\|_{\phi_T p_T}} \right) = 0.$$

We use Lemma 3.19 to obtain that

$$\frac{\|\nu^\pm(p_T)\|_{p_T}}{\|\nu^\pm(p_T)\|_{\phi_T p_T}} = \frac{\|h(p_T)v_\pm\|_{p_T}}{\|h(p_T)v_\pm\|_{\phi_T p_T}}$$

and hence for all $T > 0$, we have $\sigma_\rho(\phi_T p_T) = \sigma_\rho(p_T)$. Moreover, $\mathbf{U}\Gamma$ is compact. It follows that \mathbb{R} does not act properly on $\Gamma \backslash (\widetilde{\mathbf{U}\Gamma} \times \mathbb{H}^{n,n-1})$. Now we use Lemma 3.21 to get that Γ does not act properly on $(\widetilde{\mathbf{U}\Gamma} \times \mathbb{H}^{n,n-1})/\mathbb{R} = \partial_\infty \Gamma^{(2)} \times \mathbb{H}^{n,n-1}$ and hence Γ does not act properly on $\mathbb{H}^{n,n-1}$ a contradiction. \square

PROPOSITION 3.24. — *If $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ is such that the proto-Labourie–Margulis invariants of ρ are non-vanishing. Then the action of $\rho(\Gamma)$ on $\mathbb{H}^{n,n-1}$ is proper.*

Proof. — Suppose the action of $\rho(\Gamma)$ on $\mathbb{H}^{n,n-1}$ is not proper. Hence, there exist $\gamma_k \in \Gamma$ going to infinity and $x_k \in \mathbb{H}^{n,n-1}$ such that the sequence $\{x_k\}$ converge to some $x \in \mathbb{H}^{n,n-1}$ and $\{y_k := \rho(\gamma_k)x_k\}$ converge to some $y \in \mathbb{H}^{n,n-1}$. As $\rho(\Gamma)$ is Anosov with respect to \mathbf{P}^\pm , we use Theorem 1.7

of [40] and Remark 2.6 to get that $\rho(\Gamma)$ is AMS proximal. Hence without loss of generality we can assume that γ_k is of infinite order for all k and $\lim_{k \rightarrow \infty} \gamma_k^+ \neq \lim_{k \rightarrow \infty} \gamma_k^-$ (please see Section 8.1 of [33] for a more detailed version of this argument). Therefore, we obtain that the flowlines corresponding to γ_k converge to a flowline with boundary points $\lim_{k \rightarrow \infty} \gamma_k^\pm$ and

$$\lim_{k \rightarrow \infty} l(\gamma_k) = \infty.$$

We choose $q_k \in \widetilde{\text{U}\Gamma}$ such that q_k is a point on the flowline corresponding to γ_k and q_k converges to some $q \in \widetilde{\text{U}\Gamma}$. For $z \in \{x, y\}$ suppose $w_{z,k}^\pm \in \xi^\pm(\gamma_k^\pm)$ and $v_{z,k}^\pm \in \mathbb{R}\nu^\pm(q_k)$ are such that

$$z_k = w_{z,k}^+ + w_{z,k}^- + v_{z,k}^+ + v_{z,k}^-.$$

As $\lim_{k \rightarrow \infty} z_k = z$, we obtain that $\lim_{k \rightarrow \infty} w_{z,k}^\pm = w_z^\pm$ and $\lim_{k \rightarrow \infty} v_{z,k}^\pm = v_z^\pm$ for some finite w_z^\pm and v_z^\pm . Hence, for $z \in \{x, y\}$ it follows that,

$$\begin{aligned} \langle z \mid \nu^+(q) \rangle \langle z \mid \nu^-(q) \rangle &= \langle v_z^+ \mid v_z^- \rangle \langle \nu^+(q) \mid \nu^-(q) \rangle, \\ \langle w_z^+ \mid w_z^- \rangle + \langle v_z^+ \mid v_z^- \rangle &= \langle z \mid z \rangle = -1. \end{aligned}$$

Moreover, as ρ is P^\pm -Anosov, we deduce that $w_x^+ = 0$ and $w_y^- = 0$. Hence,

$$\langle x \mid \nu^+(q) \rangle \langle x \mid \nu^-(q) \rangle = -\langle \nu^+(q) \mid \nu^-(q) \rangle = \langle y \mid \nu^+(q) \rangle \langle y \mid \nu^-(q) \rangle.$$

We also observe that

$$\begin{aligned} \lim_{k \rightarrow \infty} \lambda^\pm(\gamma_k) &= \lim_{k \rightarrow \infty} \frac{\langle x_k \mid \nu^\pm(q_k) \rangle}{\langle x_k \mid \rho(\gamma_k)^{-1} \nu^\pm(q_k) \rangle} \\ &= \lim_{k \rightarrow \infty} \frac{\langle x_k \mid \nu^\pm(q_k) \rangle}{\langle \rho(\gamma_k) x_k \mid \nu^\pm(q_k) \rangle} = \frac{\langle x \mid \nu^\pm(q) \rangle}{\langle y \mid \nu^\pm(q) \rangle}. \end{aligned}$$

Therefore, we deduce that $\lim_{k \rightarrow \infty} \lambda^\pm(\gamma_k)$ are finite nonzero numbers. Suppose f^\pm are the proto-Labourie–Margulis invariants of ρ . As $\lim_{k \rightarrow \infty} l(\gamma_k) = \infty$, it follows that

$$\lim_{k \rightarrow \infty} \int_{\gamma_k} f^+ = \lim_{k \rightarrow \infty} \frac{\log \lambda^+(\gamma_k)}{l(\gamma_k)} = 0.$$

We also know that the space of flow invariant probability measures on $\text{U}\Gamma$ is weak* compact. Hence, there exists a flow invariant probability measure μ on $\text{U}\Gamma$ such that $\int f^+ d\mu = 0$. It follows that f^+ is not Livšic cohomologous to any non-vanishing function and we conclude our result. \square

3.4. Special Anosov representations

In this subsection we relate the existence of non-vanishing proto-Labourie–Margulis invariants with certain special Anosov representations in G .

Notation 3.25. — Suppose $f : \mathcal{U}\Gamma \rightarrow \mathbb{R}$. We define,

$$I_{a,b}(f)(p) := \frac{1}{b-a} \int_a^b f(\phi_s p) ds$$

LEMMA 3.26 (Goldman–Labourie [34]). — *Let $f : \mathcal{U}\Gamma \rightarrow \mathbb{R}$ be a Hölder continuous function which is differentiable along the flow lines of ϕ and*

$$\int f d\mu > 0$$

for all ϕ invariant measure μ . Then $I_{0,T}(f) > 0$ for some $T > 0$.

Proof. — The proof follows verbatim the proof given in the proof of Lemma 7 of [34] once we replace all the appearances of a manifold by the compact metric space $\mathcal{U}\Gamma$. \square

LEMMA 3.27. — *Suppose $\rho \in \text{Hom}(\Gamma, G, P^\pm)$ and $[f^\pm]$ are the proto-Labourie–Margulis invariants of ρ . Also, suppose f^+ is Livšic cohomologous to a strictly positive function. Then there exist positive constants C and k such that for all $t > 0$ and $p \in \mathcal{U}\Gamma$ the following hold:*

$$\begin{aligned} \|\nu^+(p)\|_{\phi_t p} &\leq C \exp(-kt) \|\nu^+(p)\|_p, \\ \|\nu^-(p)\|_{\phi_{-t} p} &\leq C \exp(-kt) \|\nu^-(p)\|_p. \end{aligned}$$

Proof. — Suppose f^+ is Livšic cohomologous to a strictly positive function. Hence $-f^-$ is also Livšic cohomologous to a strictly positive function. Then using Lemma 3.26 we get that there exists $c > 0$ such that $I_{0,c}(f^+) > 0$. Now as $\mathcal{U}\Gamma$ is compact we get that $I_{0,c}(f^+) > k^+ > 0$ for some k^+ . Hence for $t > 0$ and integer n such that $(n+1) > (t/c) \geq n$ we have

$$tI_{0,t}(f^+)(p) \geq (t - nc)I_{0,t-nc}(f^+)(p) + nck^+$$

As both $[0, c]$ and $\mathcal{U}\Gamma$ are compact, we obtain that

$$\min\{I_{0,t}(f^+)(p) \mid (t, p) \in [0, c] \times \mathcal{U}\Gamma\} > k' \text{ for some } k' < 0.$$

As $(t - nc) \in [0, c]$, we have $I_{0,t-nc}(f^+)(p) \geq k'$. Hence,

$$tI_{0,t}(f^+)(p) \geq (t - nc)k' + nck^+.$$

As $k' < 0$, we obtain that $(t - nc)k' > ck'$. We denote $\exp(k^+c - k'c)$ by C^+ and recalling the definition of f^+ deduce that

$$\|\nu^+(p)\|_{\phi_t p} \leq C^+ \exp(-k^+t) \|\nu^+(p)\|_p$$

for all $t > 0$ and for all $p \in \widetilde{\text{U}\Gamma}$. As $-f^-$ is also Livšic cohomologous to a strictly positive function, we can do a similar computation to obtain positive constants C^- and k^- such that

$$\|\nu^-(p)\|_{\phi_{-t} p} \leq C^- \exp(-k^-t) \|\nu^-(p)\|_p$$

for all $t > 0$ and for all $p \in \widetilde{\text{U}\Gamma}$. Finally, we choose $C := \max\{C^+, C^-\}$, $k := \min\{k^+, k^-\}$ and conclude our result. \square

Notation 3.28. — We denote $\text{SL}(2n, \mathbb{R})$ by SL and the stabilizer of the oriented n -dimensional subspace V_\pm inside SL by Q^\pm i. e.

$$\text{Q}^\pm := \text{Stab}_{\text{SL}} \left(\vec{V}_\pm \right).$$

Also, we consider the orientation on $\mathbb{R}\nu^\pm(p)$ coming from the direction of the vector $\nu^\pm(p)$ to be positive and denote this oriented line by $\vec{\mathbb{R}}\nu^\pm(p)$.

