1. Introduction

Classical Teichmüller theory provides links between complex analytic and dynamical quantities defined on Riemann surfaces with conformal hyperbolic metrics. More precisely, properties of the geodesic flow of a hyperbolic structure are related to holomorphic objects on the underlying Riemann surface. The Selberg trace formula is an instance of this correspondence. The goal of this paper is to extend this relationship in the context of higher rank Teichmüller theory. Specifically, in the case of Hitchin representations we find analogs to the fundamental results of Wolpert – as well as those of Hejhal and Gardiner – that compute variations of dynamical quantities for deformations of the complex structure parametrized by holomorphic differentials. In particular, we refer here to Gardiner’s formula [15] which computes the variation of the length of a geodesic in terms of Hejhal’s periods of quadratic differentials; the relation between the Thurston and Weil–Petersson metrics [42]; the computation of the variation of the cross ratio on the boundary at infinity of surface groups and the study of Fenchel–Nielsen twists [41].

Let us be more concrete. Let $X$ be a closed Riemann surface of genus at least two, and $\Sigma$ the underlying oriented differentiable manifold. Let $\delta_X$ be the monodromy of the unique conformal hyperbolic metric on $X$. Let $\iota_n$ be the irreducible representation of $\text{PSL}(2, \mathbb{R})$ in $\text{PSL}(n, \mathbb{R})$. The Fuchsian point is the representation $\delta_{X,n} = \iota_n \circ \delta_X : \pi_1(\Sigma) \to \text{PSL}(n, \mathbb{R})$.

A Hitchin representation is a homomorphism $\delta : \pi_1(\Sigma) \to \text{PSL}(n, \mathbb{R})$ that can be continuously deformed to the Fuchsian point. We call the set $\mathcal{H}(\Sigma, n)$ of conjugacy classes of Hitchin representations the Hitchin component. The Fuchsian locus is the subset of $\mathcal{H}(\Sigma, n)$ consisting of Fuchsian points obtained by varying the complex structure on $X$. By an abuse of terminology, we shall refer to these Fuchsian points as Fuchsian representations $\pi_1(\Sigma) \to \text{PSL}(n, \mathbb{R})$. Furthermore, throughout this paper we can and will assume a lift of Hitchin representations from $\text{PSL}(n, \mathbb{R})$ to $\text{SL}(n, \mathbb{R})$.

Hitchin [21] proves that $\mathcal{H}(\Sigma, n)$ can be globally parametrized by the Hitchin base: $Q(X, n) = \bigoplus_{k=2}^{n} H^0(X, K^k)$, where $K$ is the canonical bundle of $X$. Thus, the tangent space of the Fuchsian point of the Hitchin component can also be described as $Q(X, n)$. This infinitesimal parametrization, which will be crucial for our calculations, depends on some choices, and it is natural to normalize so that the restriction to the Fuchsian locus corresponds to classical deformations in Teichmüller space. In

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fact, in this paper we shall use two natural families of deformations (that is vectors in $T_{\delta} \mathcal{H}(\Sigma, n)$) associated to a point $q \in \mathcal{Q}(X, n)$ which are related by a constant depending only on $n$ and $k$:

(1) The standard deformation $\psi^0(q)$, for which the result of the computations are easier to state,

(2) The normalized deformation $\psi(q)$, for which the Atiyah–Bott–Goldman symplectic structure of the Hitchin components coincides along the Fuchsian locus with the symplectic structure inherited from the $L^2$-metric on the Hitchin base (see Corollary 5.1.2).

Recall that the moduli space of representations $\pi_1(\Sigma) \to \text{SL}(n, \mathbb{C})$ is a hyperkähler variety [20]. This structure is reflected in three algebraically distinct descriptions: the Dolbeault (Higgs bundle) moduli space, the de Rham moduli space of flat connections, and the Betti moduli space of representations. We exhibit isomorphisms of the tangent space to the Fuchsian point in each of these manifestations as $\mathcal{Q}(X, n) \oplus \overline{\mathcal{Q}(X, n)}$. We furthermore show that the different points of view actually give rise to the same parametrization of the tangent space at the Fuchsian point. A key point is that the first variation of the harmonic metric for certain variations of Higgs bundles vanishes (see Theorem 3.5.1). This result may be viewed as a generalization of Ahlfors’ lemma on variations of the hyperbolic metric under quasiconformal deformations [1]. All this occupies Section 3.

The discussion above is the complex analytic side of the Hitchin component, and we now wish to relate it to the dynamical side. In [25], the first author shows that if $\delta$ is a Hitchin representation and $\gamma$ a nontrivial element in $\pi_1(\Sigma)$, then $\delta(\gamma)$ has $n$-distinct positive eigenvalues. The underlying idea is to associate to a Hitchin representation a geodesic flow (see also [18] and [8]), thus giving a dynamical characterization of the Hitchin component.

This leads to the main motivation for this paper. In [8], Bridgeman, Canary, Sambarino and the first author constructed a pressure metric on the Hitchin component whose restriction to the Fuchsian locus is the Weil–Petersson metric. In Section 6, we shall prove the following theorem.

**Theorem 1.0.1.** Let $\delta$ be a Fuchsian representation into $\text{SL}(n, \mathbb{R})$ associated to a Riemann surface $X$ with a conformal hyperbolic metric. Let $q$ be a holomorphic $k$-differential on $X$, $2 \leq k \leq n$, and let $\psi^0(q)$ be the associated standard deformation. Then the pressure metric is proportional to the $L^2$-metric:

$$P_\delta (\psi^0(q), \psi^0(q)) = \frac{1}{2^{k-1} \pi |X|} \left[ \frac{(k-1)(n-1)!}{(n-k)!} \right] \int_X \|q\|^2 d\sigma.$$  

Moreover, two deformations associated to holomorphic differentials of different degrees are orthogonal with respect to the pressure metric.

The first ingredient in the proof of this theorem is an extension, Theorem 4.0.2, of Gardiner’s formula to Hitchin representations. This computes the first variation of the eigenvalues of $\delta(\gamma)$ as a function of $\delta$ under a standard deformation. The result, proven in Section 4, is a generalization of the classical formula for holomorphic quadratic differentials [15]. We reproduce the statement here for the highest eigenvalue.
Theorem 1.0.2. [Gardiner formula] For Hitchin representations, the first variation at the Fuchsian locus of the largest eigenvalue $\lambda_{\gamma}$ of the holonomy along a simple closed geodesic $\gamma$ of hyperbolic length $\ell_{\gamma}$ along a standard Hitchin deformation given by $q \in H^0(X,K^\gamma)$, is

$$d \log \lambda_{\gamma}(\psi^\gamma(q)) = \frac{(-1)^k(n-1)!}{2^{k-2}(n-k)!} \int_0^{\ell_{\gamma}} \Re \left( q(\gamma',...,\gamma') \right) \, dt.$$

The complete result, Theorem 4.0.2, also gives the variation of the other eigenvalues, and Corollary 4.0.4 gives the variation of the trace. The proof of Gardiner’s original formula makes use of the theory of quasiconformal maps. For Hitchin representations, no such technique is available, and our proof is purely gauge theoretic. Finally by a formula in Hejhal [19] (attributed to Petersson) the right hand side of the equation above can be interpreted as the $L^2$ pairing of $q$ with the relative Poincaré series $\Theta^{(k)}_\gamma$ associated to $\gamma$ (cf. Section 2.1.3 and Proposition 2.1.1).

The second component in the proof of Theorem 1.0.1 is a relationship, proved in Section 6, between the variance and the $L^2$-metrics for holomorphic differentials. The correspondence with the variance metric in the case of quadratic differentials has been discussed by different methods in McMullen [28]. We conclude with a remark on the family of pressure metrics that one can define on the Hitchin component.

Let us pause to note an interesting consequence of the dependence on the degree of the differential in Theorem 1.0.1. We shall see later on that the Atiyah-Bott-Goldman symplectic form $\omega_n$ on $H(\Sigma,n)$ is also related to the $L^2$-pairing. Given the pressure metric $P$ and the symplectic form $\omega$, one obtains the pressure endomorphism $A$ so that $P(u,v) = \omega_n(A \cdot u, v)$. Observe that $A$ is analytic on the Hitchin component. Our result is the following

Corollary 1.0.3. The eigenspaces of $A^2$ along the Fuchsian locus are, after the identification of the tangent space with the Hitchin base, precisely the subbundles $H^0(X,K^\gamma)$ consisting of holomorphic differentials of degree $k$. Moreover, the induced complex structure (by $A$) on these eigenspaces coincides with the complex structure on $H^0(X,K^\gamma)$.

In other words, two objects from the dynamical side, the pressure metric and symplectic form, detect the decomposition along the Fuchsian locus of the tangent space in sum of complex bundles of holomorphic differentials, a decomposition coming from the analytic side. It therefore seems interesting to study the decomposition into eigenspaces of $A^2$ everywhere on the Hitchin component, and not only along the Fuchsian locus.

We further investigate the symplectic geometry of the Hitchin component in Section 5.1. To a simple closed geodesic $\gamma$ and an element $h$ of the Cartan subalgebra of $\text{PSL}(n,\mathbb{R})$ we associate a higher Fenchel–Nielsen twist $\tau_{\gamma}(h)$, which is a vector field on $H(\Sigma,n)$. In Corollary 5.2.2, we show that these twist deformations are represented by linear combinations of the relative Poincaré series of $\gamma$ (of different degrees). The Hamiltonian vector fields of the eigenvalues of $\delta(\gamma)$, viewed as functions of the representation $\delta$, are expressed as linear combinations of the twist deformations about $\gamma$. We also prove the following generalization of [41, Theorem 2.4].

Theorem 1.0.4. [Reciprocity of the twist deformation] For all integers $k, p$, $1 \leq p \leq n$, $2 \leq k \leq n$, and any simple closed geodesics $\alpha, \beta$, the following holds at the Fuchsian locus:

$$d \log \lambda^{(p)}_{\alpha}(\psi^\beta(\Theta^{(k)}_\alpha)) = -d \log \lambda^{(p)}_{\beta}(\psi^\alpha(\Theta^{(k)}_\beta)),$$
where $\lambda^{(p)}(\delta)$ is the $p$th-largest eigenvalue of $\delta(\gamma)$.

In Section 7, we compute the variation for the cross ratio under Hitchin deformations, generalizing [41, Lemma 1.1.]. We give two formulations of this result: Theorem 7.2.1 using a generalized period that we call a rhombus function, and Theorem 7.0.1 using automorphic forms. We also comment on the triple ratios for $\text{SL}(3, \mathbb{R})$.

Finally, in Section 8, we provide two applications to large $n$-asymptotics. First, in Theorem 8.1.1 we show that (after a further normalization) the pressure metric converges for large rank to a multiple of the $L^2$-metric. In this situation, it is more natural to consider the renormalized highest eigenvalue $\mu_\gamma = \frac{\lambda_1}{n-1}$, and the associated renormalized pressure metric. The reason for this choice is so that the highest eigenvalue does not depend on $n$ along the Fuchsian locus.

We prove the following

**Theorem 1.0.5. [Large $n$-asymptotics]** The large $n$ asymptotics for the renormalized pressure metric and renormalized deformation $\psi(q)$ associated to a holomorphic $k$-differential $q$ is given by

$$P(\psi(q), \psi(q)) \sim \frac{(2k-1)!}{2^{k-1} \cdot 3^k \cdot |\chi(X)|} \int_X \|q\|^2 d\sigma .$$

This theorem provides a link between two large $n$-asymptotic theories of the Hitchin component. In [22], Hitchin argues that one can build a Higgs bundle theory for $\text{SU}(\infty)$, regarded as the group of symplectic diffeomorphisms of the sphere. The Hitchin base is $\bigoplus_{n=2}^\infty H^0(K^n)$, and thus describes a large $n$-analytic side. On the other hand, in [26] the first author has shown that there is a Hitchin component $\mathcal{H}(\infty, \mathbb{R})$ of representations of $\pi_1(\Sigma)$ in $\text{SL}(\infty, \mathbb{R})$, where the latter is considered to be (a subgroup) of the group of Hamiltonian diffeomorphisms of the annulus. In this approach, all Hitchin components embed in $\mathcal{H}(\infty, \mathbb{R})$, and all representations in $\mathcal{H}(\infty, \mathbb{R})$ are associated to geodesic flows and spectra, thus providing a dynamical side to the story. It is therefore tempting to try to understand in which way, at least formally, these versions of $\text{SL}(\infty)$ and $\text{SU}(\infty)$ have the same complexification.

The second application is motivated by the question raised at the end of the introduction of [26], a somewhat simpler version of which was posed to us by Maryam Mirzakhani: are the geodesic currents arising from Hitchin representations dense in the space of all geodesic currents? An even simpler test question is the following. Given $p$ conjugacy classes of pairwise distinct primitive elements $\gamma_1, \ldots, \gamma_p$ in $\pi_1(\Sigma)$, let

$$\Lambda_\gamma^{(p)} := \bigcup_{n=2}^\infty \left\{ (\log \lambda_{\gamma_1}(\delta), \ldots, \log \lambda_{\gamma_p}(\delta)) | \delta \in \mathcal{H}(\Sigma, n) \right\} .$$

Then what is the closure of $\Lambda_\gamma^{(p)}$ in $(\mathbb{R}^+)^p$? We prove in Theorem 8.2.1 that this set has nonempty interior. This result should be compared with a result of Lawton, Louder and McReynolds [27] which states that two elements of $\pi_1(\Sigma)$ have different traces for a certain linear representation.

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2. Preliminaries

In this section we present background material. First, in Section 2.1 we give a review of holomorphic differentials, introducing periods, relative Poincaré series and the Petersson and $L^2$ pairings. Section 2.2 is a review of what we shall need from Lie theory. Next, in Section 2.3 we present a summary of Higgs bundles for general complex semisimple Lie groups. Finally, we give the definition of opers (here only for the case of $\text{SL}(n, \mathbb{C})$) in Section 2.4.

2.1. Holomorphic differentials.

2.1.1. The $L^2$-metric. Let $X$ be a closed Riemann surface of genus $\geq 2$. We will always assume a conformal metric on $X$ is the hyperbolic metric of constant curvature $-1$. We denote the area form by $d\sigma$, and the Hodge operator by $\ast$. The metric induces a hermitian structure $\langle \cdot, \cdot \rangle$ on $K$, and hence on $K^k$. For $q_1, q_2$ smooth $k$-differentials, we define the $L^2$-metric

$$\langle q_1, q_2 \rangle_X := \int_X \langle q_1, q_2 \rangle \, d\sigma.$$ 

In local holomorphic coordinates $z = x + iy$, and with slight abuse of notation the area form may be written $d\sigma = \sigma(z) dx \wedge dy$ for a locally defined function $\sigma(z)$. Let $h = (\langle dz, dz \rangle)^{-1}$. Then $2h = \sigma$, so that $d\sigma = i\hbar(z) dz \wedge d\bar{z}$. It follows that if $q_i = q_i(z) dz^k$,

$$\langle q_1, q_2 \rangle_X = \int_X q_1(z) \overline{q_2(z)} h^{-k}(z) \, d\sigma. \quad (1)$$

Warning: The definition above differs from the usual Petersson pairing [31]:

$$\langle q_1, q_2 \rangle_P := \int_X q_1(z) \overline{q_2(z)} \sigma^{-k}(z) \, d\sigma = 2^{-k} \langle q_1, q_2 \rangle_X.$$

Let $X = \Gamma \backslash \mathbb{H}^2$ be the uniformization of $X$ coming from the hyperbolic metric. We will assume that $\Gamma$ has been lifted (once and for all) to a discrete subgroup of $\text{SL}(2, \mathbb{R})$.