PROPOSITION 3.29. — *Suppose $\rho \in \text{Hom}(\Gamma, \text{G}, \text{P}^\pm)$ and the left proto-Labourie–Margulis invariant $[f^+]$ is positive. Then $\rho \in \text{Hom}(\Gamma, \text{SL}, \text{Q}^\pm)$.*

Proof. — Suppose $p = (p_+, p_-, t) \in \widetilde{\text{U}\Gamma}$ and $\xi^\pm : \partial_\infty \Gamma \rightarrow \text{G}/\text{P}^\pm$ are the limit maps of $\rho \in \text{Hom}(\Gamma, \text{G}, \text{P}^\pm)$. We define

$$\begin{aligned} \eta^\pm : \partial_\infty \Gamma &\longrightarrow \text{SL}/\text{Q}^\pm \\ p_\pm &\longmapsto \xi^\pm(p_\pm) \oplus \vec{\mathbb{R}}\nu^\pm(p). \end{aligned}$$

and let $\eta(p) := (\eta^+(p_+), \eta^-(p_-))$. We use Lemmas 3.6 and 3.7 to observe that η^\pm is well defined and

$$\eta(\widetilde{\text{U}\Gamma}) \subset \text{SL}/(\text{Q}^+ \cap \text{Q}^-).$$

Moreover, using properties of ν^\pm we obtain that η^\pm are $\rho(\Gamma)$ -equivariant. Therefore, to show that $\rho \in \text{Hom}(\Gamma, \text{SL}, \text{Q}^\pm)$ we only need to produce equivariant metrics with contraction properties on:

$$\text{T}_{\eta(p)}(\text{SL}/(\text{Q}^+ \cap \text{Q}^-)) \cong \text{T}_{\eta^+(p_+)}(\text{SL}/\text{Q}^+) \oplus \text{T}_{\eta^-(p_-)}(\text{SL}/\text{Q}^-).$$

As the underlying subspaces of $\eta^\pm(p_\pm)$ add up to \mathbb{R}^{2n} , we use Proposition 10.1 of [71] to obtain that

$$\text{T}_{\eta^\pm(p_\pm)}(\text{SL}/\text{Q}^\pm) \cong \text{Hom}(\eta^\pm(p_\pm), \eta^\mp(p_\mp)).$$

We use Corollary 3.3 of [33] and Lemma 3.27 to obtain equivariant metrics with contraction properties on $\eta^\pm(p_\pm) = \xi^\pm(p_\pm) \oplus \overrightarrow{\mathbb{R}}\nu^\pm(p)$. Indeed, we use the following recipe: for $v \in \xi^\pm(p_\pm)$ and $w \in \mathbb{R}\nu^\pm(p)$ we define,

$$\|(v, w)\|_p^2 := \|v\|_p^2 + \|w\|_p^2.$$

Now for $A \in \text{Hom}(\eta^\pm(p_\pm), \eta^\mp(p_\mp))$ we define

$$\|A\|_p = \sup_{v \neq 0} \frac{\|A(v, w)\|_p}{\|(v, w)\|_p}.$$

As $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ and the left proto-Labourie–Margulis invariant is positive, we obtain that $\|(v, w)\|_{\phi_{\pm t}p} \leq C e^{-kt} \|(v, w)\|_p$ and $\|A(v, w)\|_{\phi_{\mp t}p} \leq C e^{-kt} \|A(v, w)\|_p$ for all $t > 0$. Hence, for all $v \neq 0$ we deduce that

$$\frac{\|A(v, w)\|_{\phi_{\mp t}p}}{\|(v, w)\|_{\phi_{\mp t}p}} \leq C^2 e^{-2kt} \frac{\|A(v, w)\|_p}{\|(v, w)\|_p} \leq C^2 e^{-2kt} \|A\|_p$$

and take the supremum over all $v \neq 0$ in the left hand side to conclude that $\|A\|_{\phi_{\mp t}p} \leq C^2 e^{-2kt} \|A\|_p$ for all $t > 0$. \square

PROPOSITION 3.30. — Suppose $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm) \cap \text{Hom}(\Gamma, \mathbf{SL}, \mathbf{Q}^\pm)$. Then the left proto-Labourie–Margulis invariant $[f^+]$ is positive.

Proof. — Suppose $\xi^\pm : \partial_\infty \Gamma \rightarrow \mathbf{G}/\mathbf{P}^\pm$ and $\eta^\pm : \partial_\infty \Gamma \rightarrow \mathbf{SL}/\mathbf{Q}^\pm$ are the two pairs of limit maps. As $\rho(\Gamma) \subset \mathbf{G}$ and the action of $\rho(\gamma)$ (resp. $\rho(\gamma)^{-1}$) is contracting on $\eta^+(\gamma^+)$ (resp. $\eta^-(\gamma^-)$), we obtain that $\eta^\pm(\gamma^\pm)$ are isotropic subspaces of $\mathbb{R}^{n,n}$. As $\eta^\pm(\gamma^\pm)$ are of dimension n , it follows that they are maximally isotropic. Hence by continuity of the limit maps we deduce that the image of η^\pm lies inside the space of all maximal isotropic subspaces of $\mathbb{R}^{n,n}$. Also, as the action of $\rho(\gamma)$ (resp. $\rho(\gamma)^{-1}$) is contracting on $\xi^+(\gamma^+)$ (resp. $\xi^-(\gamma^-)$), we obtain that $\xi^\pm(\gamma^\pm) \subset \eta^\pm(\gamma^\pm)$. Now using the orientation on $\eta^\pm(p_\pm)$ and the orientations already on $\xi^\pm(p_\pm)$ we get positive orientations on $\eta^\pm(p_\pm) \cap (\xi^\mp(p_\mp))^\perp$. As $\eta^\pm(\gamma^\pm)$ are respectively in the orbits of V_\pm , we use Lemmas 2.15 and 2.16 to obtain that

$$\eta^\pm(p_\pm) \cap (\xi^\mp(p_\mp))^\perp = \overrightarrow{\mathbb{R}}\nu^\pm(p).$$

Moreover, we know that the contraction property of an Anosov representation does not depend on a particular choice of norms up to Hölder equivalence. Hence we can choose the collection of norms to be smooth along flow lines (please see Remark 3.10) and the contraction property would still hold. Now using Lemma 5.3 of [12] we can choose $C = 1$. Hence, we get that there exists a positive constant k and a collection $\{\|\cdot\|_p \mid p \in \widetilde{\text{U}\Gamma}\}$ of Euclidean norms on \mathbb{R}^{2n} such that:

- (1) it is Hölder continuous in the variable $p \in \widetilde{\text{U}\Gamma}$,

- (2) it is smooth along the flow lines of $\{\phi_s\}_{s \in \mathbb{R}}$,
- (3) it is equivariant i. e. $\|\rho(\gamma)v\|_{\gamma p} = \|v\|_p$ for all $v \in \mathbb{R}^{2n}$ and $\gamma \in \Gamma$,
- (4) it is contracting i. e. for all $p \in \widetilde{\text{U}\Gamma}$, $v \in \eta^+(p_+)$, $w \in \eta^-(p_-)$ and $s > 0$:

$$\frac{\|v\|_{\phi_s p}}{\|w\|_{\phi_s p}} \leq \exp(-ks) \frac{\|v\|_p}{\|w\|_p}.$$

Therefore, for all $p \in \widetilde{\text{U}\Gamma}$ we obtain

$$\begin{aligned} g(p) &:= \left. \frac{\partial}{\partial s} \right|_{s=0} \log \frac{\|\nu^+(p)\|_{\phi_s p}}{\|\nu^-(p)\|_{\phi_s p}} = \lim_{s \rightarrow 0} \frac{1}{s} \log \left(\frac{\|\nu^+(p)\|_{\phi_s p}}{\|\nu^-(p)\|_{\phi_s p}} \frac{\|\nu^-(p)\|_p}{\|\nu^+(p)\|_p} \right) \\ &\leq \lim_{s \rightarrow 0} \frac{1}{s} \log (\exp(-ks)) = -k. \end{aligned}$$

Hence for all flow invariant probability measure μ on $\text{U}\Gamma$ we get that

$$\int \frac{1}{2} (f^+ - f^-) d\mu = - \int g d\mu \geq k > 0.$$

Finally, using Remark 3.18 we conclude our result. \square

Remark 3.31. — Suppose $\rho \in \text{Hom}(\Gamma, \text{SL})$. We observe that

$$JQ^\pm J = \text{Stab}_G(J\vec{V}_\pm).$$

- (1) If $\rho \in \text{Hom}(\Gamma, G, \mathbb{P}^\pm)$, then we use Lemmas 2.15 and 2.16 to deduce that the left proto-Labourie–Margulis invariant $[f^+]$ is negative if and only if $\rho \in \text{Hom}(\Gamma, \text{SL}, JQ^\pm J)$.
- (2) If $\rho(\Gamma) \subset G$, then by Lemma 2.15 we deduce that $\rho \in \text{Hom}(\Gamma, \text{SL}, Q^\pm)$ (resp. $\text{Hom}(\Gamma, \text{SL}, JQ^\pm J)$) if and only if ρ is Anosov with respect to the stabilizers of V_\pm (resp. JV_\pm).
- (3) If ρ is Anosov in SL with respect to the stabilizers of a pair of transverse n dimensional subspaces and $\rho(\Gamma) \subset G$, then we use Lemma 2.16 to deduce that either ρ is Anosov with respect to the stabilizers of V_\pm or ρ is Anosov with respect to the stabilizers of JV_\pm .

THEOREM 3.32. — *Suppose $\rho \in \text{Hom}(\Gamma, G, \mathbb{P}^\pm)$. Then the action of $\rho(\Gamma)$ on $\mathbb{H}^{n, n-1}$ is proper if and only if ρ is Anosov in SL with respect to the stabilizer of an n -dimensional subspace.*

Proof. — We use Propositions 3.23, 3.24, 3.29, 3.30 and Remark 3.31 to conclude our result. \square

4. Infinitesimal proper actions

In this section we show that proper affine actions on $\mathbb{R}^{n,n-1}$ can be seen as infinitesimal versions of proper actions on $\mathbb{H}^{n,n-1}$. We do this by recalling the notion of affine Anosov representations from [33] and relating them with the special Anosov representations appearing in Section 3.

4.1. Affine Anosov representations

In this subsection we recall the notion of an affine Anosov representation. Suppose $G_a := G_0 \ltimes \mathbb{R}^{2n-1}$. We call an element $h \in G_0$ *pseudo-hyperbolic* if the unit eigenspace of h is one dimensional and h does not have -1 as an eigenvalue. We call an element $g = (h, u) \in G_a$ *pseudo-hyperbolic* if its linear part h is pseudo-hyperbolic.

We call an n dimensional subspace $V \subset \mathbb{R}^{n,n-1}$ a *null space* if $(V^\perp) \cap \mathbb{R}^{2n-1}$ is a maximal isotropic subspace of $\mathbb{R}^{n,n-1}$. Affine subspaces which are parallel to null spaces are called *affine null spaces*. We recall Lemmas 2.12 and 2.15 and note that we can consistently provide the maximal isotropic spaces and the null spaces with positive orientations.

Now suppose $g = (h, u) \in G_a$ is such that h is pseudo-hyperbolic. Let

$$W_\pm^h := \left\{ v \mid \lim_{k \rightarrow \infty} h^{\mp k} v = 0 \right\} \subset \mathbb{R}^{2n-1}.$$

We note that W_\pm^h are maximal isotropic subspaces of $\mathbb{R}^{n,n-1}$ and let v_0^h be the unique eigenvector of h with eigenvalue 1 such that $\langle v_0^h \mid v_0^h \rangle = 1$ and which is positively oriented with respect to the orientations on W_+^h and $(W_+^h)^\perp \cap \mathbb{R}^{n,n-1}$ (for more details please see [1]). Then the *Margulis invariant* of g is defined as:

$$\alpha(g) := \langle u \mid v_0^h \rangle.$$

Let Γ be a word hyperbolic group and let $(\rho, u) : \Gamma \rightarrow G_a$ be an injective homomorphism such that the linear part ρ is P_0^\pm -Anosov with limit maps given by

$$\xi_\rho^\pm : \partial_\infty \Gamma \longrightarrow G_0/P_0^\pm$$

and $\xi_\rho(p) = (\xi_\rho^+(p_+), \xi_\rho^-(p_-))$ for all $p = (p_+, p_-, t) \in \widetilde{U}\Gamma$.

Remark 4.1. — Suppose $(\rho, u) : \Gamma \rightarrow G_a$ is an injective homomorphism such that $\rho \in \text{Hom}(\Gamma, G, P_0^\pm)$ and $\gamma^\pm \in \partial_\infty \Gamma$ are respectively the attracting and repelling points of the action of any infinite order element $\gamma \in \Gamma$.

We observe that $\xi_\rho(\gamma^+, \gamma^-, t)$ is independent of t . Henceforth, we denote $\nu(\xi_\rho(\gamma^+, \gamma^-, t))$ by $\nu_\rho(\gamma^+, \gamma^-)$ and observe that $\nu_\rho(\gamma^+, \gamma^-) = v_0^{\rho(\gamma)}$. It follows that

$$\alpha((\rho, u)(\gamma)) = \langle u_\rho(\gamma) \mid \nu_\rho(\gamma^+, \gamma^-) \rangle.$$

DEFINITION 4.2. — *Let Γ be a word hyperbolic group and let (ρ, u) be an injective homomorphism from Γ to G_a such that $\rho \in \text{Hom}(\Gamma, G_0, P_0^\pm)$. Then the Labourie–Margulis invariant of this representation is a Livšic cohomology class $[f_{\rho, u}]$ of Hölder continuous functions $f_{\rho, u}$ such that*

$$\int f_{\rho, u} d\mu_\gamma = \frac{\alpha(\rho, u)(\gamma)}{l(\gamma)}$$

where μ_γ is a flow invariant probability measure supported on the periodic orbit of $U\Gamma$ corresponding to γ and $l(\gamma)$ is the period of this orbit.

Remark 4.3. — The existence of the Labourie–Margulis invariants follow from the constructions in Appendix 8.2 and Lemma 7.2 of [33]. Moreover, the uniqueness follows from Livšic’s theorem [57].

Now we define the notion of affine Anosov representations. Suppose $W_\pm \subset \mathbb{R}^{2n-1}$ be as in the previous section. We observe that

$$\text{Stab}_{G_0}(W_\pm) = \text{Stab}_{G_0}(W_\pm^\perp \cap \mathbb{R}^{2n-1}).$$

Henceforth, in this section we treat $(W_\pm^\perp \cap \mathbb{R}^{2n-1})$ as affine subspaces in \mathbb{R}^{2n-1} and call the stabilizers,

$$P_a^\pm := \text{Stab}_{G_a}(W_\pm^\perp \cap \mathbb{R}^{2n-1}),$$

of these subspaces under the action of the affine group G_a as *pseudo parabolic* subgroups. These subgroups of G_a are used in the definition of an affine Anosov representation in the same way as parabolic subgroups are used in the definition of an Anosov representation. We observe that for L_0 and e_n as defined in the previous section we have

$$L_a := P_a^+ \cap P_a^- = L_0 \ltimes \mathbb{R}e_n.$$

Let \mathcal{X}_a be the space of all affine null subspaces in \mathbb{R}^{2n-1} and let \mathcal{Y}_a be the space of all transverse pairs of affine null subspaces. Then \mathcal{Y}_a is an open and dense subset of $\mathcal{X}_a \times \mathcal{X}_a$. The group G_a acts transitively on the space \mathcal{X}_a and we have $\mathcal{X}_a \cong G_a/P_a^\pm$. Moreover, the diagonal action of G_a is transitive on \mathcal{Y}_a and we have $\mathcal{Y}_a \cong G_a/L_a$.

DEFINITION 4.4. — *Let $(\rho, u) : \Gamma \rightarrow G_a$ be an injective homomorphism. Then (ρ, u) is called affine Anosov with respect to P_a^\pm if and only if the following conditions hold:*

- (1) (a) *There exist a continuous, injective, $(\rho, u)(\Gamma)$ -equivariant limit maps $\xi^\pm : \partial_\infty \Gamma \rightarrow \mathbf{G}_a / \mathbf{P}_a^\pm$ such that $\xi(p) := (\xi^+(p_+), \xi^-(p_-)) \in \mathbf{G}_a / \mathbf{L}_a$ for $p = (p_+, p_-, t) \in \widetilde{\mathbf{U}\Gamma}$.*
- (b) *There exist positive constants C, c and for $p \in \widetilde{\mathbf{U}\Gamma}$ a continuous collection of $(\rho, u)(\Gamma)$ -equivariant Euclidean metrics $\|\cdot\|_p$ on $\mathbf{T}_{\xi(p)}(\mathbf{G}_a / \mathbf{L}_a)$ such that for all $v^\pm \in \mathbf{T}_{\xi^\pm(p_\pm)}(\mathbf{G}_a / \mathbf{P}_a^\pm)$ and $t \geq 0$:*

$$\|v^\pm\|_{\phi_{\pm t} p} \leq C e^{-ct} \|v^\pm\|_p.$$

- (2) *There exists a $(\rho, u)(\Gamma)$ -equivariant map $s : \widetilde{\mathbf{U}\Gamma} \rightarrow \mathbb{R}^{2n-1}$ which is Hölder continuous and is differentiable along the flow lines of ϕ . Moreover, for all $p \in \widetilde{\mathbf{U}\Gamma}$ the function*

$$f(p) := \left\langle \frac{\partial}{\partial t} \Big|_{t=0} s(\phi_t p) \Big| \nu_\rho(p) \right\rangle \neq 0.$$

Remark 4.5. — We note that whenever the first condition of the above definition is satisfied one can use a partition of unity type argument to guarantee the existence of a $(\rho, u)(\Gamma)$ -equivariant map $s : \widetilde{\mathbf{U}\Gamma} \rightarrow \mathbb{R}^{2n-1}$ which is Hölder continuous and is differentiable along the flow lines of ϕ (for more details please see the Appendix 8.2 of [33]).

We denote the space of all representations in \mathbf{G}_a which are affine Anosov with respect to \mathbf{P}_a^\pm by $\mathbf{Hom}(\Gamma, \mathbf{G}_a, \mathbf{P}_a^\pm)$.

Remark 4.6. — Suppose $(\rho, u) \in \mathbf{Hom}(\Gamma, \mathbf{G}_a)$. Then we note that (ρ, u) satisfy Condition (1) in the above definition if and only if $\rho \in \mathbf{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$. Moreover, Condition (2) in the above definition is equivalent to saying that the Labourie–Margulis invariant of (ρ, u) is non-vanishing.

We state the key property of affine Anosov representations:

THEOREM 4.7 (Ghosh–Treib [33]). — *Suppose $\rho \in \mathbf{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$ and $(\rho, u) \in \mathbf{Hom}(\Gamma, \mathbf{G}_a)$. Then $(\rho, u) \in \mathbf{Hom}(\Gamma, \mathbf{G}_a, \mathbf{P}_a^\pm)$ if and only if the action of $(\rho, u)(\Gamma)$ on \mathbb{R}^{2n-1} is proper.*

It is important to mention here the existence and non-existence results due to Abels–Margulis–Soifer [1]:

THEOREM 4.8 (Abels–Margulis–Soifer [1]). — *The following holds:*

- (1) *There exist free subgroups of $\mathbf{Aff}(2n-1, \mathbb{R})$ with linear part Zariski dense in $\mathbf{SO}(n, n-1)$ which act properly discontinuously on \mathbb{R}^{2n-1} , when n is even.*
- (2) *There does not exist any subgroup of $\mathbf{Aff}(2n-1, \mathbb{R})$ with linear part Zariski dense in $\mathbf{SO}(n, n-1)$ which acts properly discontinuously on \mathbb{R}^{2n-1} , when n is odd.*

4.2. Margulis invariants as derivatives

In this subsection we relate Margulis invariants with derivatives of certain eigenvalues.

Remark 4.9. — Suppose \mathfrak{g} (resp. \mathfrak{g}_0) denote the Lie algebra of G (resp. G_0) and $G \in \mathfrak{g}$. We observe that $G \in \mathfrak{g}_0$ if and only if $Ge_{2n} = 0$.

As $\langle ge_{2n} \mid ge_{2n} \rangle$ is constant, we obtain that $\langle Ge_{2n} \mid e_{2n} \rangle = 0$ for all $G \in \mathfrak{g}$. Hence, $Ge_{2n} \in \mathbb{R}^{n,n-1}$. Moreover, for $v \in \mathbb{R}^{n,n-1} \subset \mathbb{R}^{n,n}$ we consider

$$G = \begin{bmatrix} 0 & v \\ v^t I_{n,n-1} & 0 \end{bmatrix}$$

and observe that $G \in \mathfrak{g}$ with $Ge_{2n} = v$.

Remark 4.10. — Suppose $g \in G_0$ is a pseudo-hyperbolic element and W_{\pm}^g be the attracting and repelling subspaces of g (see Section 4.1). Observe that the dimensions of W_{\pm}^g are $(n-1)$. Also, suppose $g_t \in G$ is an analytic one parameter family with $g = g_0$. Then, for t small enough we obtain a pair of attracting and repelling subspaces of g_t denoted by $W_{\pm}^{g_t}$ and whose dimensions are also $(n-1)$. Moreover, using Lemma 2.16 we obtain that there is a unique maximal isotropic space inside $(W_+^{g_t})^{\perp}$ which is in the orbit of V_+ (see Section 3.1). We denote this space by $V_+^{g_t}$. We observe that g_t preserves the line $(W_-^{g_t})^{\perp} \cap V_+^{g_t}$. Let $\lambda(g_t)$ be the eigenvalue of the action of g_t on this line. We observe that the definition of λ given here is compatible with the definition of λ^+ given in Notation 3.9.

Finally, we start using the following notation:

$$d\lambda(g, G) := \left. \frac{d}{dt} \right|_{t=0} \lambda(g_t).$$

LEMMA 4.11. — Suppose $g \in G_0$ is a pseudo-hyperbolic element and $g_t \in G$ is an analytic one parameter family whose tangent direction at $g = g_0$ is G and $Ge_{2n} = v$. Then

$$\alpha(g, v) = d\lambda(g, G).$$

Proof. — We consider the orientation on $V_{\pm}^{g_t}$ coming from \vec{V}_{\pm} and the orientation on $W_{\pm}^{g_t}$ coming from \vec{W}_{\pm} . Let $v_t^{\pm} \in (W_{\mp}^{g_t})^{\perp} \cap V_{\pm}^{g_t}$ respectively be the $\|\cdot\|$ -unit vectors which is positively oriented with respect to the orientations on $V_{\pm}^{g_t}$ and $W_{\pm}^{g_t}$. Hence,

$$\lambda(g_t) = \langle g_t^{\pm 1} v_t^{\pm} \mid v_t^{\mp} \rangle / \langle v_t^+ \mid v_t^- \rangle.$$

We take derivative on both sides and deduce that

$$\left. \frac{d}{dt} \right|_{t=0} \lambda(g_t) = \langle \pm G v_0^{\pm} \mid v_0^{\mp} \rangle / \langle v_0^+ \mid v_0^- \rangle.$$

We observe that $e_{2n} = (av_0^+ - bv_0^-)$ for some $a, b \in \mathbb{R}$. As $\langle e_{2n} \mid e_{2n} \rangle = -1$, we obtain that $2ab\langle v_0^+ \mid v_0^- \rangle = 1$. It follows that

$$d\lambda(g, G) = 2ab\langle Gv_0^+ \mid v_0^- \rangle.$$

Moreover, as $\langle g_tv_t^\pm \mid v_t^\pm \rangle = 0$, we obtain that $\langle Gv_0^\pm \mid v_0^\pm \rangle = 0$. Therefore,

$$d\lambda(g, G) = \langle G(av_0^+ - bv_0^-) \mid (av_0^+ + bv_0^-) \rangle = \langle v \mid (av_0^+ + bv_0^-) \rangle.$$

Finally, we conclude by observing that $(av_0^+ + bv_0^-) = v_0^g$. \square

Remark 4.12. — If $\alpha(g, v) \neq 0$, then for t small enough $\lambda(g_t) \neq \lambda^{-1}(g_t)$.

PROPOSITION 4.13. — *Suppose $\rho : \Gamma \rightarrow G_0$ is Anosov with respect to P_0^\pm . Then ρ is Anosov in G with respect to P^\pm i. e. $\text{Hom}(\Gamma, G_0, P_0^\pm) \subset \text{Hom}(\Gamma, G, P^\pm)$.*

Proof. — Let $\rho : \Gamma \rightarrow G_0$ be Anosov with respect to P_0^\pm with limit map

$$\xi : \widetilde{U\Gamma} \longrightarrow G_0/L_0.$$

We recall that $G_0/L_0 \subset G/L$. Hence we get limit maps

$$\xi : \widetilde{U\Gamma} \longrightarrow G_0/L_0 \subset G/L.$$

Therefore, to show that $\rho \in \text{Hom}(\Gamma, G, P^\pm)$ we need only to show that the contraction properties hold true.