2.1.2. Integration along geodesics. Let $\pi : UX \to X$ be the unit tangent bundle of $X$, equipped with the Riemannian metric induced from $X$. Let $\phi_s$ be the geodesic flow and $\mu$ the Liouville measure normalized to be a probability measure. In general, if $f$ is a function on $UX$, and $\gamma : [0, \ell] \to UX$ a geodesic arc, we shall write

$$\int_{\gamma} f \, ds =: \int_0^\ell f(\phi_s(\gamma(0))) \, ds \quad (2)$$

For integers $k$, a smooth section $q$ of $K^k$ defines a complex valued function $\hat{q} : UX \to \mathbb{C}$ that is homogeneous of degree $k$ with respect to the $S^1$-action on $UX$. If $\gamma$ is a unit speed geodesic with parameter $s$ in $X$, we will also use the alternative notations

$$\int_0^\ell q(\gamma^*, \ldots, \gamma^*) \, ds := \int_\gamma \hat{q} \, ds.$$ 

In this notation we regard $q \in K^k$ as a $k$-$\mathbb{C}$-multilinear form on $TX$. Then $\overline{q} \in \overline{K}^k$ is defined by

$$\overline{q}(u_1, \ldots, u_k) = \overline{q(u_1, \ldots, u_k)}.$$
so that
\[
\int_0^t \overline{q}(\gamma', \ldots, \gamma') \, ds = \int_0^t q(\gamma', \ldots, \gamma') \, ds.
\] (3)

2.1.3. Relative Poincaré series. Let \( q \) be a holomorphic differential on \( \mathbb{H}^2 \). Let \( \Gamma \) be a Fuchsian group then the following series, when it exists and converges,
\[
\Theta(q) = \sum_{\eta \in \Gamma} \eta^*(q),
\]
is called the Poincaré series of \( q \) and is \( \Gamma \) invariant. For any pair of distinct points \( u, U \in \partial_{\infty} \mathbb{H}^2 \cup \mathbb{H}^2 \), let
\[
\theta_{u,U}(z) = \frac{(\bar{U} - \bar{u}) dz}{(z - \bar{u})(z - \bar{U})},
\] (4)
where \( \bar{u} \) and \( \bar{U} \) are the endpoints at infinity of the geodesic joining \( u \) and \( U \) with the obvious convention. Observe that for any \( \eta \in \text{PSL}(2, \mathbb{R}) \), we have
\[
\eta^* \theta_{u,U} = \theta_{\eta^{-1}(u), \eta^{-1}(U)}.
\]

Let \( \gamma \) be a closed geodesic associated to an element (also called \( \gamma \)) of \( \Gamma \), and \( u, U \) the repelling (resp. attracting) fixed points of \( \gamma \) in \( \partial \mathbb{H}^2 \). Observe that \( \gamma^* \Theta_{u,U} = \Theta_{u,U} \).
The relative Poincaré series (of order \( k \)) of \( \gamma \), is
\[
\Theta^{(k)}_{\gamma} = \sum_{\eta \in \Gamma/\langle \gamma \rangle} \eta^* \theta_{u,U}^{k}.
\]
From the definition, one immediately sees that
\[
\Theta^{(k)}_{\gamma^{-1}} = (-1)^k \Theta^{(k)}_{\gamma}.
\] (5)

By direct computation (see [19, eq. (77)]) we find

**Proposition 2.1.1.** For any automorphic form \( q \) of degree \( k \),
\[
\int_{\gamma} \overline{q} \, ds = r_k \cdot \langle q, \Theta^{(k)}_{\gamma} \rangle_{X},
\] (6)
where
\[
r_k = \frac{(-1)^k 2^{k-2} ((k - 1)!)^2}{(2k - 2)! \pi}.
\] (7)

**Remark 2.1.2.** Note that if \( \ell_{\gamma} \) denotes the length of \( \gamma \), then the theta series in [19, eq. (56)] is \( = (2 \sinh(\ell_{\gamma}/2))^{-k} \cdot \Theta^{(k)}_{\gamma} \).

2.2. Lie theory.

2.2.1. Principal \( \text{sl}(2) \) subalgebras. Here we review material from Kostant [23] using (partly) his convention about the generators. Let \( \mathbb{G} \) be a complex semisimple Lie group \( \mathbb{G} \) of rank \( l \) with Lie algebra \( \mathfrak{g} \). Let \( (\cdot, \cdot)_h \) denote the Killing form of \( \mathfrak{g} \), normalized so that the squared length of a longest root is 2. Fix generators for \( \text{sl}(2) \):
\[
[a, x] = x, \ [a, y] = -y, \ [x, y] = -a.
\] (8)
Here we use a different convention for the sign of \( x \) from that of Kostant. With this understood, the defining representation is denoted \( \kappa_2 : \text{sl}(2) \simeq \text{sl}(\mathbb{C}^2) \). Notice here
that, for the sake of coherence with [23], we use a (nonalgebraic) convention, which differs from that of Bourbaki.

\[
\kappa_2(a) = \begin{pmatrix} 1/2 & 0 \\ 0 & -1/2 \end{pmatrix}, \quad \kappa_2(x) = \begin{pmatrix} 0 & -1/\sqrt{2} \\ 0 & 0 \end{pmatrix}, \quad \kappa_2(y) = \begin{pmatrix} 0 & 1/\sqrt{2} \\ 1 & 0 \end{pmatrix}. \tag{9}
\]

Consider a principal \(\mathfrak{sl}(2, \mathbb{C})\) embedding \(\kappa_g : \mathfrak{sl}(2) \hookrightarrow \mathfrak{g}\). Then \(\mathfrak{g}\) decomposes into irreducible representations of the principal \(\mathfrak{sl}(2)\),

\[
\mathfrak{g} = \bigoplus_{i=1}^{m} \mathfrak{v}_i, \tag{10}
\]

with \(\dim_{\mathbb{C}} \mathfrak{v}_i = 2m_i + 1\). The numbers \(m_i\) are all distinct and ordered so that \(m_i < m_{i+1}\).

They are called the exponents of the group \(G\). The grading by the element \(a\) gives the decomposition

\[
\mathfrak{g} = \bigoplus_{m=-m_l}^{m_l} \mathfrak{g}_m. \tag{11}
\]

Given a Cartan involution \(\rho\) of \(\mathfrak{g}\), fixing the principal subalgebra, and such that \(y = \rho(x)\), we shall say that a basis \(\{e_1, \ldots, e_l\}\) (resp. \(\{f_1, \ldots, f_l\}\)) of highest (resp. lowest) weight vectors for the \(\mathfrak{sl}(2, \mathbb{C})\) action is normalized if

1. \(f_i = \rho(e_i)\), and
2. \((-e_k, f_k)_h = d(\mathfrak{g})\),

where \(d(\mathfrak{g})\) is the Dynkin index for the principal embedding, which is defined by

\[
d(\mathfrak{g}) := \frac{\langle \kappa_2(a), \kappa_2(a) \rangle_h}{\text{Tr}(\kappa_2(a)^2)}. \tag{13}
\]

Observe that a normalized highest weight vector is uniquely determined up to multiplication by an element of \(S^1\). The justification for the introduction of this normalization will appear in Section 5.1.

In the case \(G = \text{SL}(n, \mathbb{C})\), we can make a more explicit choice of highest weight vectors. Let \(\kappa_n : \mathfrak{sl}(2) \rightarrow \mathfrak{sl}(n)\) denote the principal embedding (unique up to conjugation), viewed as a linear embedding via the defining representation of \(\text{SL}(n, \mathbb{C})\). We will use two choices of highest and lowest weight vectors \(E^0_k, F^0_k \in \mathfrak{sl}(n, \mathbb{C})\), for the representation of \(\mathfrak{sl}(2)\). First, let \(X = \kappa_n(x), Y = \kappa_n(y) = \rho(X)\), where \(\rho(M) = -M^t\) is the standard Cartan involution. Then define the standard highest and lowest weight vectors

\[
E^0_k := \left(-\sqrt{2}X\right)^k, \quad F^0_k := \rho(E^0_k) = -\left(\sqrt{2}Y\right)^k. \tag{14}
\]

We now renormalize these highest and lowest weights with respect to the trace to obtain the normalized vectors:

\[
E_k = \eta_k E^0_k, \quad F_k = \rho(E_k), \tag{15}
\]

where \(\eta_k\) is a positive number so that \(E_k\) and \(F_k\) are normalized in the sense discussed above. That is,

\[
\eta_k^2 = -\frac{d(n)}{\text{Tr}(E^0_k F^0_k)}., \tag{16}
\]
where the Dynkin index $d(n)$ of the principal $\mathfrak{sl}(2) \hookrightarrow \mathfrak{sl}(n)$ is (cf. [30])

$$d(n) = \frac{\text{Tr}(\kappa_2(a)^2)}{\text{Tr}(\kappa_2(a)^2)} = \left(\frac{n+1}{3}\right).$$  \hspace{1cm} (17)

An explicit calculation of the normalization and the constant $\eta_k$ is given in Corollary A.2 in the Appendix:

$$\eta_k = \frac{1}{k!} \left(\frac{n+1}{3}\right)^{1/2} \cdot \left(\frac{n+k}{2k+1}\right)^{-1/2}.$$  \hspace{1cm} (18)

2.2.2. A useful basis. We will need the following technical result which is probably well-known.

**Lemma 2.2.1.** Given a choice of highest weight vectors $\{e_i\}_{i=1}^\ell$, there is a basis $\{h_j\}_{j=1}^\ell$ of the centralizer $\mathfrak{z}(X-Y)_R$, such that $(e_i, h_j)_0 = \delta_{ij}$.

We shall call $(h_j)_{j=1,\ldots,\ell}$ the principal basis of $\mathfrak{z}(X-Y)$.

**Proof.** Recall the decomposition (10). This is an orthogonal decomposition for the Killing form. First, observe that since $X-Y$ is conjugate to a multiple of $a$, $\dim \mathfrak{z}(X-Y) \cap v_i = \dim \mathfrak{z}(a) \cap v_i = 1$.

To conclude, we need to prove that the Killing pairing with $e_i$ defines a nonzero form on $\mathfrak{z}(X-Y) \cap v_i$. Recall that $v_i$ is generated by $e_i$ and its iterated images under $\text{ad}(X-Y)$. In particular, if $w \in \mathfrak{z}(X-Y)$ and $(w, e_i)_0 = 0$, then for all $k$, $(w, (\text{ad}(X-Y))^k e_i)_0 = 0$ and thus $w = 0$. We can therefore find an element $h_i$ in $v_i$ so that $(h_i, e_i)_0 = 1$. The result follows. \hfill $\Box$

2.2.3. The involution associated to the real split form. We briefly recall material from [23], which is also explained in [21, 3]. The choice of a principal subalgebra with its standard generators defines a Cartan subalgebra and a (C-antilinear) Cartan involution $\rho$. Moreover, one can define a C-linear involution $\sigma$ characterized by

$$\sigma(e_i) = -e_i, \quad \sigma(Y) = -Y.$$ 

Then $\sigma$ and $\rho$ commute. Furthermore, the real split involution $\lambda = \sigma \circ \rho$ is such that its set of fixed points is a real split subalgebra $\mathfrak{g}_0$ of $\mathfrak{g}$. The complexification of $\mathfrak{g}_0$ is $\mathfrak{g}$.

2.3. Higgs bundles.

2.3.1. The Hitchin component. A representation of $\pi_1(\Sigma)$ into the real split form $\mathbb{G}_R$ of $\mathbb{G}$ is Fuchsian if it is conjugated to a discrete faithful representation taking values in the principal $\text{SL}(2, \mathbb{R})$. A representation is Hitchin if it can be continuously deformed into a Fuchsian representation. The Hitchin component is then (a connected component) of the space whose points are equivalence classes of Hitchin representations up to conjugacy by an element of $\mathbb{G}_R$. We denote the Hitchin component by $\mathcal{H}(\Sigma, \mathbb{G}_R)$. 
2.3.2. Higgs bundles and the self duality equations.

**Definition 2.3.1.** A $G$-Higgs bundle is a pair $(P, \Phi)$, where $P$ is a holomorphic principal $G$-bundle $P \to X$, $G_P = P \times_G g$ is the associated holomorphic adjoint bundle, and $\Phi \in H^0(X, G_P \otimes K)$.

Let $K \subset G$ be a maximal compact subgroup. We will regard a reduction of $P$ to a principal $K$-bundle $P_K$ as arising from a smooth family of Cartan involutions on the fibers of $G_P$. By a slight abuse of notation, we denote such a family by $\rho$. A connection on $P_K$ induces a covariant derivative $\nabla$ on the vector bundle $G_P$ satisfying $\nabla(\rho) = 0$. Conversely, given $\rho$ and a connection on $P$ that is compatible with the holomorphic structure, there is a uniquely determined connection on $P_K$ called the Chern connection. The curvature $F_{\nabla}$ of such a connection is a section of $\Omega^2(X, G_{P_K})$.

With this understood, we introduce Hitchin’s equations.

\[
F_{\nabla} - [\Phi, \rho(\Phi)] = 0,
\]
\[
\nabla(\rho) = 0.
\]

For a given $G$-Higgs bundle $(P, \Phi)$, eq. (19) may be viewed as equations for $\rho$. A solution $\rho$ to eq. (19) will be called a harmonic metric \[13, 10, 24\]. As a consequence of the Hitchin equations, the associated connection $D = \nabla + \Phi - \rho(\Phi)$, is flat, and the associated representation of $\pi_1(\Sigma) \to G$ is called the monodromy of the solution of the self duality equations.

In this paper we will not need the details of the notion of (semi)stability of Higgs bundles, other than to note the following (cf. \[20, 33, 6\]).

**Theorem 2.3.2 (Hitchin-Kobayashi correspondence).** A stable $G$-Higgs bundle admits a harmonic metric, i.e. a solution to (19).