We observe that $\xi^\pm(p_\pm)^\perp = (\xi^\pm(p_\pm)^\perp \cap \mathbb{R}^{2n-1}) \oplus \mathbb{R}e_{2n}$,

$$T_{(\xi^+(p_+), \xi^-(p_-))}(G_0/L_0) = T_{\xi^+(p_+)}(G_0/P_0^+) \oplus T_{\xi^-(p_-)}(G_0/P_0^-),$$

$$T_{(\xi^+(p_+), \xi^-(p_-))}(G/L) = T_{\xi^+(p_+)}(G/P^+) \oplus T_{\xi^-(p_-)}(G/P^-),$$

Hence, using Proposition 10.1 of [71] we obtain that

$$\begin{aligned} T_{\xi^\pm(p_\pm)}(G/P^\pm) &\cong \text{Hom}_{\text{skew}}(\xi^\pm(p_\pm), \xi^\mp(p_\mp)^\perp) \\ &\cong \text{Hom}_{\text{skew}}(\xi^\pm(p_\pm), \xi^\mp(p_\mp)^\perp \cap \mathbb{R}^{2n-1}) \oplus \text{Hom}(\xi^\pm(p_\pm), \mathbb{R}e_{2n}) \\ &\cong T_{\xi^\pm(p_\pm)}(G_0/P_0^\pm) \oplus \text{Hom}(\xi^\pm(p_\pm), \mathbb{R}e_{2n}). \end{aligned}$$

As $\rho \in \text{Hom}(\Gamma, G_0, P_0^\pm)$, we have a collection of norms satisfying contraction properties on $T_{\xi^\pm(p_\pm)}(G_0/P_0^\pm)$. We use Corollary 3.3 of [33] to obtain a norm satisfying contraction properties on $\text{Hom}(\xi^\pm(p_\pm), \mathbb{R}e_{2n})$. Finally, we obtain equivariant norms satisfying contraction properties on $T_{\xi^\pm(p_\pm)}(G/P^\pm)$ by using the following recipe:

$$\|(v, w)\|^2 := \|v\|^2 + \|w\|^2$$

where $v \in T_{\xi^\pm(p_\pm)}(G_0/P_0^\pm)$ and $w \in \text{Hom}(\xi^\pm(p_\pm), \mathbb{R}v_0)$. \square

Let $\{\rho_t\}_{t \in (-1,1)} \subset \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ be an analytic one parameter family and $U : \Gamma \rightarrow \mathfrak{g}$ be such that $\rho_0 = \rho$ and for all $\gamma \in \Gamma$,

$$U(\gamma) = \left. \frac{d}{dt} \right|_{t=0} \rho_t(\gamma) \rho(\gamma)^{-1}.$$

We call U a *tangent vector* of $\text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ at ρ . We observe that

$$U \in Z_{\text{Ad} \circ \rho}^1(\Gamma, \mathfrak{g}) := \{V \mid V(\gamma\eta) = \text{Ad}(\rho(\gamma))V(\eta) + V(\gamma) \text{ for all } \gamma, \eta \in \Gamma\}.$$

We consider, $u(\gamma) := U(\gamma)e_{2n}$ for all $\gamma \in \Gamma$ and observe that

$$u \in Z_\rho^1(\Gamma, \mathbb{R}^{2n-1}) := \{v \mid v(\gamma\eta) = \rho(\gamma)V(\eta) + V(\gamma) \text{ for all } \gamma, \eta \in \Gamma\}.$$

Remark 4.14. — We observe that for all $\gamma \in \Gamma$ and $p_\gamma = (\gamma^+, \gamma^-, 0)$, the action of $\rho(\gamma)$ fixes both $\nu_0^\pm(p_\gamma)$. Hence $\lambda_0^\pm(\gamma) = 1$ for all $\gamma \in \Gamma$.

PROPOSITION 4.15. — *Suppose $\rho \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm) \subset \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ and $U \in Z_{\text{Ad} \circ \rho}^1(\Gamma, \mathfrak{g})$ is a tangent vector to $\text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ at ρ . Then for $u = Ue_{2n}$,*

$$\alpha(\rho, u) = d\lambda^+(\rho, U).$$

Proof. — Follows directly from Lemma 4.11 □

Remark 4.16. — Suppose $\rho \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm) \subset \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$. Hence, for t small enough $\rho_t \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$. We use Remark 2.6 and Theorem 2.4 to obtain that the limit maps of ρ_t vary analytically along the variable t (for more details please see Theorem 6.1 of [12], Theorem 3.8 of [41] and Theorem 5.18 of [66]).

Suppose f_t^\pm be the functions whose Livšic cohomology classes are proto-Labourie–Margulis invariants and which are obtained via the construction of the collection of norms on \mathbb{R}^{2n} , indexed by $p \in \widetilde{\text{U}\Gamma}$ and ρ_t (see Remark 3.10. We recall that f_t^\pm are Hölder continuous in the variable p and vary analytically over a neighborhood of ρ inside the representation variety i. e. for some collection of Hölder continuous functions $\{h_n^\pm\}_{n=0}^\infty$ over $\text{U}\Gamma$ we have

$$f_t^\pm = \sum_{n=0}^{\infty} t^n h_n^\pm.$$

Henceforth, to simplify our notations we denote h_1^\pm by h^\pm and $\sum_{n=2}^\infty t^{n-2} h_n^\pm$ by h_t^\pm . We note that $h_0^\pm = f_0^\pm$ and h_t^\pm is analytic in the variable t .

PROPOSITION 4.17. — *Suppose $\{\rho_t\}_{t \in (-2,2)} \subset \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ be an analytic one parameter family with $\rho = \rho_0 \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$, U is the corresponding tangent vector to $\text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ at ρ and $u = Ue_{2n}$. Then the*

derivative of the left proto-Labourie–Margulis invariants of ρ_t at $t = 0$ is the Labourie–Margulis invariant of $(\rho, u) \in \text{Hom}(\Gamma, \mathbf{G}_a)$, i. e.

$$[f_{\rho, u}] = \left[\frac{\partial}{\partial t} \Big|_{t=0} f_t^+ \right].$$

Proof. — Suppose $h^+, h_t^+ : \text{U}\Gamma \rightarrow \mathbb{R}$ be as mentioned in Remark 4.16. Hence $f_t^+ = f_0^+ + th^+ + t^2 h_t^+$. It follows that for all flow invariant probability measures μ_γ supported on the closed orbits corresponding to infinite order elements $\gamma \in \Gamma$,

$$\int f_t^+ d\mu_\gamma = \int (f_0^+ + th^+ + t^2 h_t^+) d\mu_\gamma.$$

We use Proposition 3.17 and deduce that

$$\log \lambda_t^+(\gamma) - \log \lambda_0^+(\gamma) = l(\gamma)t \int h^+ d\mu_\gamma + l(\gamma)t^2 \int h_t^+ d\mu_\gamma.$$

We recall that $\text{U}\Gamma$ is compact, the functions $\{h_t^+\}_{t \in [-1, 1]}$ vary analytically in the variable t and $[-1, 1]$ is a compact set. Hence, for $t \in [-1, 1]$,

$$\int h_t^+ d\mu_\gamma \leq \max_{t \in [-1, 1]} \max_{p \in \text{U}\Gamma} |h_t^+(p)| = K \in \mathbb{R}.$$

It follows that $\lim_{t \rightarrow 0} t \int h_t^+ d\mu_\gamma = 0$. Also, we use Remark 4.14 to obtain

$$\int h^+ d\mu_\gamma = \frac{d\lambda^+((\rho, U)(\gamma))}{l(\gamma)\lambda_0^+(\gamma)} = \frac{\alpha((\rho, u)(\gamma))}{l(\gamma)}.$$

Finally, using Livšic’s theorem [57] we conclude our result. \square

Remark 4.18. — Similarly, we can show that the derivative of the right proto-Labourie–Margulis invariants of ρ_t at $t = 0$ is the negative of the Labourie–Margulis invariant of (ρ, u) .

Also, we recall that the negative of the Labourie–Margulis invariant of (ρ, u) is the Labourie–Margulis invariant of $(\rho, -u)$.

4.3. Margulis spacetimes and quotients of $\mathbb{H}^{n, n-1}$

In this subsection we relate elements of $\text{Hom}(\Gamma, \mathbf{G}_a)$ with deformations in $\text{Hom}(\Gamma, \mathbf{G})$ of elements in $\text{Hom}(\Gamma, \mathbf{G}_0)$. Moreover, we use this to relate proper affine actions on $\mathbb{R}^{n, n-1}$ with proper actions on $\mathbb{H}^{n, n-1}$.

LEMMA 4.19. — Suppose $\rho \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$ and

$$E : Z_{\text{Ad}\rho}^1(\Gamma, \mathfrak{g}) \longrightarrow Z_\rho^1(\Gamma, \mathbb{R}^{2n-1})$$

is the map which sends $U \in Z_{\text{Ad}\rho}^1(\Gamma, \mathfrak{g})$ to $u = Ue_{2n}$. Then E is surjective and the kernel of E is $Z_{\text{Ad}\rho}^1(\Gamma, \mathfrak{g}_0)$.

Proof. — Suppose $u \in Z_\rho^1(\Gamma, \mathbb{R}^{2n-1})$. We consider $U : \Gamma \rightarrow \mathfrak{g}$ such that,

$$U(\gamma) := \begin{bmatrix} 0 & u(\gamma) \\ u(\gamma)^t I_{n,n-1} & 0 \end{bmatrix}.$$

We observe that $U \in Z_{\text{Ad}_\rho}^1(\Gamma, \mathfrak{g})$. As $Ue_{2n} = u$, we deduce E is surjective.

Suppose $U \in Z_{\text{Ad}_\rho}^1(\Gamma, \mathfrak{g})$ is such that $Ue_{2n} = 0$. As $U(\gamma) \in \mathfrak{g}$, we compute and conclude that $U(\gamma) \in \mathfrak{g}_0$ and our result follows. \square

LEMMA 4.20. — *Let $\{\rho_t\}_{t \in (-1,1)} \subset \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ be an analytic one parameter family with $\rho = \rho_0 \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$, U be the tangent vector to $\{\rho_t\}_{t \in (-1,1)}$ at ρ and $u = Ue_{2n}$. Suppose the Labourie–Margulis invariant of (ρ, u) is non-vanishing. Then there exists $\epsilon \in (0, 1)$ such that for all t with $|t| \in (0, \epsilon)$, the proto-Labourie–Margulis invariants of ρ_t are also non-vanishing.*

Proof. — We use Remark 4.16 to deduce the existence of Hölder continuous functions $h^\pm, h_t^\pm : \mathbf{U}\Gamma \rightarrow \mathbb{R}$ such that h_t^\pm vary analytically over t in a neighborhood of zero (i. e. $|t| < \epsilon$ for some $\epsilon \in (0, 1)$) with

$$f_t^\pm = f_0^\pm + th^\pm + t^2 h_t^\pm$$

and $[f_t^\pm]$ are respectively the proto-Labourie–Margulis invariants of ρ_t . We use Proposition 4.17 to deduce that h^\pm are Livšic cohomologous to non-vanishing functions. Without loss of generality suppose h^+ is Livšic cohomologous to a function $h > 0$. We define

$$f_t := th + t^2 h_t^+.$$

Moreover, let us consider $c > \max\{|h_s^+(p)| \mid p \in \mathbf{U}\Gamma, s \in [-\epsilon, \epsilon]\}$. Then for $|t| \in (0, \epsilon)$ with $c|t| < \min\{h(p) \mid p \in \mathbf{U}\Gamma\}$ we have

$$h + th_t^+ \geq h - |th_t^+| > h - c|t| > 0.$$

It follows that for all t with $|t| \in (0, \epsilon)$, the functions f_t are non-vanishing. We recall that h is Livšic cohomologous to h^+ and use Proposition 3.17 to deduce that

$$\int (f_t^+ - f_t) d\mu_\gamma = \int f_0 d\mu_\gamma + t \int (h^+ - h) d\mu_\gamma = \frac{\log 1}{l(\gamma)} + 0 = 0.$$

Hence, using Livšic’s Theorem [57] we obtain that $[f_t^+] = [f_t]$ and our result follows. \square

THEOREM 4.21. — *Let $\{\rho_t\}_{t \in (-1,1)}$ be an analytic one parameter family of representations of Γ in \mathbf{G} with $\rho_0(\Gamma) \subset \mathbf{G}_0$. Let U be the tangent vector to $\{\rho_t\}_{t \in (-1,1)}$ at $\rho = \rho_0$ and $u = Ue_{2n}$. Suppose $(\rho, u) \in \text{Hom}(\Gamma, \mathbf{G}_a, \mathbf{P}_a^\pm)$.*

Then there exists $\epsilon > 0$ such that for all t with $|t| \in (0, \epsilon)$, ρ_t is Anosov in $\mathrm{SL}(2n, \mathbb{R})$ with respect to the stabilizer of an n -dimensional subspace.

Proof. — As $(\rho, u) \in \mathrm{Hom}(\Gamma, \mathbf{G}_a, \mathbf{P}_a^\pm)$, we obtain that $\rho \in \mathrm{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$. We use Proposition 4.13 to obtain an $\epsilon > 0$ such that for all $|t| < 2\epsilon$, $\rho_t \in \mathrm{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$. Finally, we use Lemma 4.20, Proposition 3.29 and Remark 3.31 to conclude our result. \square

COROLLARY 4.22. — *Let $\{\rho_t\}_{t \in (-1,1)}$ be an analytic one parameter family of representations of Γ in \mathbf{G} with $\rho = \rho_0 \in \mathrm{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$, U be the tangent vector to $\{\rho_t\}_{t \in (-1,1)}$ at ρ and $u = Ue_{2n}$. Suppose (ρ, u) is a Margulis spacetime. Then there exists $\epsilon > 0$ such that for all t with $|t| \in (0, \epsilon)$, $\rho_t(\Gamma)$ acts properly on $\mathbb{H}^{n,n-1}$.*

Proof. — As (ρ, u) is a Margulis spacetime, $(\rho, u)(\Gamma)$ acts properly on \mathbb{R}^{2n-1} . We use Theorem 4.7 to obtain that $(\rho, u) \in \mathrm{Hom}(\Gamma, \mathbf{G}_a, \mathbf{P}_a^\pm)$. Finally, we use Theorems 4.21 and 3.32 to conclude our result. \square

COROLLARY 4.23. — *Suppose n is even. Then there exists a non-abelian free subgroup with finitely many generators inside \mathbf{G} which act properly on $\mathbb{H}^{n,n-1}$.*

Proof. — The result follows from Theorem B of [1] and Corollary 4.22. \square

Appendix

In this section we introduce affine crossratios corresponding to Margulis invariants and crossratios corresponding to the eigenvalues whose derivative give rise to Margulis invariants. Moreover, we relate these affine crossratios (resp. crossratios) with a limiting result corresponding the Margulis invariants (resp. eigenvalues).

A.1. Affine Crossratios and Margulis Invariants

In this subsection we define affine crossratios.

Suppose $\{V_i\}_{i=1}^4$ are four null vector subspaces of \mathbb{R}^{2n-1} which are mutually transverse to each other and $\{A_i\}_{i=1}^4$ are four affine subspaces in \mathbb{R}^{2n-1} such that V_i is respectively parallel to A_i . Moreover, for $i \neq j$ suppose $x_{i,j}$ be a point in $A_i \cap A_j$ and suppose $v_{i,j} := \nu(V_i^\perp, V_j^\perp)$.

LEMMA A.1. — Suppose V_i, V_j, V_k are three null vector subspaces of \mathbb{R}^{2n-1} which are mutually transverse to each other. Then $v_{i,j} = (-1)^{n-1}v_{j,i}$ and

$$\langle v_{i,j} \mid v_{i,k} \rangle = 1 = \langle v_{i,j} \mid v_{k,j} \rangle.$$

Proof. — As $\dim(V_i) = n$, $\dim(V_i^\perp) = (n-1)$ and $\langle v_{i,*} \mid v_{i,*} \rangle = 1$, we obtain that $(v_{i,j} - av_{i,k}) \in V_i^\perp$ for some non-zero constant a . It follows that

$$a = \langle v_{i,j} \mid v_{i,k} \rangle = a^{-1}.$$

Hence, $a^2 = 1$. As $\langle v_{i,j} \mid v_{i,j} \rangle = 1$, using continuity we conclude that $a = 1$.

We note that $-I$ and I lie in the same connected component of the orthogonal group $O(n)$ if and only if n is even. Hence, $I_{n,n-1} \in SO_0(n, n-1)$ for n odd and $-I_{n,n-1} \in SO_0(n, n-1)$ for n even. It follows that for n odd, $\nu(W_-, W_+) = \nu([I_{n,n-1}]) = \nu(W_+, W_-) = (-1)^{n-1}\nu(W_+, W_-)$ and for n even, $\nu(W_-, W_+) = \nu([-I_{n,n-1}]) = -\nu(W_+, W_-) = (-1)^{n-1}\nu(W_+, W_-)$. Our result follows. \square

We define the *affine crossratio* as

$$\beta_{1,2,3,4} = \beta(A_1, A_2, A_3, A_4) := \langle x_{1,3} - x_{2,4} \mid v_{1,4} - v_{2,3} \rangle.$$

In particular, for $n = 2$, null subspaces of $\mathbb{R}^{2,1}$ are planes which are tangent to the light-cone. Now given four affine subspaces of \mathbb{R}^3 which are mutually transverse to each other and whose underlying vector spaces are null subspaces, we obtain that their mutual intersections give us four affine lines. Then $x_{1,3}$ (resp. $x_{2,4}$) is a point on the line of intersection between the first (resp. second) and the third (resp. fourth) affine subspace and $v_{1,4}$ (resp. $v_{2,3}$) are vectors which are unit with respect to the bilinear form $\langle \cdot \mid \cdot \rangle$, parallel with the intersection between the first (resp. second) and the fourth (resp. third) affine subspace and whose directions are consistent with the choice made in Remark 3.3.

The above definition is well defined since using Lemma A.1 it follows that for all $a, b \in \mathbb{R}$, $\langle av_{1,3} - bv_{2,4} \mid v_{1,4} - v_{2,3} \rangle = 0$. We also observe that the following equality holds for any other points $x_{3,1} \in A_3 \cap A_1$ and $x_{4,2} \in A_4 \cap A_2$:

$$\langle x_{1,3} - x_{2,4} \mid v_{1,4} - v_{2,3} \rangle = \langle x_{3,1} - x_{4,2} \mid v_{1,4} - v_{2,3} \rangle.$$

Now for $i \neq j$ we consider the following decomposition:

$$\mathbb{R}^{2n-1} = V_i^\perp \oplus V_j^\perp \oplus (V_i \cap V_j).$$

Let x_i^j be the projection of $x_{i,j}$ on V_i^\perp with respect to this decomposition. We observe that as $x_{i,j}$ varies along $A_i \cap A_j$ the projection x_i^j stays fixed.

Moreover, $x_i^j + V_j = A_j$. Using these observations we obtain:

$$\begin{aligned}\beta_{1,2,3,4} &= \langle x_{1,3} - x_{2,4} \mid v_{1,4} - v_{2,3} \rangle = \langle x_3^1 + x_1^3 - x_2^4 - x_4^2 \mid v_{1,4} - v_{2,3} \rangle \\ &= \langle x_3^1 \mid v_{1,4} \rangle - \langle x_1^3 \mid v_{2,3} \rangle - \langle x_2^4 \mid v_{1,4} \rangle + \langle x_4^2 \mid v_{2,3} \rangle \\ &= \langle x_3^1 \mid v_{1,4} - v_{1,3} \rangle - \langle x_1^3 \mid v_{2,3} - v_{1,3} \rangle - \langle x_2^4 \mid v_{1,4} - v_{2,4} \rangle \\ &\quad + \langle x_4^2 \mid v_{2,3} - v_{2,4} \rangle.\end{aligned}$$

Hence for any $x_i \in A_i$ we obtain the following identity:

$$\begin{aligned}(\text{A.1}) \quad \beta_{1,2,3,4} &= \langle x_1 \mid v_{1,4} - v_{1,3} \rangle + \langle x_2 \mid v_{2,3} - v_{2,4} \rangle \\ &\quad + \langle x_3 \mid v_{1,3} - v_{2,3} \rangle + \langle x_4 \mid v_{2,4} - v_{1,4} \rangle.\end{aligned}$$

PROPOSITION A.2. — *Let β be defined as above. Then for any five affine null spaces $A_*, \{A_i\}_{i=1}^4$ which are mutually transverse to each other and for any $(g, u) \in G_a$ the following identities hold:*

- (1) $\beta((g, u)A_1, (g, u)A_2, (g, u)A_3, (g, u)A_4) = \beta(A_1, A_2, A_3, A_4)$,
- (2) $\beta_{1,2,3,4} = \beta_{2,1,4,3} = (-1)^n \beta_{3,4,1,2} = (-1)^n \beta_{4,3,2,1}$,
- (3) $\beta_{1,2,3,4} + \beta_{1,2,4,3} = 0$,
- (4) $\beta_{1,*,3,4} + \beta_{*,2,3,4} = \beta_{1,2,3,4}$.

Moreover, for n even, $\beta_{1,2,3,4} + \beta_{1,3,4,2} + \beta_{1,4,2,3} = 0$.

Proof. — We use the definition of β to deduce that for all $(g, u) \in G_a$,

$$\beta((g, u)A_1, (g, u)A_2, (g, u)A_3, (g, u)A_4) = \beta(A_1, A_2, A_3, A_4).$$

We recall that $\nu(V, W) = (-1)^{n-1} \nu(W, V)$. Moreover, exploiting the symmetries in the definition of β we obtain the identity (2). Now interchanging A_3 and A_4 in the identity (A.1) and adding them up we obtain that $\beta_{1,2,3,4} + \beta_{1,2,4,3} = 0$. Suppose A_* is another affine null space which is mutually transverse with the other null spaces $\{A_i\}_{i=1}^4$. We observe that

$$\begin{aligned}&\langle x_1 \mid v_{1,4} - v_{1,3} \rangle + \langle x_* \mid v_{*,3} - v_{*,4} \rangle + \langle x_3 \mid v_{1,3} - v_{*,3} \rangle \\ &\quad + \langle x_4 \mid v_{*,4} - v_{1,4} \rangle + \langle x_* \mid v_{*,4} - v_{*,3} \rangle + \langle x_2 \mid v_{2,3} - v_{2,4} \rangle \\ &\quad + \langle x_3 \mid v_{*,3} - v_{2,3} \rangle + \langle x_4 \mid v_{2,4} - v_{*,4} \rangle \\ &= \langle x_1 \mid v_{1,4} - v_{1,3} \rangle + \langle x_2 \mid v_{2,3} - v_{2,4} \rangle + \langle x_3 \mid v_{1,3} - v_{2,3} \rangle \\ &\quad + \langle x_4 \mid v_{2,4} - v_{1,4} \rangle.\end{aligned}$$

Therefore, we conclude that $\beta_{1,*,3,4} + \beta_{*,2,3,4} = \beta_{1,2,3,4}$.