2.3.3. Hitchin sections. We recall here the construction of the Hitchin section, which depends on a normalization. First, the Hitchin base is defined as

\[
Q(X, g) := \bigoplus_{i=1}^l H^0(X, K^{m_i+1}), \tag{20}
\]

For $q = (q_1, \ldots, q_l) \in Q(X, g)$, we define the normalized Hitchin deformation

\[
\phi(q) = \sum_{i=1}^l q_i \otimes e_i \in \Omega^{1,0}(X, \mathcal{G}), \tag{21}
\]

where $\{e_i\}_{i=1}^l$ are normalized highest weight vectors for the principal $\text{sl}(2) \hookrightarrow g$ (see equation (12)). In the case of $\text{SL}(n, \mathbb{R})$ it is convenient to introduce the standard Hitchin deformation

\[
\phi^0(q) = \sum_{i=1}^l q_i \otimes E_i^0 \in \Omega^{1,0}(X, \mathcal{G}), \tag{22}
\]

We shall choose our Hitchin section $\mathcal{L} : Q(X, g) \to \mathcal{M}_{\text{Dol}}(X, G)$ to be

\[
\mathcal{L}(q_1, \ldots, q_l) := (\mathcal{G}, Y + \phi(q)). \tag{23}
\]

Then by [21] we have
Theorem 2.3.3 (Hitchin section). For any \( q \), the Higgs bundle \( L(q) \) admits a unique solution of the self-duality equations (19). If \( \delta(q) \) the monodromy of the corresponding solution modulo conjugacy, then the map \( q \to \delta(q) \) is a diffeomorphism from the Hitchin base \( Q(X, g) \) to the Hitchin component \( \mathcal{H}(\Sigma, G_R) \).

2.3.4. The Fuchsian Higgs bundle. Let us consider the holomorphic vector bundle

\[ \mathcal{H} = K^{-1} \oplus O \oplus K. \]

Choosing a spin bundle \( S \) on \( X \) identifies \( \mathcal{H} \) with the \( \mathfrak{sl}(2, \mathbb{C}) \) bundle of trace free endomorphisms of \( S \oplus S^{-1} \). Actually, this identification is independent on the choice of \( S \). The hyperbolic metric defines a metric on \( S \oplus S^{-1} \), and hence a Cartan involution on \( \mathcal{H} \). Finally, the canonical bracket map viewed as a holomorphic section of \( K \otimes K^{-1} \subset K \otimes \mathcal{H} \), defines a holomorphic section \( \Phi_0 \in \Omega^{1,0}(X, \mathcal{H}) \). The hyperbolic metric defines a connection on \( \mathcal{H} \), and all together \( \nabla, \rho, \Phi_0 \) satisfy (19).

2.3.5. The Fuchsian \( G \)-Higgs bundle. Let \( G \) be a complex Lie group equipped with a choice of a principal \( \text{SL}(2, \mathbb{C}) \) with its canonical generators \( a, X, Y \). We use the grading defined in eq. (11) to define a holomorphic bundle \( \mathcal{G} \) by

\[ \mathcal{G} := \bigoplus_{m=-m_1}^{m_1} \mathfrak{g}_m \otimes K^m. \]  

(24)

Since the complex vector bundle underlying \( \mathcal{G} \) will later be equipped with another holomorphic structure, we will refer to (24) as the split holomorphic structure. Observe that now \( \mathcal{H} \) maps into \( \mathcal{G} \) by

\[ (u, v, w) \mapsto (Y \otimes u, a \otimes v, W \otimes w_0). \]

Thus the Higgs field \( \Phi_0 \in \Omega^{1,0}(\mathcal{H}) \) defined in the previous paragraph gives rise to a Higgs field, also denoted by \( \Phi_0 \in \Omega^{1,0}(\mathcal{G}) \). By definition, the equivalence class of the Higgs bundle \( (\mathcal{G}, \Phi_0) \) is the Fuchsian point in \( \mathcal{M}_{\text{hol}}(X, \mathcal{G}) \).

Observe that the family of Cartan involutions on \( \mathcal{H} \) defined in the previous paragraph extends to a section of Cartan involutions on \( \mathcal{G} \) (also denoted \( \rho \)). Similarly the hyperbolic metric connection extends to a connection. This connection, also denoted \( V \) on \( \mathcal{G} \), is compatible with the holomorphic structure and metric: \( V(\rho) = 0 \).

In other words, \( V \) is the unique \( \rho \)-compatible Chern connection on \( \mathcal{G} \). Then, altogether the connection \( (V, \Phi_0, \rho) \) solves Hitchin equations (19).

In the special case of \( \text{SL}(n, \mathbb{C}) \), it is useful to consider the vector bundle

\[ \mathcal{E} = \text{Sym}^{n-1}(S \oplus S^*) = \bigoplus_{p=1}^{n} S^{2p-n-1}, \]  

(25)

where \( S \) is a spin bundle. Then \( \mathcal{G} \) will be the (trace free) endomorphism bundle of \( \mathcal{E} \).

2.4. Nonslipping connections and opers. In this section, we restrict ourselves to the case \( G = \text{GL}(n, \mathbb{C}) \) and refer to [4, 5] for general \( G \) (see also: [12], [37], [17] the original reference [14] and the geometric version [32] for further discussion). Let \( \mathcal{P} \to X \) be a holomorphic vector bundle. A holomorphic filtration of \( \mathcal{P} \) is a family \( \{ \mathcal{F}_p \}_{1 \leq p \leq n} \) of holomorphic subbundles of \( \mathcal{P} \) such that

- \( \mathcal{F}_n = \mathcal{P} \),
• $\mathcal{F}_{p-1} \subset \mathcal{F}_p$
• $\dim(\mathcal{F}_p) = p$.

**Definition 2.4.1.** [OPERS] A holomorphic connection $D$ on $\mathcal{P}$ equipped with a holomorphic filtration $\{F_p\}_{1 \leq p \leq n}$ is nonslipping if it satisfies the following conditions:

• $\forall F_p \subset F_{p+1}$ for all $p$,
• If $\alpha_p$ is the projection from $F_{p+1}$ to $F_{p+1}/F_p$, then the map 
  $$(X, u) \rightarrow \alpha_p(D_X(u)),$$
  considered as a linear map from $K^* \otimes F_p/F_{p-1} = S^{n+1-2p}$ to $F_{p+1}/F_p = S^{n+1-2p}$ is the identity.

A nonslipping connection is also called a $\text{GL}(n, \mathbb{C})$-oper.

2.4.1. The Veronese oper and the nonsplit holomorphic structure. Let $S$ be a spin bundle on the Riemann surface $X$, so that $S^2 = K$. Let $\mathcal{E}_{\text{op}} := I^{n-1}(S^{1-n})$ be the holomorphic rank $n$ bundle of $(n - 1)$-jets of holomorphic sections of $S^{1-n}$. Let $F_p$ be the vector subbundle of $\mathcal{G}_{\text{op}}$ defined by

$$F_p := \{j^{n-1}\sigma \mid j^{n-1}\sigma = 0\}.$$ 

The family $\{F_p\}_{1 \leq p \leq n}$ is a holomorphic filtration of $\mathcal{G}_{\text{op}}$: we have $F_n = \mathcal{E}_{\text{op}}, F_{p-1} \subset F_p$ and $\dim(F_p) = p$. Observe furthermore that

$$F_p/F_{p-1} = K^{n-p-1} \otimes S^{1-n} = S^{n-1-2p}.$$ 

In particular, the graded bundle associated to the filtration is given by $\mathcal{E}$ in (25). We let $\overline{\mathcal{E}}_{\text{op}}$ be the same underlying complex vector bundle, but with the holomorphic structure induced by $\overline{\mathcal{F}}_D$. We will call this the oper holomorphic structure. The Veronese oper or Fuchsian oper is $\mathcal{E}_{\text{op}}$ equipped with the above filtration and the Fuchsian holomorphic connection $D$.

3. Moduli spaces and the Tangent Space at the Fuchsian point

By construction, every Riemann surface defines a Hitchin parametrization of the Hitchin component by the Hitchin base. In particular, the tangent space at the Fuchsian point (the image of zero under this parametrization) is identified with the Hitchin base.

However, working directly with this description of the tangent space at the Fuchsian point might not be very handy since, at first sight, it involves solving an elliptic PDE. We rather describe another approach which will be more helpful in the sequel: roughly speaking we will describe the tangent space of Hitchin component at the Fuchsian point as the real part of an oper deformation.

Let us summarize here the construction in the following

**Proposition 3.0.2.** Let $q \in Q(X, g)$ be an element of the Hitchin base. Let $\lambda$ be the real split involution, $D$ the flat connection at the $\mathcal{G}$-Fuchsian point, and $\psi(q) = \phi(q) + \lambda(\phi(q)) \in \Omega^1(\Sigma, g)$. Then,

$$\text{d}_D\phi(q) = 0, \text{ d}_D\psi(q) = 0.$$

Furthermore, passing to cohomology, the map $\psi$ realizes an isomorphism $\psi : Q(X, g) \rightarrow H^1_D(g)$, which coincides with the isomorphism coming from the Hitchin parametrization.
In this proposition, we refer to the renormalized Hitchin section, but the same statement clearly holds for the standard Hitchin section.

By the previous proposition $\phi(q)$ can be considered as an element of $H^1_{dR}(\Sigma, G)$. It is the tangent vector to a one parameter family of flat connections: $D_t = D + t\phi(q)$, which we will call an oper deformation.

In the course of proving Proposition 3.0.2, we will actually describe and relate the Fuchsian points in various moduli spaces, parametrize their tangent spaces and spend some time describing intermediate results of independent interest.

3.1. Moduli spaces. We define the following moduli spaces (see [35, 36]):

1. The Dolbeault moduli space $M_{Dol}(X, G)$ of $s$-equivalence classes of semistable $G$-Higgs bundles on $X$,
2. The de Rham moduli space $M_{dR}(\Sigma, G)$ of gauge equivalence classes of reductive flat $G$-connections,
3. The Betti moduli space $M_B(\Sigma, G)$ of conjugacy classes of completely reducible representations $\pi_1(\Sigma) \to G$.

There are homeomorphisms, which are diffeomorphisms in the neighborhood of the Fuchsian point

$$M_{Dol}(X, G) \overset{HK}{\cong} M_{dR}(\Sigma, G) \overset{RH}{\cong} M_B(\Sigma, G)$$

where $HK$ (resp. $RH$) is the Hitchin–Kobayashi (resp. Riemann–Hilbert) correspondence.

We also introduce two distinguished submanifolds of the moduli space:

4. The Hitchin component $\mathcal{H}(\Sigma, G_R)$ of conjugacy classes of Hitchin representations into a split real form of $G$;
5. The oper moduli space $\mathcal{Op}(X, G)$ of gauge equivalence classes of $G$-opers.

This section focuses on the common (smooth) Fuchsian point in the moduli spaces that we have encountered before: the Fuchsian $G$-Higgs bundle, the Fuchsian or Veronese oper etc. The Fuchsian point is a point of transverse intersection of $\mathcal{H}(\Sigma, G_R)$ and $\mathcal{Op}(X, G)$. The outcome will be to describe this tangent space in its various guises using the Hitchin base and the Hitchin section.

More precisely we have several goals in this section.

- In Propositions 3.2.1 and 3.3.1 and 3.4.1, we describe, using Hitchin deformations, the tangent space at the Fuchsian point for the three moduli spaces $M_{Dol}(X, G)$, $M_{dR}(\Sigma, G)$, and $M_B(\Sigma, G)$. The description is in terms of isomorphisms with $Q(X, g) \oplus Q(X, \tilde{g})$, where $Q(X, g)$ is the Hitchin base (20). The corresponding tangent vectors are called Dolbeault, Hodge and Betti deformations, respectively.
- Then in Corollary 3.5.2, we show that all these descriptions coincide. In other words, the isomorphisms with the space of holomorphic differentials commute with the Riemann–Hilbert and Hitchin–Kobayashi correspondences. This follows from the vanishing of the first variation of the harmonic metric for Dolbeault deformations.
- Finally, we use this (now unambiguous) description to achieve our main goal: describing the tangent space of the Fuchsian point in the Hitchin component as the Hitchin base in Proposition 3.6.1. In particular, we relate this tangent space to the handy oper deformations.
We remark here that this discussion extends the various isomorphic descriptions of the cotangent space of Teichmüller space as holomorphic quadratic differentials [41], as well as the global parametrization via harmonic maps due to Wolf [40]. We also point out the thesis of Dalakov [11], which also studies the germ of the moduli space at the Fuchsian point.

3.2. Dolbeault deformations for Higgs bundles. Let us first consider the Dolbeault moduli space \( M_{\text{Dol}}(X, G) \). The first order deformations of the Fuchsian point can be described as follows.

The tangent space \( T_{(\Phi,q)}M_{\text{Dol}}(X, G) \) is given by the first cohomology of a deformation complex \( C_{\text{Dol}}(\mathcal{G}, \Phi) \). In the presence of the solution \((\nabla, \rho)\) of Hitchin self-duality equation, we may take harmonic representatives for \( H^1(C_{\text{Dol}}(\mathcal{G}, \Phi)) \) – and denoting the corresponding vector space by \( H^1(C_{\text{Dol}}(\mathcal{G}, \Phi)) \).

Proposition 3.2.1. [Dolbeault infinitesimal parametrization] At the Fuchsian point, the map

\[
(\phi, \beta) : Q(X, g) \oplus \overline{Q(X, g)} \rightarrow H^1(C_{\text{Dol}}(\mathcal{G}, \Phi))
\]

where, \( \rho \) being the Cartan involution,

\[
\beta(b) := \rho(\phi(b)),
\]

defines a complex linear isomorphism.

Proof. Notice that for \((\phi, \beta) = (\phi(q), \beta(b))\) as above, \( \phi \) is holomorphic and \( \beta \) is harmonic with respect to \( \partial_Y \). Moreover, \([\Phi, \beta(b)] = [Y, \beta(b)] = 0\), and \([\rho(\Phi), \phi(q)] = -[X, \phi(q)] = 0\). Hence, \((\phi(q), \beta(b)) \in H^1(C_{\text{Dol}}(\mathcal{G}, \Phi))\). This gives an inclusion \( Q(X, g) \oplus \overline{Q(X, g)} \hookrightarrow H^1(C_{\text{Dol}}(\mathcal{G}, \Phi)) \), and now the result follows for dimensional reasons.

We could actually use Serre duality, the hyperbolic metric on \( X \) and the Dolbeault isomorphism to identify \( \overline{Q(X, g)} \) with \( \bigoplus_{l=1}^{l} H^{0,1}(X, K^{-m}) \), where the script indicates harmonic forms.
3.3. **Hodge parametrization in the de Rham picture.** Recall that $M_{dR}(\Sigma, G)$ is the moduli space of reductive flat $G$-connections. By Corlette’s theorem [10, 13, 24], for any reductive flat connection $D$ and conformal structure $X$ on $\Sigma$, there exists a unique harmonic metric $\rho$. This completes the Hitchin-Kobayashi correspondence. Fixing a metric $\rho$, we can take harmonic representatives for the first cohomology of deformation complex at a flat connection $D$, and write:

$$\begin{align*}
T_{[D]}M_{dR}(\Sigma, G) &\cong H^1(C_{dR}(D)) = \{ B \in \Omega^1(X, \mathcal{G}) : DB = D^\ast B = 0 \},
\end{align*}$$

where $D^\ast$ is the formal adjoint of $D$ for the metric

$$A, B \mapsto \int_{\Sigma} (A, \rho(B \circ f))_0,$$

so that $D^\ast = \rho(D \circ f)$. Then at the Fuchsian point we have the following identification

**Proposition 3.3.1.** [Hodge infinitesimal parametrization] At the Fuchsian point $[D]$, the map

$$\begin{align*}
(\psi, \alpha) : Q(X, \phi) &\oplus \overline{Q(X, \rho)} \longrightarrow H^1(C_{dR}(D)).
\end{align*}$$

given by

$$\begin{align*}
\psi(q) &:= \phi(q) - \rho(\phi(q)), \\
\alpha(b) &:= \beta(b) + \rho(\beta(b)),
\end{align*}$$

defines a real linear isomorphism.