Finally, for even n , we cyclically permute A_2, A_3, A_4 in the identity (A.1) and add them up to conclude our result. \square

PROPOSITION A.3. — Suppose $(g, u) \in \mathbf{G}_a$ be such that its action on the space of affine null subspaces has an attracting (resp. repelling) fixed point A_+ (resp. A_-) and A_\pm are transverse to each other. Then for any affine null space A which is transverse to both A_\pm the following holds:

- (1) $\beta(A_-, A_+, (g, u)A, A) = 2\alpha(g, u)$ when n is even,
- (2) $\beta(A_-, A_+, (g, u)A, A) = 0$ when n is odd.

Proof. — Suppose $h = (g, u)$ and x_\pm, x, x_h are any four points respectively in A_\pm, A and hA . We observe that $\langle x_\pm \mid v_{A_\pm, hA} \rangle = \langle g^{-1}x_\pm \mid v_{A_\pm, A} \rangle$. Hence,

$$\begin{aligned} \langle x_\pm \mid v_{A_\pm, A} - v_{A_\pm, hA} \rangle &= \langle x_\pm - g^{-1}x_\pm \mid v_{A_\pm, A} \rangle \\ &= \langle x_\pm - h^{-1}x_\pm - g^{-1}u \mid v_{A_\pm, A} \rangle. \end{aligned}$$

As $x_h = hx'$ for some $x' \in A$, we obtain $g^{-1}x_h = g^{-1}u + x'$. Suppose V is the underlying vector subspace of A . We observe that $x' - x \in V$ and $(v_{A_-, A} - v_{A_+, A}) \in V^\perp$. It follows that

$$\langle x_h \mid v_{A_-, hA} - v_{A_+, hA} \rangle = \langle g^{-1}u + x \mid v_{A_-, A} - v_{A_+, A} \rangle.$$

We use the identity (A.1) and deduce that

$$\beta(A_-, A_+, hA, A) = \langle x_- - h^{-1}x_- \mid v_{A_-, A} \rangle + \langle h^{-1}x_+ - x_+ \mid v_{A_+, A} \rangle.$$

Suppose V_\pm are the vector spaces which are respectively parallel to A_\pm . We recall that h fixes A_\pm and hence $(x_\pm - h^{-1}x_\pm) \in V_\pm$. On the other hand $(v_{A_\pm, A} - v_{A_\pm, A_\mp}) \in V_\pm^\perp$ and therefore we deduce that

$$\beta(A_-, A_+, (g, u)A, A) = \langle u \mid v_{A_-, A_+} - v_{A_+, A_-} \rangle.$$

As $v_{i,j} = (-1)^{n-1}v_{j,i}$, our result follows. \square

Remark A.4. — Suppose $\rho \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$ and $(\rho, u) \in \text{Hom}(\Gamma, \mathbf{G}_a)$. Hence, (ρ, u) admits limit maps $\xi^\pm : \partial_\infty \Gamma \rightarrow \mathbf{G}_a / \mathbf{P}_a^\pm$ which satisfy the first two properties of being an affine Anosov representation (for more details please see Proposition 5.3 of [33]). In general, it only fails to satisfy the third property. Hence, for all infinite order elements $\gamma \in \Gamma$ we obtain that the action of $(\rho, u)(\gamma)$ on \mathcal{X}_a has an attracting fixed point and a repelling fixed point. We abuse notation and let $\xi_{\rho, u}(\gamma^+)$ (resp. $\xi_{\rho, u}(\gamma^-)$) denote the attracting (resp. repelling) fixed point. Henceforth, we fix the representation (ρ, u) and omit the subscripts (ρ, u) from the notation of the Margulis invariants and the affine crossratios. Also, when there is no confusion of notation, for $a, b, c, d \in \partial_\infty \Gamma$ all distinct, we denote $\beta(\xi(a), \xi(b), \xi(c), \xi(d))$ by $\beta(a, b, c, d)$.

PROPOSITION A.5. — Suppose n is even, $\rho \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$ and $(\rho, u) \in \text{Hom}(\Gamma, \mathbf{G}_a)$. Also, suppose $\gamma, \eta \in \Gamma$ are two infinite order elements such that the four points $\gamma^\pm, \eta^\pm \in \partial_\infty \Gamma$ are distinct and the sequence $\{\gamma^m \eta^k\}_{m \in \mathbb{N}} \subset \Gamma$ contains a subsequence $\{\gamma^{n_i} \eta^k\}_{i \in \mathbb{N}}$ consisting only of infinite order elements. Then the following identity holds:

$$\begin{aligned} \beta(\eta^-, \gamma^-, \gamma^+, \eta^k \gamma^+) + \beta(\eta^+, \gamma^+, \gamma^-, \eta^{-k} \gamma^-) \\ = 2 \lim_{i \rightarrow \infty} [\alpha(\gamma^{n_i} \eta^k) - \alpha(\gamma^{n_i})] - 2\alpha(\eta^k). \end{aligned}$$

Proof. — Suppose ξ is the affine limit map as mentioned in the previous remark and A, B, C are affine null spaces such that A is transverse to both $\xi((\gamma^{n_i} \eta^k)^\pm)$, B is transverse to both $\xi(\gamma^\pm)$ and C is transverse to both $\xi(\eta^\pm)$. We use Proposition A.3 and obtain the following three identities:

$$\begin{aligned} 2\alpha(\gamma^{n_i} \eta^k) &= \beta(\xi((\gamma^{n_i} \eta^k)^-), \xi((\gamma^{n_i} \eta^k)^+), (\rho, u)(\gamma^{n_i} \eta^k)A, A), \\ 2\alpha(\gamma^{n_i}) &= \beta(\xi(\gamma^-), \xi(\gamma^+), (\rho, u)(\gamma^{n_i})B, B), \\ 2\alpha(\eta^k) &= \beta(\xi(\eta^-), \xi(\eta^+), (\rho, u)(\eta^k)C, C). \end{aligned}$$

We observe that $\lim_{i \rightarrow \infty} (\gamma^{n_i} \eta^k)^+ = \gamma^+$ and $\lim_{i \rightarrow \infty} (\gamma^{n_i} \eta^k)^- = \eta^{-k} \gamma^-$. Also, $\eta^\pm \neq \gamma^\pm$. It follows that $\eta^- \neq \lim_{i \rightarrow \infty} (\gamma^{n_i} \eta^k)^\pm$. Hence, we can choose $A = B = \xi(\eta^-)$, $C = \xi(\gamma^-)$. We use Proposition A.2(4) to obtain:

$$\begin{aligned} 2\alpha(\gamma^{n_i} \eta^k) &= \beta(\gamma^+, (\gamma^{n_i} \eta^k)^+, \gamma^{n_i} \eta^-, \eta^-) + \beta((\gamma^{n_i} \eta^k)^-, \gamma^+, \gamma^{n_i} \eta^-, \eta^-), \\ 2\alpha(\gamma^{n_i}) &= \beta((\gamma^{n_i} \eta^k)^-, \gamma^+, \gamma^{n_i} \eta^-, \eta^-) + \beta(\gamma^-, (\gamma^{n_i} \eta^k)^-, \gamma^{n_i} \eta^-, \eta^-). \end{aligned}$$

Also, using Proposition A.2(1) we deduce that

$$\begin{aligned} \beta(\gamma^+, (\gamma^{n_i} \eta^k)^+, \gamma^{n_i} \eta^-, \eta^-) &= \beta(\gamma^+, (\eta^k \gamma^{n_i})^+, \eta^-, \gamma^{-n_i} \eta^-), \\ 2\alpha(\eta^k) &= \beta(\eta^-, \eta^+, \gamma^-, \eta^{-k} \gamma^-). \end{aligned}$$

Hence, taking the limit and then using Proposition A.2(2) and (3) we obtain:

$$\begin{aligned} 2 \lim_{i \rightarrow \infty} [\alpha(\gamma^{n_i} \eta^k) - \alpha(\gamma^{n_i})] &= \beta(\gamma^+, \eta^k \gamma^+, \eta^-, \gamma^-) - \beta(\gamma^-, \eta^{-k} \gamma^-, \gamma^+, \eta^-) \\ &= \beta(\eta^-, \gamma^-, \gamma^+, \eta^k \gamma^+) + \beta(\eta^-, \gamma^+, \gamma^-, \eta^{-k} \gamma^-). \end{aligned}$$

Now we use Proposition A.2(4) to deduce that

$$\beta(\eta^-, \gamma^+, \gamma^-, \eta^{-k} \gamma^-) = \beta(\eta^-, \eta^+, \gamma^-, \eta^{-k} \gamma^-) + \beta(\eta^+, \gamma^+, \gamma^-, \eta^{-k} \gamma^-).$$

Finally, our result follows from combining the last three identities. \square

PROPOSITION A.6. — Suppose n is even, $\rho \in \text{Hom}(\Gamma, \mathbf{G}_0, \mathbf{P}_0^\pm)$ and $(\rho, u) \in \text{Hom}(\Gamma, \mathbf{G}_a)$. Suppose $\gamma, \eta \in \Gamma$ are two infinite order elements such that the four points $\gamma^\pm, \eta^\pm \in \partial_\infty \Gamma$ are distinct and the sequence

$\{\gamma^m \eta^m\}_{m \in \mathbb{N}} \subset \Gamma$ contains a subsequence $\{\gamma^{n_i} \eta^{n_i}\}_{i \in \mathbb{N}}$ consisting only of infinite order elements. Then the following identity holds:

$$\lim_{i \rightarrow \infty} (\alpha(\gamma^{n_i} \eta^{n_i}) - \alpha(\gamma^{n_i}) - \alpha(\eta^{n_i})) = \beta(\eta^-, \gamma^-, \gamma^+, \eta^+).$$

Proof. — Suppose ξ is the corresponding limit map and $\{A_i, B_i, C_i\}_{i \in \mathbb{N}}$ is a collection of affine null spaces such that A_i is transverse to $\xi((\gamma^{n_i} \eta^{n_i})^\pm)$, B_i is transverse to $\xi(\gamma^\pm)$ and C_i is transverse to $\xi(\eta^\pm)$. We use Proposition A.3 to obtain the following three identities:

$$\begin{aligned} 2\alpha(\gamma^{n_i} \eta^{n_i}) &= \beta(\xi((\gamma^{n_i} \eta^{n_i})^-), \xi((\gamma^{n_i} \eta^{n_i})^+), (\rho, u)(\gamma^{n_i} \eta^{n_i}) A_i, A_i), \\ 2\alpha(\gamma^{n_i}) &= \beta(\xi(\gamma^-), \xi(\gamma^+), (\rho, u)(\gamma^{n_i}) B_i, B_i), \\ 2\alpha(\eta^{n_i}) &= \beta(\xi(\eta^-), \xi(\eta^+), (\rho, u)(\eta^{n_i}) C_i, C_i). \end{aligned}$$

Suppose $D_i := (\rho, u)(\eta^{n_i}) A_i$. We use Proposition A.2(1) and (4) to deduce

$$\begin{aligned} &\beta(\xi((\gamma^{n_i} \eta^{n_i})^-), \xi((\gamma^{n_i} \eta^{n_i})^+), (\rho, u)(\gamma^{n_i}) D_i, (\rho, u)(\eta^{-n_i}) D_i) \\ &= \beta(\xi((\gamma^{n_i} \eta^{n_i})^-), \xi((\gamma^{n_i} \eta^{n_i})^+), (\rho, u)(\gamma^{n_i}) D_i, D_i) \\ &\quad + \beta(\xi((\gamma^{n_i} \eta^{n_i})^-), \xi((\gamma^{n_i} \eta^{n_i})^+), D_i, (\rho, u)(\eta^{-n_i}) D_i) \\ &= \beta(\xi((\gamma^{n_i} \eta^{n_i})^-), \xi((\gamma^{n_i} \eta^{n_i})^+), (\rho, u)(\gamma^{n_i}) D_i, D_i) \\ &\quad + \beta(\xi((\eta^{n_i} \gamma^{n_i})^-), \xi((\eta^{n_i} \gamma^{n_i})^+), (\rho, u)(\eta^{n_i}) D_i, D_i). \end{aligned}$$

Moreover, by applying Proposition A.2(4) twice we deduce that

$$\begin{aligned} &\beta(\xi((\gamma^{n_i} \eta^{n_i})^-), \xi((\gamma^{n_i} \eta^{n_i})^+), (\rho, u)(\gamma^{n_i}) D_i, D_i) \\ &= \beta(\xi((\gamma^{n_i} \eta^{n_i})^-), \xi(\gamma^-), (\rho, u)(\gamma^{n_i}) D_i, D_i) \\ &\quad + \beta(\xi(\gamma^+), \xi((\gamma^{n_i} \eta^{n_i})^+), (\rho, u)(\gamma^{n_i}) D_i, D_i) \\ &\quad + \beta(\xi(\gamma^-), \xi(\gamma^+), (\rho, u)(\gamma^{n_i}) D_i, D_i). \end{aligned}$$

Similarly, we also have

$$\begin{aligned} &\beta(\xi((\eta^{n_i} \gamma^{n_i})^-), \xi((\eta^{n_i} \gamma^{n_i})^+), (\rho, u)(\eta^{n_i}) D_i, D_i) \\ &= \beta(\xi((\eta^{n_i} \gamma^{n_i})^-), \xi(\eta^-), (\rho, u)(\eta^{n_i}) D_i, D_i) \\ &\quad + \beta(\xi(\eta^+), \xi((\eta^{n_i} \gamma^{n_i})^+), (\rho, u)(\eta^{n_i}) D_i, D_i) \\ &\quad + \beta(\xi(\eta^-), \xi(\eta^+), (\rho, u)(\eta^{n_i}) D_i, D_i). \end{aligned}$$

We observe that $\lim_{i \rightarrow \infty} (\gamma^{n_i} \eta^{n_i})^+ = \gamma^+$, $\lim_{i \rightarrow \infty} (\gamma^{n_i} \eta^{n_i})^- = \eta^-$ and also $\lim_{i \rightarrow \infty} (\eta^{n_i} \gamma^{n_i})^+ = \eta^+$, $\lim_{i \rightarrow \infty} (\eta^{n_i} \gamma^{n_i})^- = \gamma^-$. We recall that the four points η^\pm, γ^\pm are distinct. Let $x \in \partial_\infty \Gamma$ be such that it is distinct from all

the following four points: γ^\pm, η^\pm . Hence without loss of generality we can choose $B_i = C_i = D_i = \xi(x)$ for all $i \in \mathbb{N}$. It follows that

$$\begin{aligned} & 2(\alpha(\gamma^{n_i} \eta^{n_i}) - \alpha(\gamma^{n_i}) - \alpha(\eta^{n_i})) \\ &= \beta((\gamma^{n_i} \eta^{n_i})^-, \gamma^-, \gamma^{n_i} x, x) + \beta(\gamma^+, (\gamma^{n_i} \eta^{n_i})^+, \gamma^{n_i} x, x) \\ & \quad + \beta((\eta^{n_i} \gamma^{n_i})^-, \eta^-, \eta^{n_i} x, x) + \beta(\eta^+, (\eta^{n_i} \gamma^{n_i})^+, \eta^{n_i} x, x) \\ &= \beta((\gamma^{n_i} \eta^{n_i})^-, \gamma^-, \gamma^{n_i} x, x) + \beta(\gamma^+, (\eta^{n_i} \gamma^{n_i})^+, x, \gamma^{-n_i} x) \\ & \quad + \beta((\eta^{n_i} \gamma^{n_i})^-, \eta^-, \eta^{n_i} x, x) + \beta(\eta^+, (\gamma^{n_i} \eta^{n_i})^+, x, \eta^{-n_i} x). \end{aligned}$$

Finally, we observe that

$$\begin{aligned} & \lim_{i \rightarrow \infty} \beta((\gamma^{n_i} \eta^{n_i})^-, \gamma^-, \gamma^{n_i} x, x) = \beta(\eta^-, \gamma^-, \gamma^+, x), \\ & \lim_{i \rightarrow \infty} \beta(\gamma^+, (\eta^{n_i} \gamma^{n_i})^+, x, \gamma^{-n_i} x) = \beta(\gamma^+, \eta^+, x, \gamma^-) = \beta(x, \gamma^-, \gamma^+, \eta^+), \\ & \lim_{i \rightarrow \infty} \beta((\eta^{n_i} \gamma^{n_i})^-, \eta^-, \eta^{n_i} x, x) = \beta(\gamma^-, \eta^-, \eta^+, x) = \beta(\eta^-, \gamma^-, x, \eta^+), \\ & \lim_{i \rightarrow \infty} \beta(\eta^+, (\gamma^{n_i} \eta^{n_i})^+, x, \eta^{-n_i} x) = \beta(\eta^+, \gamma^+, x, \eta^-) = \beta(\eta^-, x, \gamma^+, \eta^+), \end{aligned}$$

and conclude our result using Proposition A.2(4). \square

A.2. Crossratios and Eigenvalues

In this subsection we define, for the linear case, appropriate counterparts of the affine crossratios. Affine crossratios can be seen as infinitesimal versions of these crossratios.

Let $\{W_i\}_{i=1}^4$ be four $(n-1)$ -dimensional isotropic subspaces in $\mathbb{R}^{n,n}$ such that their orthogonal spaces are mutually transverse to each other. We recall that for $i \neq j$, $W_i^\perp \cap W_j^\perp$ contain exactly two isotropic lines. We use Lemmas 2.12 and 2.16 to choose $v_{i,j}^\pm$ arbitrarily from one of these two lines such that $\mathbb{R}v_{i,j}^+ \oplus W_i$ (resp. $\mathbb{R}v_{i,j}^- \oplus W_i$) lies in the orbit of $\text{cspan}([I_n, I_n]^t)$ (resp. $J\text{cspan}([I_n, I_n]^t)$) under the action of $\text{SO}_0(n, n)$.

We define the following *crossratio*:

$$\theta_{1,2,3,4} = \theta(W_1, W_2, W_3, W_4) := \frac{\langle v_{1,3}^+ | v_{2,3}^- \rangle \langle v_{2,4}^+ | v_{1,4}^- \rangle}{\langle v_{2,4}^+ | v_{2,3}^- \rangle \langle v_{1,3}^+ | v_{1,4}^- \rangle}.$$

As $v_{i,j}^\pm$ are unique upto scaling, we observe that the above expression does not depend on the choice of the vectors $v_{i,j}^\pm$ and hence is well defined.

LEMMA A.7. — Let W_*, W_i, W_j, W_k be four $(n-1)$ dimensional isotropic subspaces such that their orthogonal spaces are mutually transverse to each other. Then the following identity holds:

$$\frac{\langle v_{*,i}^+ | v_{*,i}^- \rangle \langle v_{*,j}^+ | v_{*,k}^- \rangle}{\langle v_{*,j}^+ | v_{*,i}^- \rangle \langle v_{*,i}^+ | v_{*,k}^- \rangle} = 1 = \frac{\langle v_{i,*}^+ | v_{i,*}^- \rangle \langle v_{j,*}^+ | v_{k,*}^- \rangle}{\langle v_{j,*}^+ | v_{i,*}^- \rangle \langle v_{i,*}^+ | v_{k,*}^- \rangle}.$$

Moreover, for n even, $\mathbb{R}v_{i,j}^\pm = \mathbb{R}v_{j,i}^\mp$ and for n odd, $\mathbb{R}v_{i,j}^\pm = \mathbb{R}v_{j,i}^\pm$.

Proof. — As $\mathbb{R}v_{*,y}^\pm \oplus W_* \subset W_*^\perp$ is a maximal isotropic subspace for $y = i, j, k$ and $\mathbb{R}v_{*,y}^+ \oplus W_*$ (resp. $\mathbb{R}v_{*,y}^- \oplus W_*$) lie in the orbit of $\text{cspan}([I_k, I_k]^t)$ (resp. $J\text{cspan}([I_k, I_k]^t)$) under the action of $\text{SO}_0(n, n)$, we deduce that there exist non-zero constants a^\pm, b^\pm such that

$$v_{*,i}^\pm - a^\pm v_{*,j}^\pm \in W_* \ni v_{*,i}^\pm - b^\pm v_{*,k}^\pm.$$

Hence, $\langle v_{*,i}^+ | v_{*,i}^- \rangle = a^+ \langle v_{*,j}^+ | v_{*,i}^- \rangle = b^- \langle v_{*,i}^+ | v_{*,k}^- \rangle = a^+ b^- \langle v_{*,j}^+ | v_{*,k}^- \rangle$ and the left hand side of the identity follows.

Moreover, we observe that $I_{n,n}W_\pm = W_\mp = -I_{n-1,n+1}W_\pm$, $I_{n,n}v_\pm = v_\mp$ and $-I_{n-1,n+1}v_\pm = v_\pm$. As $I_{n,n} \in \text{SO}_0(n, n)$ for n even, we obtain that $\nu_\pm([I_{n,n}]) = \nu_\mp([I])$. As $-I_{n-1,n+1} \in \text{SO}_0(n, n)$ for n odd, we obtain that $\nu_\pm([-I_{n-1,n+1}]) = \nu_\pm([I])$. Suppose $g \in \text{SO}_0(n, n-1)$ is such that $gW_+ = W_i$ and $gW_- = W_j$. Then $\mathbb{R}v_{i,j}^\pm = \mathbb{R}\nu_\pm([g])$. Finally, we conclude by observing that, for n even, $\mathbb{R}v_{i,j}^\pm = \mathbb{R}\nu_\pm([g]) = \mathbb{R}\nu_\mp([gI_{n,n}]) = \mathbb{R}v_{j,i}^\mp$ and for n odd, $\mathbb{R}v_{i,j}^\pm = \mathbb{R}\nu_\pm([g]) = \mathbb{R}\nu_\pm([-gI_{n-1,n+1}]) = \mathbb{R}v_{j,i}^\pm$. \square

PROPOSITION A.8. — Let $\{W_i\}_{i=1}^4$ and W_* be five $(n-1)$ dimensional isotropic subspaces such that their orthogonal spaces are transverse to each other and let $g \in \mathbf{G}$. Then the following identities hold:

- (1) $\theta(gW_1, gW_2, gW_3, gW_4) = \theta(W_1, W_2, W_3, W_4)$,
- (2) $\theta_{1,2,3,4} = \theta_{2,1,4,3} = \theta_{3,4,1,2}^{(-1)^n} = \theta_{4,3,2,1}^{(-1)^n}$,
- (3) $\theta_{1,2,3,4}\theta_{1,2,4,3} = 1$,
- (4) $\theta_{1,*,3,4}\theta_{*,2,3,4} = \theta_{1,2,3,4}$.

Moreover, for n even, $\theta_{1,2,3,4}\theta_{1,3,4,2}\theta_{1,4,2,3} = 1$.

Proof. — The first two identity follows from the definition of θ and Lemma A.7. Also, the third identity follows from Lemma A.7 by taking $j = k$.