**Proof.** Recall: $D = \partial D + \delta D$, $D^\ast = i * (\delta D - \partial D)$. We first show that $B := \alpha(b) + \psi(q)$ is harmonic. Breaking into type, we see that the result is equivalent to $\delta D B = 0$ and $\partial D B = 0$. The first of these equations is

$$\delta D(\rho(\beta(b)) + \phi(q)) = \delta_{\mathcal{G}}(\rho(\beta(b)) + \phi(q)) - [\rho(\Phi), \rho(\beta(b)) + \phi(q)]$$

$$= \rho(\delta_{\mathcal{G}}(\rho(\beta(b)) + \phi(q)) - \rho(\Phi, \beta(b)) - [\rho(\Phi), \phi(q)])$$

$$= \rho(\delta_{\mathcal{G}}(\rho(\beta(b)) + \phi(q)) - \rho(\Phi, \beta(b)) + [X, \phi(q)])$$

$$= 0,$$

since the first two terms on the right hand side of eq. (30) vanish because the $q_k$ (resp. $b_k$) are holomorphic (resp. harmonic), and the last two terms vanish because the $e_k$ (resp. $f_k$) are highest (resp. lowest) weight vectors, and $\rho(Y) = X$. The second equation follows similarly. Note that we have used the fact that $\rho$ is parallel with respect to $\mathcal{G}$. The fact that the map is an isomorphism follows from dimensional reasons. \hfill \Box

3.4. **Oper parametrization in the Betti picture.** Let $\mathcal{V}$ be the local system determined by the holomorphic connection $D$ on $\mathcal{G}_{op}$. By Weil [38], the Zariski tangent space $T_{[D]}M_{dR}(\Sigma, G)$ is given by $H^1(\mathcal{V})$.

Now there is an exact sequence of sheaves of $\mathcal{C}$-modules, where $\mathcal{C}$ is the locally constant sheaf.

$$0 \longrightarrow \mathcal{V} \longrightarrow \mathcal{G}_{op} \longrightarrow \mathcal{G}_{op} \otimes K \longrightarrow 0.$$

Since $D$ is irreducible this gives an exact sequence in cohomology:

$$0 \longrightarrow H^0(\mathcal{G}_{op}) \longrightarrow H^0(\mathcal{G}_{op} \otimes K) \longrightarrow H^1(\mathcal{V}) \longrightarrow H^1(\mathcal{G}_{op}) \longrightarrow H^1(\mathcal{G}_{op} \otimes K) \longrightarrow 0,$$

The full tangent space to $M_{dR}(\Sigma, G)$ at the Veronese oper will be described by the next lemma.
Lemma 3.4.1. We have the following

(1) The inclusion \(Q(X, g) \subset H^0(G_{op} \otimes K)\) induces an isomorphism with the cokernel of the map \(H^0(G_{op}) \to H^0(G_{op} \otimes K)\).

(2) The inclusion \(Q(X, g) \subset H^1(G_{op})\) induces an isomorphism with the kernel of \(H^1(G_{op}) \to H^1(G_{op} \otimes K)\).

Proof. As in the proof of Proposition 3.3.1, \(Q(X, g) \subset \ker D^*,\) and hence \(Q(X, g)\) is orthogonal to the image of \(D^*\). Now by the Riemann–Roch formula,

\[
\dim \mathcal{C}H^0(G_{op} \otimes K) - \dim \mathcal{C}H^0(G_{op}) = \dim \mathcal{C}Q(X, g).
\]

This implies the first statement, and the second statement is proven similarly. □

The involution \(\lambda\) gives a splitting of the short exact sequence (31), which in turn via Lemma 3.4.1 provides a real linear isomorphism

\[
Q(X, g) \oplus Q(X, g) \to H^1(V),
\]

\[
(q, b) \mapsto (\phi(q), \lambda(\phi(b))).
\]

(32) (33)

Actually, this parametrization will prove to be different from the previous ones. The compatible choice is given in the following

Proposition 3.4.2. [Oper infinitesimal parametrization] At the Veronese oper, the mapping

\[
(\chi, \varepsilon) : Q(X; g) \oplus Q(X; g) \to T_D \mathcal{M}_{ur}(\Sigma, G)
\]

given by

\[
\chi(q) = \phi(q) + \lambda(\phi(q)), \quad \varepsilon(b) = \phi(b) - \lambda(\phi(b)),
\]

defines a real linear isomorphism of \(Q(X; g) \oplus Q(X; g)\) with \(H^1(V)\).

3.5. Identification of the different infinitesimal parametrizations. Our main result here is Corollary 3.5.2, which states that the three descriptions of the tangent space to the moduli space at the Fuchsian point given in the previous section are compatible with the Hitchin-Kobayashi and Riemann-Hilbert correspondences (26). This relies on the following theorem, which may be regarded as a generalization of the classic result of Ahlfors [1, Lemma 2].

3.5.1. Variation of the harmonic metric. In this section we prove

Theorem 3.5.1. The first variation of the harmonic metric vanishes for the Dolbeault deformations in Proposition 3.2.1.

Proof. Fix \((q, b) \in Q(X, g) \oplus Q(X, g)\). Let \(\dot{\rho}\) denote the first variation of the Cartan involution for the Dolbeault deformation \((\phi(q), \beta(b))\). Then \(\dot{\rho}\) is a family of derivations, and since \(G\) is semisimple there is a smooth section \(Z\) of \(G\) such that \(\dot{\rho} = \text{ad} Z\). Then for any other one parameter family of sections \(U\) of \(G\),

\[
\dot{\rho}(U) = \rho(U) + [\rho(Z), \rho(U)].
\]

(34)

For convenience, set \(k = -\rho(Z)\). To begin, we claim that the first variation of the connection satisfies:

\[
(V)^{1,0} = \partial V + \rho(\beta(b)).
\]

(35)
Indeed, since $\nabla$ is the Chern connection for $\rho$, for any fixed $U \in \Gamma(\mathcal{G})$ independent of the variational parameter, $\nabla(\rho(U)) = \rho(\nabla U)$. Hence differentiating, using eq. (34), we have

$$\nabla_p \rho(U) + \nabla[\rho(Z), \rho(U)] = [\rho(Z), \rho(\nabla U)] + \rho(\nabla_p \rho(U)).$$

Thus

$$\nabla p \rho(U) - [\nabla(k), \rho(U)] - [k, \nabla(\rho(U))] = -[\nabla-k, \rho(U)] + \rho(\nabla_p \rho(U)),$$

$$\nabla - \nabla(k), \rho(U)] = [\rho(\nabla), \rho(U)],$$

$$\nabla = \nabla(k) + \rho(\nabla_p \rho(U)).$$

(36)

since $U$ was arbitrary. Now by definition, $\nabla(\rho(0)) = \rho(0)$, so (35) follows by taking the $(1,0)$-part of (36). Notice that from (35),

$$\nabla_p = \nabla = \nabla(\rho(0)) = \rho(0) + \partial \rho(\nabla) = \rho(\nabla).$$

Differentiating (19), we then get

$$\nabla_p \partial \rho k - [\phi(q), \rho(\Phi)] - [\Phi, \rho(\phi(q))] + [\Phi, \rho(\phi(q))] = 0$$

$$\nabla_p \partial \rho k - [\phi(q), X] + \rho(\phi(q), X) + [\Phi, \rho(\phi(q))] = 0.$$

However, as in the proof of Proposition 3.3.1, since $\phi(q)$ is a combination of highest weight vectors,

$$[\phi(q), X] = \rho(\phi(q), X) = 0.$$

We therefore obtain

$$\nabla_p \partial \rho k + [\Phi, [k, \rho(\Phi)]] = 0.$$

(37)

Let $\langle U, V \rangle = -\langle U, \rho(V) \rangle$. Then using the Kähler identities

$$\partial \rho = i * \partial \rho, \quad \partial \rho = -i * \partial \rho.$$

(38)

Using (37) and integration by parts yields successively

$$\int_X (\nabla_p \partial \rho k, \rho(k))_g \sigma d\sigma + \int_X ([\Phi, [k, \rho(\Phi)]], \rho(k))_g \sigma d\sigma = 0,$$

$$\int_X (\partial \rho, \rho(\partial \rho))_g \sigma d\sigma + \int_X ([k, \rho(\Phi)], [\Phi, \rho(k)])_g \sigma d\sigma = 0,$$

$$\int_X (\partial \rho, \rho(\partial \rho))_g \sigma d\sigma + \int_X ([\Phi, \rho(k)], [\Phi, \rho(k)])_g \sigma d\sigma = 0.$$

(39)

Both terms on the left hand side of (39) are nonnegative; hence, both vanish. Vanishing of the second term implies that $\rho(k)$ is a linear combination of lowest
weight vectors, so \( k \) is a linear combination of highest weight vectors. Now eqs. (38) imply the Bochner formula
\[
\bar{\nabla} \nabla + \nabla \bar{\nabla} = F_V \quad \Rightarrow \quad \bar{\nabla} \nabla = \bar{\nabla} \nabla + i \ast F_V .
\]
Since \( m_i \geq 1 \), the metric has constant positive curvature, \( i \ast F_V > 0 \), on the highest weight components. This implies that \( \ker \bar{\nabla} = \{0\} \) on the highest weights. Hence, the vanishing of the first term on the left hand side of (39) implies \( k \equiv 0 \). □

3.5.2. All parametrizations coincide.

**Corollary 3.5.2.** The Dolbeault, Hodge, and oper parametrizations coincide. More precisely, the following diagram commutes:

\[
\begin{array}{ccc}
Q \oplus \overline{Q} & \xleftarrow{(\phi, \beta)} & \mathcal{H}^1(C_{Dol}(G, \Phi)) \\
& \downarrow{(\psi, \alpha)} & \downarrow{(\chi, \epsilon)} \\
& \mathcal{H}^1(C_{Dol}(G, \Phi)) & \xrightarrow{(HK)^*} \mathcal{H}^1(C_{dR}(D)) \\
& & \xrightarrow{(RH)^*} \mathcal{H}^1(V) ,
\end{array}
\]

where the vertical isomorphisms are those described in Section 3.1, and \((HK)^*, (RH)^*\) are the derivatives of the Hitchin–Kobayashi and Riemann–Hilbert maps.

**Proof.** The commutivity of the identification of de Rham and oper deformations follows from the following simple remark. By construction, we have \( \sigma(q) = -q \) where \( \sigma \) is the involution defined in Section 2.2.3; thus \( \lambda(q) = -\rho(q) \). The content of the lemma is therefore in the Hitchin–Kobayashi correspondence. For simplicity, abbreviate the notation \( \phi = \phi(q) \), \( \beta = \beta(b) \), etc. We need to show that:

\[
(HK)^*(\phi, \beta) = \alpha + \psi = \beta - \lambda(\beta) + \phi + \lambda(\phi) .
\] (40)

Because the Fuchsian point is a smooth point of the moduli space, deformations are unobstructed. We may therefore find a family of Higgs bundles \( (\nabla_0^0, \Phi_0) \), passing through the Fuchsian point at \( \varepsilon = 0 \), and satisfying

\[
(\nabla_0^0, \Phi_0) = (\phi, \beta) .
\]

Moreover, the Fuchsian bundle is stable, which is an open condition, so we may assume the Higgs bundles \( (\nabla_0^0, \Phi_0) \) are stable for \( \varepsilon \) sufficiently small. Let \( \rho_\varepsilon \) be a family of harmonic metrics, whose existence is guaranteed by Theorem 2.3.2. Then:

\[
HK((\nabla_0^0, \Phi_0)) = \nabla_\varepsilon + \Phi_\varepsilon - \rho_\varepsilon(\Phi_\varepsilon) .
\]

By Theorem 3.5.1, the first variation of the harmonic metrics \( \rho_\varepsilon \) vanishes. We therefore conclude that

\[
(HK)^*(\phi, \beta) = \nabla_\varepsilon + \phi - \rho_\varepsilon(\phi)
= \beta + \rho(\beta) + \phi - \rho(\phi)
= \beta + \lambda(\beta) + \phi + \lambda(\phi)
= \alpha + \psi ,
\]

which verifies (40). The result follows. □
3.6. The tangent spaces to opers and the Hitchin component. We explain in this section our main technical tool, which we state in the de Rham picture:

**Proposition 3.6.1.** At the Fuchsian point $D$,

1. The map $\phi$ (regarded as taking values in $T_D M_{\text{dr}}(\Sigma, G)$) defines an isomorphism of $Q(X, g)$ with the tangent space $TD \text{Op}(X, G)$.
2. The map $\phi + \lambda(\phi)$ defines an isomorphism of $Q(X, g)$ with the tangent space of the Hitchin component $T_D H(\Sigma, G_R)$, which coincide with the infinitesimal version of the Hitchin parametrization.

**Proof.** The first point follows from the fact that the tangent space of opers is the set of variations of flat connections fixing the holomorphic structure $\mathcal{G}_{\text{op}}$. Now the variations $\phi(q)$ defines such a variation, and for dimensional reasons $\phi(Q(X, g))$ identifies with $TD \text{Op}(X, G)$.

The second point follows at once from Corollary 3.5.2. Indeed, the Hitchin infinitesimal parametrization is interpreted in the Dolbeault parametrization as the map $q \to \phi(q)$, but by Corollary 3.5.2, $(HK)_k(\phi(q)) = \phi(q) + \lambda(\phi(q))$. □

4. First Variation of Holonomy

The main result of this section is a Gardiner type formula for the variation of the eigenvalues of the holonomy under deformations of the Fuchsian point. Although the approach can be generalized to linear representations for all split groups we shall concentrate here on the case of $\text{SL}(n, \mathbb{R})$. Then, if $\gamma$ is a closed geodesic of length $\ell$, the largest eigenvalue $\lambda_{\gamma}$ (resp. $p$-th largest $\lambda^{(p)}_{\gamma}$) at the Fuchsian point is

$$\lambda_{\gamma} = e^{(n-1)\ell/2} \quad \text{(resp. } \lambda^{(p)}_{\gamma} = e^{(n+1-2p)\ell/2} \text{).} \quad (41)$$

Recall that in this context, we have associated to an element $q = (q_2, \ldots, q_n) \in Q(X, n)$:

- The standard oper deformation $\phi^0(q) = \sum_{k=2}^n q_k \otimes E^0_{k-1}$,
- The normalized oper deformation $\phi(q) = \sum_{k=2}^n q_k \otimes E_{k-1}$,
- The standard Hitchin deformation $\psi^0(q) = \phi^0(q) + \lambda(\phi^0(q))$,
- The normalized Hitchin deformation $\psi(q) = \phi(q) + \lambda(\phi(q))$.

Our main results in this section are the following.

**Theorem 4.0.2.** [Gardiner formula] Along the Fuchsian locus, the first variation of the largest eigenvalue $\lambda_{\gamma}$ of the holonomy along a simple closed geodesic $\gamma$ of length $\ell$, along a standard Hitchin deformation given by $q_k \in H^0(X, K^r)$, is

$$d \log \lambda_{\gamma}(\psi^0(q_k)) = \frac{(-1)^k(n-1)!}{2^{k-2}(n-k)!} \int_0^{\ell_{\gamma}} \Re \left( q_k \gamma^*, \ldots, \gamma \right) dt .$$

More generally, the first variation of the $p$-th-largest eigenvalue is

$$d \log \lambda^{(p)}_{\gamma}(\psi(q_k)) = \epsilon^{(p)}_{n,k} \int_0^{\ell_{\gamma}} \Re \left( q_k \gamma^*, \ldots, \gamma \right) dt .$$

where

$$\epsilon^{(p)}_{n,k} = \frac{(p-1)!(n-p)!}{2^{k-2}(n-k)!} \sum_{j=\max(1,k+p-n)}^{\min(k,p)} \binom{n-k}{p-j} \binom{k-1}{j-1} (-1)^{j+k+1} . \quad (42)$$
In particular, for $k = n$,

$$d \log \lambda^{(p)}_\gamma (\psi(q_n)) = (-1)^{p+n+1} \frac{(n-1)!}{2n-2} \left( \frac{n-1}{p-1} \right) \int_0^{\ell_\gamma} \mathbb{R} \left( q_n(\gamma', \ldots, \gamma') \right) \, dt .$$

(43)

**Remark 4.0.3.** For the special case of deforming in the direction of quadratic differentials, notice that $c_{n,2} = n + 1 - 2p$, and by (7), $r_2 = 1/2\pi$. It follows from Theorem 4.0.2, Proposition 2.1.1, and eq. (41), that if $q \in H^0(X, K^2)$,

$$d \ell_\gamma (\psi(q)) = \frac{1}{n} \mathbb{R} \langle q, \Theta^{(2)}_\gamma \rangle = \frac{4}{\pi} \int_X \mu(z) \cdot \Theta^{(2)}_\gamma(z) dxdy ,$$

where $\mu(z) = \overline{q(z)}(x(z))^{-1}$ is the harmonic Beltrami differential associated to $q$. This is the formula in [15, Theorem 2]. Hence, Theorem 4.0.2 is indeed a higher rank generalization of Gardiner’s formula.

Summing over in the index $p$ in (43), and using (41), we have the following

**Corollary 4.0.4.** [Variation of the trace] For a standard Hitchin deformation given by an element $q_n \in H^0(X, K^n)$, the first variation of the trace of the holonomy along a closed geodesic $\gamma$ of length $\ell_\gamma$ at the Fuchsian locus is

$$d \text{Tr}(\delta(\gamma))(\psi^0(q_n)) = 2(-1)^n (n-1)! \left( \sinh(\ell_\gamma/2) \right)^{n-1} \int_0^{\ell_\gamma} \mathbb{R} \left( q_n(\gamma', \ldots, \gamma') \right) \, dt .$$

As another corollary, one gets

**Corollary 4.0.5.** [Variation of consecutive eigenvalues] Using the notation of Theorem 4.0.2, we have

$$d \log \frac{\lambda^{(p)}_\gamma (\psi(q_n))}{\lambda^{(p+1)}_\gamma (\psi(q_n))} = A(p, k, n) \int_0^{\ell_\gamma} \mathbb{R} \left( q_n(\gamma', \ldots, \gamma') \right) \, dt ,$$

where $A(p, n, k)$ only depends on the integers $(p, n, k)$ and is nonzero.

Finally, we also have

**Theorem 4.0.6.** [Oper deformations] The first variation of the $p^{th}$-largest eigenvalue of the holonomy along a closed geodesic $\gamma$ of length $\ell_\gamma$ along a standard oper deformation given by an element $q_k \in H^0(X, K^3)$, is

$$d \log \frac{\lambda^{(p)}_\gamma (\phi^0(q_k))}{\lambda^{(p+1)}_\gamma (\phi^0(q_k))} = \frac{c^{(p)}_{n,k}}{2} \int_0^{\ell_\gamma} q_k(\gamma', \ldots, \gamma') \, dt ,$$

where $c^{(p)}_{n,k}$ is defined in (42).

In order to prove these results, in Section 4.1 we first recall a general formula computing the variation of eigenvalues by the variation of parameters method. In Section 4.2, we give an explicit description of the Fuchsian bundle along a geodesic and in Section 4.3 prove our results.
4.1. A general formula. Consider a connection $\nabla$ and a closed curved $\gamma$ in $\Sigma$ so that the holonomy of $\nabla$ along $\gamma$ has an eigenvalue $\lambda_\gamma$ of multiplicity 1. We denote by $L_\gamma$ the corresponding eigenline along $\gamma$, $H_\gamma$ the supplementary hyperplane stable by the holonomy, and $\pi$ the projection on $L_\gamma$ along $H_\gamma$. In the sequel, let us use the the general notation,

$$\dot{f} = \frac{d}{dt} \Big|_{t=0} f(t).$$

Then we have

**Lemma 4.1.1.** let $\nabla^t$ be a smooth one parameter family of connections. Then there exists (for $t$ small enough) a unique smooth function $\lambda_\gamma(t)$ such that

- $\lambda_\gamma(0) = \lambda_\gamma$,
- $\lambda_\gamma(t)$ is an eigenvalue of the holonomy of $\nabla^t$ of multiplicity 1.

Moreover,

$$\lambda_\gamma = -\lambda_\gamma \cdot \int_0^{\gamma_\ell} \text{Tr}(\nabla(s) \cdot \pi) \, ds. \quad (44)$$

**Proof.** The first part of the lemma is obvious. Observe now that the left hand side of equation (44) is invariant under gauge transform of the family $\nabla^t$. For the left hand side, let $\mathcal{V}_0 = A^t \mathcal{V}$, where $A_0 = \text{Id}$.

Then $\nabla_0 = \mathcal{V} + \nabla A$. In particular,

$$\int_0^{\gamma^\ell} \text{Tr}(\nabla_0(s) \cdot \pi) \, ds = \int_0^{\gamma^\ell} \text{Tr}(\nabla(s) \cdot \pi) \, ds + \int_0^{\gamma^\ell} \text{Tr}(\nabla(A)(s) \cdot \pi) \, ds.$$  

But, since $\pi$ is $\mathcal{V}$-parallel,

$$\int_0^{\gamma^\ell} \text{Tr}(\nabla(A)(s) \cdot \pi) \, ds = \int_0^{\gamma^\ell} \frac{d}{ds} \text{Tr}(A(s) \cdot \pi) \, ds = \left[ \text{Tr}(A(s) \cdot \pi) \right]_0^{\gamma^\ell} = 0.$$  

Thus the left hand side of (44) is also invariant under gauge transform.

We may now use a gauge transform so that the eigenline $L_\gamma^t$ and the transverse parallel hyperplane $H_\gamma^t$ are constant. Let $e$ be a nowhere vanishing section of $L_\gamma$ (after possibly taking an irrelevant double cover of $\gamma$), let $f$ be a section of $H_\gamma^t \subset \mathcal{E}^t$. Then for all $t$

$$\log \lambda_\gamma^t = -\int_0^{\gamma^\ell} \frac{\langle f(s) | \nabla^t e(s) \rangle}{\langle f(s) | e(s) \rangle} \, ds,$$

where the bracket denotes the duality. Then a standard derivation yields

$$\lambda_\gamma = -\lambda_\gamma \cdot \int_0^{\gamma^\ell} \frac{\langle f(s) | \nabla^t e(s) \rangle}{\langle f(s) | e(s) \rangle} \, ds = -\lambda_\gamma \cdot \int_0^{\gamma^\ell} \text{Tr}(\nabla(s) \cdot \pi) \, ds.$$  

We have completed the proof of formula (44). \qed
4.2. The Fuchsian bundle along a geodesic. Consider the case of the principal $SL(2)$ in $SL(n, \mathbb{C})$. Let $S \rightarrow X$ be a choice of spin structure on $X$, so that $S^{\otimes 2} = K$. Then we can describe the associated bundle to the defining representation of $SL(n, \mathbb{C})$ by (25). Let $\gamma$ be a geodesic in $X$, we then have the following structure on $E$ along $\gamma$.

(1) a harmonic metric: If $\Phi_0$ is defined as in Section 2.3.4, then the hyperbolic metric on $X$ induces a (split) hermitian structure on $E$ which solves Hitchin’s equations (19).

(2) a trivialization: We have a canonical trivialization (up to sign) of $S$ given by a section $s$ along $\gamma$ so that $s^2(\gamma) = 1$. We also have an identification of $S^*$ with $S$ using the metric and denote by $\hat{s}$ the dual section of $S$ corresponding to $s$. Both give a trivialization of $E$ along $\gamma$ by the frame

$$\{\hat{w}_p := s^p s^{n-p}\}_{p=1, \ldots, n}, \hat{w}_p \in S^{2p-n-1}.$$ 

(3) a real structure: We have a real structure on $E$ characterized by $\hat{w}_j \mapsto \overline{\hat{w}_j} := \hat{w}_{n+1-j}$.

Observe that if

$$A(w) = B(w) + B(\overline{w}), \quad \overline{C(\overline{w})} = C(w),$$

then

$$\text{Tr}(A \cdot C) = 2 \Re \left( \text{Tr}(B \cdot C) \right). \quad (45)$$

Finally let $x = \frac{1}{2}(s + \bar{s})$ and $y = \frac{1}{2}(s - is)$, so that $x + iy = s$ and $x - iy = \bar{s}$ then $x$ and $y$ are real sections of $S$.

With respect to the Chern connection $\nabla$ on $E$, the frame $\{w_p\}$ described above is parallel. The flat connection $D = \nabla + \Phi_0 - \rho(\Phi_0)$ along the geodesic $\gamma$ is may now be expressed

$$\gamma^* D = \gamma^* \nabla + (Y - X) \otimes \frac{dt}{\sqrt{2}},$$

where $Y$ and $X$ are endomorphisms satisfying

$$(-X)(\hat{w}_p) = c_p \cdot \hat{w}_{p+1},$$

$$Y(\hat{w}_p) = c_{p-1} \cdot \hat{w}_{p-1}.$$ 

where $c_p = (p(n-p)/2)!^{1/2}$ (see [23, p. 978]). The factor of $1/\sqrt{2}$ is due to the fact that the hyperbolic metric is twice the real part of the hermitian metric on $K^{-1}$. If we change to a new $\nabla$-parallel frame defined by

$$w_p = \{(p-1)!(n-p)!\}^{1/2} \hat{w}_p$$

then

$$(-X)(w_p) = \frac{1}{\sqrt{2}}(p-1) \cdot w_{p-1}, \quad (46)$$

$$Y(w_p) = \frac{1}{\sqrt{2}}(n-p) \cdot w_{p+1}. \quad (47)$$

The action of $(Y - X) \otimes dt/\sqrt{2}$ on $E$ identified with the action of a first order differential operator on homogeneous polynomials of degree $n - 1$ in $s$ and $\bar{s}$. It is given by

$$(Y - X) \otimes \frac{dt}{\sqrt{2}} = \frac{1}{2} \left( s \frac{\partial}{\partial s} + \bar{s} \frac{\partial}{\partial \bar{s}} \right) \otimes dt.$$
We observe that the system with eigenvalues \(2\), which vanishes for \(n\),

\[ (Y - X) \otimes \frac{dt}{\sqrt{2}} = \frac{1}{2} \left( x \frac{\partial}{\partial x} - y \frac{\partial}{\partial y} \right) \otimes dt. \]

In particular, let us now define

\[ u_p := x^{p-1} y^{n-p}. \]  

(48)

We observe that the system \(\{u_p\}_{p=1,\ldots,n}\) is a \(V\)-parallel eigenbasis of \((Y - X) \otimes dt/ \sqrt{2}\), with eigenvalues \(2p - n - 1)/2\). Thus,

\[ D_{\gamma} u_p = \frac{2p - n - 1}{2} u_p . \]  

(49)

Furthermore, if \(\gamma\) is a closed geodesic of length \(\ell\), the \(p\)-th largest for the holonomy of \(D\) is given by (41).

The lines \(L_p\) generated by \(u_p\) form a parallel frame along any geodesic for both \(V\) and \(D\). Let \(p\) be the projection on the eigenline \(L_j\) orthogonal to \(\bigoplus_{p \neq j} L_j\). For \(k = 1, \ldots, n - 1\), let \(E^0_k = (-\sqrt{2} Y)^k\). The proof of the following proposition is in the Appendix (see Proposition A.3).

**Proposition 4.2.1.** For all \(k, p\), we have \(\text{Tr}(E^0_k \pi_p) = -\text{Tr}(E^0_k \pi_p)\). Moreover,

\[ \text{Tr}(E^0_k \pi_1) = \frac{(-1)^k (n - 1)!}{2^k (n - k - 1)!} \]  

(50)

\[ \text{Tr}(E^0_k \pi_p) = \frac{(p - 1)! (n-p)!}{2^k (n-k-1)!} \sum_{j=\max(0,k+p-n)}^{\min(k,p-1)} \left( n-k-1 \right)^2 \left( \begin{array}{c} k \end{array} \right)^2 (-1)^{j+k}. \]  

(51)

It is interesting to remark that \(\text{Tr}(E^0_k \pi_1) \neq 0\) for all \(k\) and \(n\). On the other hand, for \(p = n - 1, k \leq n - 2\) the second equation yields

\[ \text{Tr}(E^0_k \pi_p) = \frac{(n - 2)!}{2^k (n-k-1)!} \sum_{j=1-k}^{k} \left( \frac{n-k-1}{n-j-1} \right)^2 \left( \begin{array}{c} k \end{array} \right)^2 (-1)^{j+k} \]  

\[ = \frac{(n - 2)!}{2^k (n-k-1)!} \left( \left( \frac{n-k-1}{n-k-2} \right) - \left( \frac{k}{k-1} \right)^2 \right) \]  

\[ = \frac{(n - 2)!}{2^k (n-k-1)!} \left( n-k^2 - k - 1 \right) , \]

which vanishes for \(n = k^2 + k + 1\) (e.g. \(n = 3, k = 1\)).

4.3. **Proof of Theorems 4.0.2 and 4.0.6.** We now have the following consequence of the calculations in the previous section. For the standard Hitchin variation \(\dot{q}_k(q_k)\), the corresponding variation of flat connection is

\[ \dot{D} = q_k \otimes E^0_{k-1} - q_k \otimes E^0_{k-1} . \]

In particular, for any endomorphism \(B\) commuting with the real involution, as a function on \(UX\),

\[ \int_0^t \text{Tr} \left( \dot{D} \cdot B \right) \, ds = \int_0^t 2 \text{Re} \left( q_k \right) \text{Tr}(E^0_{k-1} \cdot B) \, ds . \]  

(52)

The theorem now follows from Proposition 4.2.1 and eqs. (52) and (44).
We note that the same argument applies to oper deformations at the Veronese oper, where now \( \mathring{D} = q_k \otimes E^0_{k-1} \). In particular, we immediately obtain 4.0.6.