We use the definition of θ and cancel the terms appearing both in the numerator and denominator to see that

$$\frac{\theta_{*,2,3,4}\theta_{1,*,3,4}}{\theta_{1,2,3,4}} = \frac{\langle v_{*,3}^+ | v_{2,3}^- \rangle \langle v_{2,4}^+ | v_{*,4}^- \rangle \langle v_{1,3}^+ | v_{*,3}^- \rangle \langle v_{*,4}^+ | v_{1,4}^- \rangle}{\langle v_{*,3}^+ | v_{*,4}^- \rangle \langle v_{*,4}^+ | v_{*,3}^- \rangle \langle v_{1,3}^+ | v_{2,3}^- \rangle \langle v_{2,4}^+ | v_{1,4}^- \rangle}.$$

The fourth identity follows by replacing the above formula by the following identities which are obtained by repeated application of Lemma A.7:

$$\begin{aligned}\langle v_{*,3}^+ | v_{*,4}^- \rangle \langle v_{*,4}^+ | v_{*,3}^- \rangle &= \langle v_{*,3}^+ | v_{*,3}^- \rangle \langle v_{*,4}^+ | v_{*,4}^- \rangle, \\ \langle v_{*,3}^+ | v_{*,3}^- \rangle \langle v_{1,3}^+ | v_{2,3}^- \rangle &= \langle v_{1,3}^+ | v_{*,3}^- \rangle \langle v_{*,3}^+ | v_{2,3}^- \rangle, \\ \langle v_{*,4}^+ | v_{*,4}^- \rangle \langle v_{2,4}^+ | v_{1,4}^- \rangle &= \langle v_{2,4}^+ | v_{*,4}^- \rangle \langle v_{*,4}^+ | v_{1,4}^- \rangle.\end{aligned}$$

Moreover, by repeated use of Lemma A.7 we obtain the following identities:

$$\begin{aligned}\langle v_{1,3}^+ | v_{2,3}^- \rangle \langle v_{4,3}^+ | v_{1,3}^- \rangle &= \langle v_{1,3}^+ | v_{1,3}^- \rangle \langle v_{4,3}^+ | v_{2,3}^- \rangle, \\ \langle v_{2,4}^+ | v_{1,4}^- \rangle \langle v_{1,4}^+ | v_{3,4}^- \rangle &= \langle v_{1,4}^+ | v_{1,4}^- \rangle \langle v_{2,4}^+ | v_{3,4}^- \rangle, \\ \langle v_{3,2}^+ | v_{1,2}^- \rangle \langle v_{1,2}^+ | v_{4,2}^- \rangle &= \langle v_{1,2}^+ | v_{1,2}^- \rangle \langle v_{3,2}^+ | v_{4,2}^- \rangle.\end{aligned}$$

Plugging it in we obtain that

$$\theta_{1,2,3,4} \theta_{1,3,4,2} \theta_{1,4,2,3} = \frac{\langle v_{4,3}^+ | v_{2,3}^- \rangle \langle v_{2,4}^+ | v_{3,4}^- \rangle \langle v_{3,2}^+ | v_{4,2}^- \rangle}{\langle v_{2,4}^+ | v_{2,3}^- \rangle \langle v_{3,2}^+ | v_{3,4}^- \rangle \langle v_{4,3}^+ | v_{4,2}^- \rangle}.$$

As $v_{i,j}^+ = v_{j,i}^-$ for n even, we conclude that $\theta_{1,2,3,4} \theta_{1,3,4,2} \theta_{1,4,2,3} = 1$. \square

Remark A.9. — Suppose $g \in G$ is such that its action on the space of $(n-1)$ dimensional isotropic subspaces has an attracting fixed point W_a and a repelling fixed point W_r and suppose W_a^\perp and W_r^\perp are transverse to each other. We call such elements *proto-pseudo-hyperbolic* and recall that

$$g^{\pm 1} v_{a,r}^+ = \lambda(g)^{\pm 1} v_{a,r}^+.$$

PROPOSITION A.10. — Suppose $g \in G$ is a proto-pseudo-hyperbolic element with attracting fixed point W_a and a repelling fixed point W_r . Then for any $(n-1)$ dimensional isotropic subspace W_* whose orthogonal space is transverse to both W_a^\perp and W_r^\perp the following holds:

- (1) $\theta(W_r, W_a, gW_*, W_*) = \lambda(g)^2$ when n is even,
- (2) $\theta(W_r, W_a, gW_*, W_*) = 1$ when n is odd.

Proof. — We denote gA_* by A_{g*} and use Lemma A.7 to deduce that

$$\begin{aligned}\theta(A_r, A_a, gA_*, A_*) &= \frac{\langle v_{r,g*}^+ | v_{a,g*}^- \rangle \langle v_{a,*}^+ | v_{r,*}^- \rangle}{\langle v_{a,*}^+ | v_{a,g*}^- \rangle \langle v_{r,g*}^+ | v_{r,*}^- \rangle} = \frac{\langle gv_{r,*}^+ | gv_{a,*}^- \rangle \langle v_{a,*}^+ | v_{r,*}^- \rangle}{\langle v_{a,*}^+ | gv_{a,*}^- \rangle \langle gv_{r,*}^+ | v_{r,*}^- \rangle} \\ &= \frac{\langle v_{r,*}^+ | v_{a,*}^- \rangle \langle v_{a,*}^+ | v_{r,*}^- \rangle}{\langle v_{a,*}^+ | gv_{a,*}^- \rangle \langle gv_{r,*}^+ | v_{r,*}^- \rangle} = \frac{\langle v_{a,*}^+ | v_{a,*}^- \rangle \langle v_{r,*}^+ | v_{r,*}^- \rangle}{\langle v_{a,*}^+ | gv_{a,*}^- \rangle \langle gv_{r,*}^+ | v_{r,*}^- \rangle}.\end{aligned}$$

Again using Lemma A.7 twice more we obtain the following two identities:

$$\begin{aligned}\langle v_{r,g*}^+ | v_{r,*}^- \rangle \langle v_{r,a}^+ | v_{r,a}^- \rangle &= \langle v_{r,g*}^+ | v_{r,a}^- \rangle \langle v_{r,*}^- | v_{r,a}^+ \rangle, \\ \langle v_{a,g*}^- | v_{a,*}^+ \rangle \langle v_{a,r}^- | v_{a,r}^+ \rangle &= \langle v_{a,g*}^- | v_{a,r}^+ \rangle \langle v_{a,*}^+ | v_{a,r}^- \rangle.\end{aligned}$$

Therefore, we deduce that:

$$\begin{aligned}\langle gv_{r,*}^+ | v_{r,*}^- \rangle \langle v_{r,a}^+ | v_{r,a}^- \rangle &= \langle gv_{r,*}^+ | v_{r,a}^- \rangle \langle v_{r,*}^- | v_{r,a}^+ \rangle \\ &= \langle v_{r,*}^+ | g^{-1}v_{r,a}^- \rangle \langle v_{r,*}^- | v_{r,a}^+ \rangle \\ &= \lambda(g)^{(-1)^{n-1}} \langle v_{r,*}^+ | v_{r,a}^- \rangle \langle v_{r,*}^- | v_{r,a}^+ \rangle \\ &= \lambda(g)^{(-1)^{n-1}} \langle v_{r,*}^+ | v_{r,*}^- \rangle \langle v_{r,a}^- | v_{r,a}^+ \rangle, \\ \langle gv_{a,*}^- | v_{a,*}^+ \rangle \langle v_{a,r}^- | v_{a,r}^+ \rangle &= \langle gv_{a,*}^- | v_{a,r}^+ \rangle \langle v_{a,*}^+ | v_{a,r}^- \rangle \\ &= \langle v_{a,*}^- | g^{-1}v_{a,r}^+ \rangle \langle v_{a,*}^+ | v_{a,r}^- \rangle \\ &= \lambda(g)^{-1} \langle v_{a,*}^- | v_{a,r}^+ \rangle \langle v_{a,*}^+ | v_{a,r}^- \rangle \\ &= \lambda(g)^{-1} \langle v_{a,*}^- | v_{a,*}^+ \rangle \langle v_{a,r}^+ | v_{a,r}^- \rangle.\end{aligned}$$

Hence, $\theta(A_r, A_a, gA_*, A_*) = \lambda(g)\lambda(g)^{(-1)^n}$ and our result follows. \square

Remark A.11. — Suppose $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$. Hence, for all infinite order elements $\gamma \in \Gamma$ we obtain that the action of $\rho(\gamma)$ on the space of $(n-1)$ -dimensional isotropic subspaces of $\mathbb{R}^{n,n}$ has an attracting fixed point and a repelling fixed point. We abuse notation and denote the attracting fixed point by $\xi_\rho(\gamma^+)$ and the repelling fixed point by $\xi_\rho(\gamma^-)$. Henceforth, we fix ρ and omit the subscripts ρ from the eigenvalues and crossratios. Also, when there is no confusion of notation, for $a, b, c, d \in \partial_\infty \Gamma$ all distinct, we denote $\theta(\xi(a), \xi(b), \xi(c), \xi(d))$ by $\theta(a, b, c, d)$.

PROPOSITION A.12. — Suppose n is even, $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ and $\gamma, \eta \in \Gamma$ are two infinite order elements such that the four points $\gamma^\pm, \eta^\pm \in \partial_\infty \Gamma$ are distinct and the sequence $\{\gamma^m \eta^k\}_{m \in \mathbb{N}} \subset \Gamma$ contains a subsequence $\{\gamma^{n_i} \eta^k\}_{i \in \mathbb{N}}$ consisting only of infinite order elements. Then the following identity holds:

$$\lim_{i \rightarrow \infty} \frac{\lambda(\gamma^{n_i} \eta^k)^2}{\lambda(\gamma^{n_i})^2 \lambda(\eta^k)^2} = \theta(\eta^-, \gamma^-, \gamma^+, \eta^k \gamma^+) \theta(\eta^+, \gamma^+, \gamma^-, \eta^{-k} \gamma^-).$$

Proof. — The proof follows exactly word to word as in the proof of Proposition A.5 by replacing the appearances of α by $\log \lambda$, β by $\log \theta$ and replacing the appearances of Proposition A.2 and Proposition A.3 respectively by Proposition A.8 and Proposition A.10. \square

PROPOSITION A.13. — Suppose n is even, $\rho \in \text{Hom}(\Gamma, \mathbf{G}, \mathbf{P}^\pm)$ and $\gamma, \eta \in \Gamma$ are two infinite order elements such that the four points $\gamma^\pm, \eta^\pm \in \partial_\infty \Gamma$ are distinct and the sequence $\{\gamma^m \eta^m\}_{m \in \mathbb{N}} \subset \Gamma$ contains a subsequence $\{\gamma^{n_i} \eta^{n_i}\}_{i \in \mathbb{N}}$ consisting only of infinite order elements. Then the following identity holds:

$$\lim_{i \rightarrow \infty} \frac{\lambda(\gamma^{n_i} \eta^{n_i})^2}{\lambda(\gamma^{n_i})^2 \lambda(\eta^{n_i})^2} = \theta(\eta^-, \gamma^-, \gamma^+, \eta^+)^2.$$

Proof. — The proof follows exactly word to word as in the proof of Proposition A.6 by replacing the appearances of α by $\log \lambda$, β by $\log \theta$ and replacing the appearances of Proposition A.2 and Proposition A.3 respectively by Proposition A.8 and Proposition A.10. \square

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Sourav GHOSH
Ashoka University, Rajiv Gandhi Education City,
Sonapat, Haryana 131029, India
sourav.ghosh@ashoka.edu.in
sourav.ghosh.bagui@gmail.com