5. The Symplectic Structure and Twist Deformations

The purpose of this section is to investigate some relations between the Hitchin parametrization along the Fuchsian locus and the “symplectic nature of the fundamental group of a surface,” to paraphrase Goldman [16].

(1) In Corollary 5.1.2, we show that via the (normalized) Hitchin parametrization the symplectic structure of the \( L^2 \) pairing coincides with (a multiple of) the Atiyah–Bott–Goldman symplectic structure.

(2) In Section 5.2, we relate the twist or bending deformations to relative Poincaré series, in Section 5.2 we prove the Reciprocity Theorem 5.3.1.

(3) In Theorem 5.4.1 we study the Hamiltonian flow of the length functions in term of twists and the relative Poincaré series.

5.1. Symplectic structure. Recall [2, 16] that the de Rham moduli space \( \mathcal{M}_{dR}(\Sigma, G) \) carries the Atiyah–Bott–Goldman symplectic form \( \omega^0_G \) given by:

\[
\omega^0_G(A, B) = \int_{\Sigma} (A \wedge B)_h .
\]  

(53)

It will also be convenient to consider the normalized symplectic form:

\[
\omega_G(A, B) = \frac{1}{d(g)} \omega^0_G(A, B) ,
\]  

(54)

where \( d(g) \) is the Dynkin index of a principal \( \mathfrak{sl}(2) \hookrightarrow \mathfrak{sl}(n) \) subalgebra (see (13)). Then restricting to the Fuchsian locus one has \( \omega_G = \omega_2 \), where \( \omega_2 \) is the standard Weil–Petersson symplectic form on Teichmüller space. In general, when \( G = \text{SL}(n, \mathbb{C}) \), we denote the Atiyah–Bott–Goldman symplectic form by \( \omega^0_n \) and the normalized one by \( \omega_n \).

As will become clear from the computations below, the Hitchin space \( \mathcal{H}(\Sigma, G) \) a symplectic submanifold of \( \mathcal{M}_{dR}(\Sigma, G) \), and the oper moduli space \( \mathcal{O}(X, G) \) is Lagrangian. To illustrate, let us evaluate \( \omega^0_G \) on real de Rham deformations.

**Lemma 5.1.1.** Let \( q_1, q_2 \) be differentials of weight \( m_i + 1 \) and \( m_j + 1 \), respectively. Then if \( i = j \),

\[
\omega_G(\psi(q_1), \psi(q_2)) = -2\Im \langle q_1, q_2 \rangle .
\]

if \( i \neq j \),

\[
\omega_G(\psi(q_1), \psi(q_2)) = 0 .
\]

**Proof.** Observe first that if \( a = \alpha \otimes A \) and \( b = \beta \otimes B \) then \( (a \wedge b)_h = \alpha \wedge \beta \cdot (A, B)_h \). Thus if \( v \in K^{m+1} \) and \( w \in K^{m+1} \)

\[
((v \otimes E^0_i) \wedge (\rho(w \otimes e_j)))_h = i\delta_{ij}d(g) \langle v, w \rangle \, d\sigma ,
\]

\[
((v \otimes E^0_i) \wedge (w \otimes e_j))_h = 0 .
\]  

(55)
It then follows that
\[
\omega^0_G(\psi(q_1), \psi(q_2)) = \int_X \left( (\langle q_1 \otimes E^0_i \rangle - \rho(q_1 \otimes e_i)) \right. \\
\left. \wedge (q_2 \otimes e_j - \rho(q_2 \otimes e_j)) \right) \delta (xy) \\
= -\int_X \langle q_1 \otimes E^0_i \rangle \wedge \rho(q_2 \otimes e_j) \delta - \int_X \rho(q_1 \otimes E^0_i) \wedge (q_2 \otimes e_j) \delta \\
= i \cdot \delta_{ij} \cdot d(\rho) \left( -\int_X \langle q_1, q_2 \rangle d\sigma + \int_X \langle q_2, q_1 \rangle d\sigma \right) \\
= -\delta_{ij} \cdot d(\rho) \cdot \mathcal{J} \langle q_1, q_2 \rangle_X .
\]

As an immediate consequence of this lemma we have the following corollary, which explains the normalization that we have chosen.

**Corollary 5.1.2.** At the origin in the Hitchin base, the pull-back of the normalized Atiyah–Bott–Goldman symplectic form \( \omega_G \) on \( \mathcal{H}(\Sigma, G) \) by the normalized Hitchin section coincides with the symplectic form on \( Q(X, g) \) associated to the \( L^2 \)-metric.

5.2. **Twist deformations.** Let \( D \) be a flat connection on a bundle \( E \) over \( \Sigma \). Let \( \gamma \) be a closed simple curve, whose holonomy is \( g \), and let finally \( h \) be an element in the Lie algebra \( \mathfrak{z}(g) \) of the centralizer of \( g \). We can then construct an element \( \tau_\gamma(h) \) in \( H^1_\partial(\Sigma, \text{ad}(E)) \) by the following construction. Choose a tubular neighborhood \( U = S^1 \times [0, 1] \) of \( \gamma \), let \( f \) be a function on \( U \) which is constant equal to 1 on a neighborhood of \( S^1 \times \{1\} \) and 0 on a neighborhood of \( S^1 \times \{0\} \). Observe finally that the parallel sections \( \tilde{h} \) of \( \text{ad}(E) \) restricted to \( U \) are identified canonically to elements \( h \in \mathfrak{z}(g) \). Then we define
\[
\tau_\gamma(h) := df \cdot \tilde{h} = D(f \cdot \tilde{h}) \in \Omega^1(U, \text{ad}(E)) ,
\]
and observe that \( \tau_\gamma(h) \) extends (by zero) to all of \( \Sigma \).

For Hitchin representations, the monodromy is regular [25] and so using the previous notation, \( \mathfrak{z}(g) \) is identified with a Cartan subalgebra \( \mathfrak{h} \). Thus the previous construction associates to every element \( h \in \mathfrak{h} \) a vector field \( \tau_\gamma(h) \), called the **twist vector field**, on the Hitchin component.

At the Fuchsian locus we furthermore have an identification \( \mathfrak{z}(g) = \mathfrak{z}(X-Y)_R \). This leads to

**Proposition 5.2.1.** For any \( B \in T_{\mathfrak{D}}\mathcal{M}_{\mathfrak{D}}(\Sigma, G) \) we have
\[
\omega^0_G(B, \tau_\gamma(h)) = \int_\gamma (B, h) \delta .
\]

In particular, if \( q_k \) is a \( k \)-holomorphic differential, \( k = m_i + 1 \) for some \( i \), then
\[
\omega_G(\psi(q_k), \tau_\gamma(h)) = \frac{-2r_k \cdot (h_i e_i) \delta \cdot \mathcal{J}(q_k, i \cdot \Theta_\gamma)}{d(\rho)} ,
\]
where \( r_k \) is given in (7).
Proof. The first assertion follows from eqs. (53) and (56) from integration by parts (recall that $DB = 0$). For the second, since $\lambda(h) = h$, and using Proposition 2.1.1,

$$\omega_G^0(\psi(q_k), \tau_\gamma(h)) = \int_\gamma (h, (q_k \otimes e_{k-1} - \rho(q_k \otimes e_{k-1}))_h$$

$$= (h, e_i)_h \left( \int_\gamma q_k + \int_\gamma \frac{\tau}{h} \right)$$

$$= (h, e_i)_h r_k \cdot \left( \langle q_k, \Theta^{(k)}_\gamma \rangle + \langle q_k, \Theta^{(k)}_\gamma \rangle \right)$$

$$= -2(h, e_i)_h r_k \cdot \Im \langle q_k, i \cdot \Theta^{(k)}_\gamma \rangle .$$

The statement now follows from (54). \hfill \Box

Recall Lemma 2.2.1 and the notion of a principal basis.

**Corollary 5.2.2.** The space of twist deformations at a simple geodesic $\gamma$ is the $\mathbb{R}$-span of the relative $k$-Poincaré series of $\gamma$, for degrees $k = m_1 + 1, \ldots, m_i + 1$. More precisely, if $(h_i)_{i \in 1, \ldots, l}$ is the principal basis of $\tilde{z}(X - Y)_{\mathbb{R}}$, then

$$\tau_\gamma(h_i) = \frac{r_{m_i+1}}{d(g)} \cdot \psi \left( i \cdot \Theta^{(m_i+1)}_\gamma \right) .$$

Proof. Write

$$\tau_\gamma(h) = \sum_{i=1}^l \psi(Q_i) ,$$

where $Q_i \in H^0(X, K^{m_i})$ (as a real vector space), and let $q_k$ be an arbitrary $k$-differential, $k = m_i + 1$ for some $i$. Then using Lemma 5.1.1 and eq. (57), we find

$$\omega_G(\psi(q_k), \tau_\gamma(h)) = -2 \Im \langle q_k, Q_i \rangle = -2r_k \cdot (h, e_i)_h \cdot \Im \langle q_k, i \cdot \Theta^{(k)}_\gamma \rangle ,$$

and,

$$\omega_G(\psi(iq_k), \tau_\gamma(h)) = -2(e_i, f_i)_h \Re \langle q_k, Q_i \rangle = -2r_k \cdot (h, e_i)_h \cdot \Re \langle q_k, i \cdot \Theta^{(k)}_\gamma \rangle ,$$

from which

$$\langle q_k, Q_i \rangle = \frac{r_k \cdot (h, e_i)_h}{d(g)} \langle q_k, i \cdot \Theta^{(k)}_\gamma \rangle .$$

Since $k$ and $q_k$ were arbitrary, the result follows immediately. \hfill \Box

5.3. Reciprocity: twists and lengths. By Corollary 5.2.2, it follows that each relative Poincaré series $i \cdot \Theta^{(k)}_\gamma$ corresponds to a unique twist deformation $\psi(i \cdot \Theta^{(k)}_\gamma)$. Recall that $\lambda^{(p)}$ denotes the $p$-th largest eigenvalue for the holonomy about $\gamma$.

**Theorem 5.3.1.** [Reciprocity of the twist deformation] For any $k$, $p$, and any simple closed geodesics $\alpha, \beta$, then at the Fuchsian locus,

$$\text{d} \log \lambda^{(p)}_\alpha \left( \psi(i \cdot \Theta^{(k)}_\beta) \right) = -\text{d} \log \lambda^{(p)}_\beta \left( \psi(i \cdot \Theta^{(k)}_\alpha) \right) .$$
Theorem 6.1.1. \[ P \]

Fix a simple closed curve $\omega$ with respect to the normalized Atiyah-Bott-Goldman symplectic form orthogonal with respect to the pressure metric. Moreover two deformations associated to holomorphic differentials of degree $k$ on $X$, and let $q$ be a holomorphic differential of degree $k$ on $X$, and let $\psi^0(q)$ be the associated standard deformation, then

$$ P_\delta (\psi^0(q), \psi^0(q)) = \frac{1}{2^{k-1} \pi |\lambda(X)|} \left[ \frac{(k-1)!(n-1)!}{(n-k)!} \right]^2 \langle q, q \rangle_X .$$

Moreover two deformations associated to holomorphic differentials of different degrees are orthogonal with respect to the pressure metric.

5.4. Hamiltonian functions and twists. Recall that the Hamiltonian vector field $H_f$ of a $C^1$ function $f$ on a symplectic manifold $(M, \omega)$ is defined by

$$ df(V) = \omega(V, H_f),$$

for any tangent vector field $V$ on $M$.

Theorem 5.4.1. Fix a simple closed curve $\gamma$. Then along the Fuchsian locus the Hamiltonian vector field of the function

$$ \log \lambda_\gamma^0 : \mathcal{H}(\Sigma, n) \to \mathbb{R} $$

with respect to the normalized Atiyah-Bott-Goldman symplectic form $\omega_n$ is given by

$$ H_{\log \lambda_\gamma^0} = \sum_{k=2}^{n} c_{i,n}^{-1} r_k \cdot \eta_{k-1} \cdot \frac{\psi(i \cdot \Theta_\gamma^0)}{2} , \quad (58) $$

where $c_{i,n}$ is defined in (42), $r_k$ in (7), and $\eta_{k-1}$ in (18).

Proof. Let $q_k$ be a $k$-differential. From Theorem 4.0.2 and Proposition 2.1.1,

$$ d \log \lambda_\gamma^0 (\psi(q_k)) = c_{i,n}^{-1} \cdot \eta_{k-1} \cdot \int_{\gamma} q_k \cdot \frac{r_k \cdot \eta_{k-1} \cdot \mathcal{I} \langle q_k, i \cdot \Theta_\gamma^0 \rangle}{2} . $$

On the other hand, from Lemma 5.1.1,

$$ \omega_n(\psi(q_k), \psi(i \cdot \Theta_\gamma^0)) = -2 \mathcal{I} \langle q_k, i \cdot \Theta_\gamma^0 \rangle . $$

Since $k$ and $q_k$ are arbitrary, the result follows. \qed

Remark 5.4.2. We point out that for the highest eigenvalue $\lambda_\gamma$ (i.e. $p = 1$), the coefficients in the expression (58) are nonzero for all $2 \leq k \leq n$ (see Theorem 4.0.2).

6. The Variance and the Pressure Metrics

6.1. The pressure metric. In the following paragraphs we shall recall some of the results of [8], where the authors introduced the pressure metric $P$ on the Hitchin component. We will prove the following

Theorem 6.1.1. [Pressure metric for standard deformations] Let $\delta$ be a Fuchsian representation into $\text{SL}(n, \mathbb{R})$ associated to a Riemann surface $X$ with conformal hyperbolic metric. Let $q$ be a holomorphic differential of degree $k$ on $X$, and let $\psi^0(q)$ be the associated standard deformation, then

$$ P_\delta (\psi^0(q), \psi^0(q)) = \frac{1}{2^{k-1} \pi |\lambda(X)|} \left[ \frac{(k-1)!(n-1)!}{(n-k)!} \right]^2 \langle q, q \rangle_X . $$
The pressure metric for normalized deformations $\psi(q) = \eta_{k-1} \cdot \psi^0(q)$, when $q$ has degree $k$ now follows:

**Corollary 6.1.2.** [Pressure metric for normalized deformations] Let $\delta$ be a Fuchsian representation into $\text{SL}(n, \mathbb{R})$ associated to a Riemann surface $X$ with conformal hyperbolic metric. Let $q$ be a holomorphic differential of degree $k$ on $X$, and let $\psi(q)$ be the associated normalized deformation, then

$$P(\psi(q), \psi(q)) = \frac{[(n-1)!]^2}{2^{n-1} \pi |\chi(X)|}(n+1) \frac{(2k-1)!}{(n+k-1)! (n-k)!} \langle q, q \rangle_X.$$ 

Moreover two deformations associated to holomorphic differentials of different degrees are orthogonal with respect to the pressure metric.

Corollary 1.0.3, stated in the introduction, follows at once from Corollary 6.1.2 and Corollary 5.1.2.

The proof of the theorem and the structure of this section fall into several parts. First, we recall the definition of the pressure metric for projective Anosov representation in Section 6.2. More importantly, we shall introduce the notion of variation of paramatrization for a deformation of representation in Section 6.2.1 and state eq. 59 which relates the pressure metric and the variation of reparametrization. In Section 6.2.2, we identify the variation of reparametrization associated to the standard variation $\psi^0(q)$ for $q$ a holomorphic differentials, and we deduce eq. (60), which identifies the pressure metric as a multiple of the variance metric. The rest of the section is devoted to the proof of Theorem 6.3.1, which computes the variance metric for holomorphic differentials in terms of the $L^2$-metric. In a concluding Section 6.6, we sketch how one might extend these results to other pressure metrics.

The pressure metric has been discussed in the quasi-Fuchsian context in [28, 9, 7], and the construction of the pressure metric in [28] ties in Wolpert’s approach of the identification of the Thurston metric with the Weil–Petersson metric [42].

### 6.2. Projective Anosov representations.

Introduced in [25] and studied further in [18] and [8], a projective Anosov representation $\delta$ of a hyperbolic group $\Gamma$ in $\text{SL}(n, \mathbb{R})$ is characterized by the following features:

- A **spectrum** that is a map which associates to every nontrivial element $\gamma$ of $\Gamma$, the number $\ell_\gamma(\delta) := \log |\lambda_\gamma(\delta)|$, where $\lambda_\gamma(\delta)$ is the highest eigenvalue (in modulus) of $\delta(\gamma)$,
- An **entropy** defined as
  $$h(\delta) := \lim_{T \to \infty} \frac{1}{T} \log (\# L_\delta(R)),$$
  where $L_\delta(R) := \{ \gamma \in \Gamma \mid \ell_\gamma(\delta) \leq T \}$.
- Moreover, one can define the intersection of two projective Anosov representations $\delta_1$ and $\delta_2$ as
  $$I(\delta_1, \delta_2) := \lim_{T \to \infty} \frac{1}{\# L_\delta(R)} \sum_{\gamma \in L_\delta(R)} \ell_\gamma(\delta_2) \ell_\gamma(\delta_1).$$

and the renormalized intersection $J(\delta_1, \delta_2) := \frac{h(\delta_2)}{h(\delta_1)} I(\delta_1, \delta_2)$.

As proved by the first author, Hitchin representations are examples of projective Anosov representations, and indeed they were the initial motivation for the definition [25]. The article [8] introduces entropy and intersections.
which by definition is the Hessian at a representation $\delta_0$ of the function $\delta \mapsto f(\delta_0, \delta)$, is a consequence.

6.2.1. Variation of reparametrization. In [8], an explicit formula for the pressure metric was given as follows. To simplify the exposition, we work in the context of deformations of Fuchsian representations. Let $[\delta_t]_{t \in (-\varepsilon, \varepsilon)}$ be a family of deformations of a Fuchsian representation $\delta_0$, associated to a hyperbolic surface $X$. A Hölder function $f$ on $U_X$ is a variation of reparametrization associated to $\delta_0$ if for every $\gamma$ in $\pi_1(\Sigma)$,

$$\int_\gamma f \, ds = \frac{d}{dt} \bigg|_{t=0} \ell_\gamma(\delta_t).$$

Then the pressure metric is given by the variance of $f$ that is

$$\mathbf{P}_0 \left( \delta_0, \delta \right) = \text{Var}(f) := \lim_{t \to \infty} \frac{1}{t} \int_{U_X} \left( \int_0^t f(\phi_s(x)) \, ds \right)^2 \, d\mu(x),$$

where $\mu$ is the Liouville measure of $U_X$.

6.2.2. A consequence of Gardiner Formula. From the Gardiner formula, Theorem 4.0.2, and the definition of variation of reparametrization in Section 6.2.1, we immediately obtain that for $q$ be a holomorphic $k$-differential, $\hat{q}$ the associated complex valued function on $U_X$ and $\check{q} = \Re(\hat{q})$, the function

$$\frac{(-1)^k(n-1)!}{2^{k-2}(n-k)!} \cdot \check{q},$$

is the variation of parametrization associated to the standard deformation $\psi^0(q)$. Thus it follows from (59), that

$$\mathbf{P}(\psi^0(q), \psi^0(q)) = \left[ \frac{(n-1)!}{2^{k-2}(n-k)!} \right]^2 \cdot \text{Var}(q),$$

where by a slight abuse of notation we write $\text{Var}(q) := \text{Var}(\check{q})$. We note again that Curt McMullen has treated the case of quadratic differentials in [28].

6.3. Variance and the $L^2$ metric. The goal of this section is to prove

**Theorem 6.3.1.**

$$\text{Var}(q) = \frac{2^{k-2}[(k-1)!]^2}{2\pi [\chi(X)]} \langle q, q \rangle_X.$$  

Moreover, if two holomorphic forms have different degrees, then they are orthogonal for the bilinear form underlying the quadratic form $\text{Var}$.

This theorem, together with eq. (60), concludes the proof of Theorem 6.1.1.

6.4. Preliminary computations.
6.4.1. Hypergeometric integrals. For nonnegative integers \( m \) and \( d \), and \( R \in [0,1] \), the following incomplete Beta functions will play a technical role in the sequel:

\[
I_{n,d}(R) := \int_0^R t^m (1-t^2)^d \, dt .
\]

We postpone the proof of the following result to the Appendix.

**Theorem 6.4.1.**

\[
\lim_{R \to 1} \left( \frac{1}{|\log (1-R)|} \sum_{n=0}^\infty \frac{(n+2p-1)!}{n!} I_{n,p}^2(R) \right) = 2^{2p-2}(p-1)^2 .
\]

6.4.2. In the Poincaré disk model. Let us define the function

\[
r(R) := \frac{1}{2} \log \left( \frac{1 + R}{1 - R} \right) .
\]

We recall that if \( d_H \) is the hyperbolic distance in the Poincaré disk model, then

\[
d_H(0, R e^{it}) = r(R) .
\]

Let

\[
q(z) = dz^k \sum_{n=0}^\infty a_n z^n ,
\]

be a holomorphic \( k \)-differential on the Poincaré disk. Let us also consider the real valued function

\[
\tilde{q}(z) := \Re \left( q(z) \left( \frac{\partial}{\partial r}, \ldots, \frac{\partial}{\partial r} \right) \right) \quad \text{where} \quad \frac{\partial}{\partial r} = (1 - R^2) \frac{\partial}{\partial R} .
\]

If furthermore \( q \) and \( q^0 \) are two holomorphic forms of degree \( k \) and \( k_0 \), respectively, we will consider the complex valued function

\[
q \odot q^0(z) = q(z) \left( \frac{\partial}{\partial r}, \ldots, \frac{\partial}{\partial r} \right) q^0(z) \left( \frac{\partial}{\partial r}, \ldots, \frac{\partial}{\partial r} \right) .
\]

The main result in this section computes integrals related to the holomorphic differentials \( q \) and \( q^0 \).

**Proposition 6.4.2.** We have

\[
\frac{1}{2\pi} \int_0^{2\pi} \left( \int_0^R \left| q(T e^{i\theta}) \right|^2 \frac{dT}{1 - T^2} \right) d\theta = \sum_{n \in \mathbb{N}} a_n \overline{a}_{n'} I_{2n,2k-1}(R) ,
\]

\[
\frac{1}{2\pi} \int_0^{2\pi} \left( \int_0^R \left| \tilde{q}(T e^{i\theta}) \frac{dT}{1 - T^2} \right|^2 \right) d\theta = \sum_{n \in \mathbb{N}} a_n \overline{a}_{n'} I_{n,k}^2(R) .
\]

Moreover, if \( q^0 = (dz)^{k_0} \sum_{n=0}^\infty b_n z^n \) is of degree \( k_0 \), then

\[
\frac{1}{2\pi} \int_0^{2\pi} \left( \int_0^R \left| \tilde{q}(T e^{i\theta}) \frac{dT}{1 - T^2} \right|^2 \right) \left( \int_0^R \left| q^0(T e^{i\theta}) \frac{dT}{1 - T^2} \right|^2 \right) d\theta
\]

\[
= \sum_{n=0}^\infty \Re \left( b_{n+k-k_0} \overline{a}_n \right) I_{n,k-1}(R) I_{n+k-k_0,k-1}(R) .
\]
Finally
\[
\frac{1}{2\pi} \int_0^{2\pi} \left( \int_0^R q \odot q^0(Te^{i\theta}) \frac{dT}{1 - T^2} \right) d\theta = \frac{1}{2} \sum_{n=0}^{\infty} b_{n+k-k_0} a_{n-l_{2n+k-k_0+1}}(R) . \tag{66}
\]

Proof. We first prove
\[
\| q(Re^{i\theta}) \|^2 = 2^k(1 - R^2)^{2k} \left( \sum_{n\in\mathbb{N}} a_n T_n(R) + \sum_{k=1}^{\infty} \sum_{n\in\mathbb{N}} a_n a_{n+m} R^{n+m} e^{i\theta} \right) . \tag{67}
\]

The hyperbolic metric of the Poincaré disk model is
\[
\sigma = \frac{1}{1 - R^2} (dx^2 + dy^2) .
\]

Thus \( \|dz\|^2 = 2(1 - R^2)^2 \) (recall our conventions from Section 2.1.1). It follows that \( \|f(z)dz\|^2 = |f(z)|^22^k(1 - R^2)^{2k} \), and hence (67). Now (63) is an immediate consequence of (67).

Let us now move to the proof of (64). By (61), we have that
\[
dz \left( \frac{\partial}{\partial \sigma} \right) \bigg|_{Re^{i\theta}} = (1 - R^2) \cdot dz \left( \frac{\partial}{\partial R} \right) \bigg|_{Re^{i\theta}} = (1 - R^2)e^{i\theta} .
\]

Thus
\[
\partial(Re^{i\theta} \hat{\phi}_f(x)) = \Re \left( \sum_{n=0}^{\infty} a_n T_n(R) e^{i(n+k)\theta} \right) .
\]

It follows that
\[
\int_0^R \partial(Re^{i\theta}) \frac{dT}{1 - T^2} = \Re \left( \sum_{n=0}^{\infty} a_n e^{i(n+k)\theta} \int_0^R T^{n-1}(1 - T^2)^{k-1} dT \right) = \frac{1}{2} \sum_{n=0}^{\infty} I_{n+k-1}(R) \left( a_n e^{i(n+k)\theta} + \pi a_n e^{-i(n+k)\theta} \right) .
\]

Thus we obtain after taking the square and integrating over \( \theta \), that
\[
\frac{1}{2\pi} \int_0^{2\pi} \left( \int_0^R \partial(Re^{i\theta}) \frac{dT}{1 - T^2} \right)^2 d\theta = \frac{1}{2} \sum_{n=0}^{\infty} a_n a_{n-l_{2n+k-k_0+1}}(R) .
\]

Moreover, if \( q \) of degree \( k \), and \( q^0 \) of degree \( k_0 \geq k \), then the same argument yields
\[
\frac{1}{2\pi} \int_0^{2\pi} \left( \int_0^R \partial(Re^{i\theta}) \frac{dT}{1 - T^2} \right) \left( \int_0^R \partial(Re^{i\theta}) \frac{dT}{1 - T^2} \right) d\theta = \frac{1}{2} \sum_{n=0}^{\infty} \Re \left( b_{n+k-k_0} a_{n-l_{2n+k-k_0+1}}(R) \right) .
\]

Let us move to (66). We have that
\[
q \odot q^0(Re^{i\theta}) = \left( \sum_{n=0}^{\infty} a_n T_n(R) \right) \left( \sum_{m=0}^{\infty} b_m R^m(1 - R^2)^{k_0} e^{i(m+k_0)\theta} \right) .
\]

Thus again, integrating over the circle yields
\[
\frac{1}{2\pi} \int_0^{2\pi} q \odot q^0(Re^{i\theta}) d\theta = \sum_{n=0}^{\infty} a_n b_{n+k-k_0} R^{2n+k-k_0}(1 - R^2)^{k+k_0} .
\]
Further integration now gives
\[ \frac{1}{2\pi} \int_0^{2\pi} \int_0^R q \circ q^0(Re^{i\theta}) \frac{dT d\theta}{1 - T^2} = \sum_{n=0}^\infty a_n b_{n+k-k_0} \int_0^R T^{2n+k-k_0} (1 - T^2)^{k-k_0-1} dT \]
\[ = \sum_{n=0}^\infty a_n b_{n+k-k_0} I_{2n+k-k_0,k+k_0+1}(R). \]

\[ \square \]

6.4.3. Functions on the unit tangent bundle. Observe that any \( x \) in \( UX \) gives an identification of \( \tilde{X} \) the universal cover of \( X \) with \( \mathbb{D} \) the Poincaré disk model: the identification sends \( \pi(x) \) to 0 and \( x \) to \((1, 0)\). We can then write the analytic expansion of a \( k \)-differential \( q \) in these coordinates:
\[ q_x(z) = dz^k \sum_{n=1}^\infty a_n(x)z^n. \] (68)

In particular, in this way we obtain complex valued functions \( a_n \) on \( UX \). The functions \( a_n \) contain all the information about \( q \) with some redundancy. We can use the family of holomorphic forms \( q_x \) to describe the functions that we wish to study on \( UX \). As we have done previously, we also view \( q \) as a function on \( UX \), homogeneous of degree \( k \). Let us furthermore recall that the circle acts on both \( UX \) and the Poincaré disk model (hence on holomorphic forms). With this understood, we have the following

**Proposition 6.4.3.**
\[ \|q(\varphi_r(\xi))\|^2 = \|q_x(Re^{i\theta})\|^2 \] (69)
\[ \bar{q}(\varphi_r(\xi)) = \bar{q}_x(Re^{i\theta}) \] (70)
\[ q \circ q^0(\varphi_r(\xi)) = q_x \circ q^0_x(Re^{i\theta}). \] (71)

**Proof.** By construction,
\[ q_x = (e^{i\theta})^*q_x. \]

Therefore,
\[ \|q_x(\xi)\|^2 = \|q_x(Re^{i\theta})\|^2, \]
\[ q_x(\xi) = \tilde{q}_x(Re^{i\theta}), \]
\[ q_x \circ q^0_x(\xi) = q_x \circ q^0_x(Re^{i\theta}). \]

We now describe the action of the geodesic flow. We have
\[ \|q(\tilde{\varphi}_r(x))\|^2 = \|q_x(R)\|^2, \]
\[ \bar{q}(\tilde{\varphi}_r(x)) = \bar{q}_x(R), \]
\[ q \circ q^0(\tilde{\varphi}_r(x)) = q_x \circ q^0_x(R). \]

Combining the two actions, one gets the proposition. \( \square \)
6.4.4. The Hilbert norm and the analytic expansion. Let
\[ A_n := \int_{UX} a_n(x) \overline{a_n(x)} \, d\mu(x). \] (72)

We now prove

**Proposition 6.4.4.** For any \( 1 > R > 0, \)
\[ \frac{1}{2\pi|X|} \cdot ||\eta||_X^2 = \frac{2^{k+1}}{\log \frac{R}{1+R}} \left( \sum_{n=0}^{\infty} A_n \cdot I_{2n,2k-1}(R) \right). \] (73)

Moreover, if \( q^0 \) is a holomorphic \( k_0 \)-differential on \( X \) with associated functions \( b_n \) on \( UX \), then for any \( 1 > R > 0, \)
\[ \int_{UX} q \circ q_0 \, d\mu = \frac{2}{\log \frac{R}{1+R}} \left( \sum_{n=0}^{\infty} I_{2n+k-k_0,k+k_0-1}(R) \cdot \int_{UX} a_n(x) b_{n+k-k_0}(x) \cdot d\mu(x) \right). \] (74)

As a corollary, we obtain

**Corollary 6.4.5.** The following hold:
\[ A_n = \frac{2^{-k}}{2\pi|X|} \left( \frac{n+2k-1}{n} \right) \cdot ||\eta||_X^2; \] (75)
and for \( k_0 \neq k, \)
\[ 0 = \int_{UX} a_n(x) b_{n+k-k_0}(x) \cdot d\mu(x). \] (76)

We can now proceed to the

**Proof of Proposition 6.4.4.** We view \( ||\eta|| \) as a function on \( UX \). Then, by the normalization of the Liouville measure,
\[ ||\eta||_X^2 = \int_X ||\eta||^2 \, d\sigma = 2\pi|X| \int_{UX} ||\eta(x)||^2 \, d\mu(x). \]

Since the Liouville measure is invariant by the geodesic flow, we have for all \( r, \)
\[ ||\eta||_X^2 = \frac{2\pi|X|}{r} \int_{UX} \int_0^r ||\eta(\phi_t(x))||^2 \, dt \cdot d\mu(x). \]

Let us further use the action of the circle on \( UX \) and the invariance of the Liouville measure to get
\[ \int_X ||\eta||^2 \, d\sigma = \frac{|X|}{r} \int_{UX} \int_0^{2\pi} \int_0^r ||\eta(\phi_t(e^{i\theta}x))||^2 \, dt \cdot d\theta \cdot d\mu(x). \] (77)

Using the notation defined in eq. (68), we have by eq. (69) that if \( t = r(T), \)
\[ ||\eta(\phi_t(e^{i\theta}x))||^2 = ||\eta_T(e^{i\theta}x)||^2. \] (78)

By (63) and using that \( dt = \frac{d\theta}{1-r^2}, \)
\[ \frac{1}{2\pi} \int_0^{2\pi} \int_0^R ||\eta_T(e^{i\theta}x)||^2 \frac{dT \, d\theta}{1-r^2} = 2^k \sum_{n \in \mathbb{N}} a_n(x) \overline{a_n(x)} \cdot I_{2n,2k-1}(R). \]

Combining eqs. (77) and (78) with a further integration over \( UX \), the first result follows. The second equation follows from the same ideas using eqs. (71) and (66). \( \square \)
We move to the proof of Corollary 6.4.5.

Proof. Near $R = 0$, $I_{m,d}(R) \sim R$. Taking the limit when $R \to 0$ of (73) therefore leads to

$$2\pi|\chi(X)| \cdot 2^k A_0 = \|q\|_X^2,$$

Thus, we can rewrite (73) as

$$r(R) \cdot A_0 = \sum_{n=0}^{\infty} A_n \int_0^R T^{2n} (1 - T^2)^{2k-1} dT.$$  

Taking the derivative with respect to $R$ leads to

$$(1 - R^2)^{-1} \cdot A_0 = \sum_{n=0}^{\infty} A_n \cdot R^{2n} (1 - R^2)^{2k-1}.$$  

Taking $L = R^2$, we get

$$\sum_{n=0}^{\infty} A_n L^n = \frac{A_0}{(1 - L)^{2k}}.$$  

From the asymptotic expansion

$$\frac{1}{(1 - L)^{N}} = \sum_{m=0}^{\infty} \binom{m + N - 1}{m} \cdot L^m,$$

(which follows inductively by differentiation) we obtain that

$$A_n = \binom{n + 2k - 1}{n} \cdot A_0.$$

Eq. (75) now follows from (79). We note that $q \odot q_0(e^{i\theta}x) = e^{i(k-k_0)\theta} q \odot q_0(x)$. Thus, for $k_0 \neq k$,

$$\int_{UX} q \odot q_0 \cdot d\mu = 0.$$  

Then (74) yields

$$0 = \sum_{n=0}^{\infty} I_{2n+k-k_0,k+k_0-1}(R) \cdot \int_{UX} a_n(x) b_{n+k-k_0}(x) \cdot d\mu(x).$$  

Taking the derivatives as a function of $R$ we get

$$0 = \sum_{n=0}^{\infty} R^{2n+k-k_0} (1 - R^{2k+k_0-1}) \cdot \int_{UX} a_n(x) b_{n+k-k_0}(x) \cdot d\mu(x).$$  

Thus,

$$0 = \sum_{n=0}^{\infty} R^{2n} \cdot \int_{UX} a_n(x) b_{n+k-k_0}(x) \cdot d\mu(x),$$

and this last assertion concludes the proof of the corollary. □
6.4.5. The variance and the analytic expansion. We perform a similar analysis for the variance to get

**Proposition 6.4.6.** We have

\[ \text{Var}(q) = \lim_{R \to 1} \frac{1}{\log(1 - R)} \sum_{n=0}^{\infty} A_n I_{n,k-1}^2(R) . \]

**Proof.** Recall (61). Then by the invariance of the Liouville measure,

\[
\text{Var}(q) = \lim_{r \to \infty} \frac{1}{r} \int_{U_X} \left( \int_0^r \hat{q}(\phi_s(x)) \, dt \right)^2 \, d\mu(x)
\]

\[
= \lim_{r \to \infty} \frac{1}{r} \int_{U_X} \frac{1}{2\pi} \int_0^{2\pi} \left( \int_0^r \hat{q}(\phi_s(e^{i\theta}x)) \, dt \right)^2 \, d\theta \, d\mu(x)
\]

\[
= \lim_{R \to 1} \frac{2}{\log(1 + R)} \int_{U_X} \int_0^{2\pi} \left( \int_0^R \hat{q}(Te^{i\theta}) \, dT \right)^2 \, d\theta \, d\mu(x), \quad \text{by (70)}
\]

\[
= \lim_{R \to 1} \frac{1}{\log(1 + R)} \sum_{n=0}^{\infty} a_n \mu I_{n,k-1}^2(R) \, d\mu(x), \quad \text{by (64)}
\]

\[
= \lim_{R \to 1} \frac{1}{\log(1 - R)} \sum_{n=0}^{\infty} A_n I_{n,k-1}^2(R), \quad \text{by (72)}
\]

where in the last line we have used the fact that as \( R \to 1 \),

\[
\log \left( \frac{1 + R}{1 - R} \right) \sim |\log(1 - R)|.
\]

\[ \square \]

6.4.6. The bilinear part of the variance and the asymptotic expansion. We have

**Proposition 6.4.7.** If \( q \) and \( q^0 \) are holomorphic forms with different degrees, then

\[
\lim_{t \to 0} \frac{1}{t} \int_{U_X} \left( \int_0^t \hat{q}(\phi_s(x)) \, ds \right) \left( \int_0^t \hat{q}^0(\phi_s(x)) \, ds \right) \, d\mu(x) = 0 .
\]

**Proof.** We use the invariance by the action of the circle to get that

\[
B := \int_{U_X} \left( \int_0^t \hat{q}(\phi_s(x)) \, ds \right) \left( \int_0^t \hat{q}^0(\phi_s(x)) \, ds \right) \, d\mu(x)
\]

\[
= \frac{1}{2\pi} \int_{U_X} \int_0^{2\pi} \left( \int_0^t \hat{q}(\phi_s(e^{i\theta}x)) \, ds \right) \left( \int_0^t \hat{q}^0(\phi_s(e^{i\theta}x)) \, ds \right) \, d\mu(x) \, d\theta .
\]

Using the identification (70), and \( s = r(S) \), we get that

\[
B = \frac{1}{2\pi} \int_{U_X} \int_0^{2\pi} \left( \int_0^T \hat{q}(S e^{i\theta}) \, dS \, \frac{dS}{1 - S^2} \right) \left( \int_0^T \hat{q}^0(S e^{i\theta}) \, dS \, \frac{dS}{1 - S^2} \right) \, d\mu(x) \, d\theta .
\]

We now use (65) to get that

\[
B = \frac{1}{2} \sum_{n=0}^{\infty} \Re \left( \int_{U_X} b_{n+k-1}(x) \bar{s}_n(x) \, d\mu(x) \right) I_{n,k-1}(T) I_{n+k-1}(T).
\]

But from (76), \( B = 0 \). This concludes the proof of the result. \[ \square \]
6.5. **Proof of Theorem 6.3.1.** From Proposition 6.4.6, we get that

\[ \text{Var}(q) = \lim_{R \to 1} \frac{1}{|\log(1 - R)|} \sum_{n=0}^{\infty} A_n I_{n,k-1}(R). \]

By Corollary 6.4.5,

\[ \text{Var}(q) = \frac{2^{-k}}{2\pi|\chi(X)|} \|q\|_X^2 \left( \lim_{R \to 1} \frac{1}{|\log(1 - R)|} \sum_{n=0}^{\infty} \left( \frac{n + 2k - 1}{n} \right) I_{n,k-1}(R) \right). \]

Thus the first part of the theorem follows from Theorem 6.4.1. The second part follows from Proposition 6.4.7.

6.6. **A family of pressure metrics.** The pressure metric (which may sometimes be degenerate) was defined in [8] for projective Anosov representations using the logarithm of the highest eigenvalue as an initial input for “length”. The spectrum of a representation is defined as the indexed collection of lengths. Then, from this data, the intersection and pressure metric are defined. In the context of Hitchin representations, one has many choices for lengths. For instance, one could take the quotient of two consecutive eigenvalues, or any polynomial in these. Via the Gardiner formula that we have obtained, we can in principle compute all the associated pressure metrics at the Fuchsian locus. Jörgen Andersen and Scott Wolpert have suggested to us that for one of these lengths the pressure metric should actually be the \( L^2 \)-metric. A count of parameters is consistent with this idea. Starting from the work presented here, a somewhat tedious calculation might verify this conjecture.

7. **First Variation of Cross Ratios**

We recall that one can associate a cross ratio to a Hitchin representation: that is, a function on 4 points on the boundary at infinity (see [26]). More precisely a cross ratio is a function \( b \) on

\[ \partial_{\infty} \pi_1(\Sigma)^4 := \{(x, y, z, w) \in \partial_{\infty} \pi_1(\Sigma)^4 \mid x \neq y \neq z \neq w \}. \]

satisfying some functional rules. Conversely, the cross ratio determines the representation. Moreover, one can characterize those cross ratios which arise from Hitchin representations [26, Theorem 1.1]. For example, the cross ratio for a Fuchsian representation is the cross ratio that comes from the identification of \( \partial_{\infty} \pi_1(\Sigma) \) with \( \mathbb{P}^1(R) \) associated to the hyperbolic structure.

More generally, a cross ratio is associated to projective Anosov representations. Since the construction only depends on the limit curve of the Anosov representation, the value of the cross ratio at four given points \( (x, y, z, w) \) depends analytically on the representation (cf. [8, Theorem 6.1]).

The purpose of this section is to describe the variation of the cross ratio under oper and Hitchin deformations associated to a holomorphic differential of degree \( k \). This is achieved in two theorems: Theorem 7.2.1 and Theorem 7.2.2 stated in Section 7.2. The first theorem interprets the variation as a “generalized period” which we call a **rhombus function**. This requires further definition. The second theorem gives a description in terms of automorphic forms, and we can state it right now. This result is a generalization of [41, Lemma 1.1]. The latter uses techniques of quasiconformal maps which are not available in this context.
Theorem 7.0.1. Let $\delta$ be a family of Hitchin representations coming from the standard deformation associated to a $\pi_1(\Sigma)$-invariant holomorphic differential $q$ of degree $k$. Let $b_{ij}^t$ be the corresponding family of cross ratios. For $(x, y, X, Y) \in \partial_\infty \pi_1(\Sigma)^4$, 

\[
\frac{d}{dt} \mid_{t=0} \log b_{ij}^t(x, y, X, Y) = \left( \frac{(-1)^k (n-1)!}{2^{k-2} (n-k)!} \right) \cdot r_k \cdot \int_{\mathbb{H}^2} \mathbb{R} \, q(z) \, \sum_{i,j \in \{1,2\}} (-1)^{i+j} \left( \theta_{x_i,x_j} \right)^k \, d\sigma(z),
\]
where $r_k$ is defined in (7).

After a preliminary definition of the rhombus function, we state our two main results. Then we proceed through the proof. It is enough to study oper variations, since Hitchin variations are twice the "real part" of oper variations. The idea is to first study the variation of cross ratios for points inside $\mathbb{H}^2$ mapped into $\mathbb{C}P(E)$ and $\mathbb{C}P(E^*)$ by the associated Frenet immersion (see Section 7). Once this is done we extend this variational formula to points in the boundary at infinity in Section 7.4. Finally in Section 7.4 we relate the rhombus function to automorphic forms.

7.1. Preliminary: the rhombus function. We recall the following obvious and classical lemma which follows from the fact that geodesics ending at the same point at infinity in $\mathbb{H}^2$ approach each other exponentially fast.

Lemma 7.1.1. Let $f$ be a Hölder function on $\mathbb{H}^2$. Let $\gamma_1$ and $\gamma_2$ be two geodesics in $\mathbb{H}^2$ so that $\gamma_1(+\infty) = \gamma_2(+\infty)$. Then 

\[
t \to \int_0^t f(\gamma_1(s)) \, ds - \int_0^t f(\gamma_2(s)) \, ds,
\]

admits a limit when $t$ goes to $+\infty$.

The previous lemma allows us to make the following definition. Given four points $x, y, X, Y \in \partial_\infty \mathbb{H}^2$, let $h_x, h_y, h_X, h_Y$ be corresponding Busemann functions. Let $\gamma_{z,Z}$ be the geodesic going from $z$ to $Z$. For any $a \in \{x, y\}, A \in \{X, Y\}$ and $c \in [a, A]$, choose $r_{a,A}$ so that 

\[
h_c(\gamma_{a,A}(r_{a,A})) = t.
\]

For a Hölder function $f$, we define the rhombus function by 

\[
\text{Rh}(x, y, X, Y; f) := \lim_{t \to \infty} \left\{ \int_{r_{x,X}}^{r_{x,X}} f(\gamma_{x,X}) \, ds - \int_{r_{x,X}}^{r_{x,X}} f(\gamma_{y,X}) \, ds - \int_{r_{x,Y}}^{r_{x,Y}} f(\gamma_{x,Y}) \, ds + \int_{r_{x,Y}}^{r_{x,Y}} f(\gamma_{y,Y}) \, ds \right\}
\]

Intuitively speaking, $\text{Rh}(x, y, X, Y; f)$ is the alternating sum of integrals along geodesics that pairwise meet at points at infinity. See Figure 1.

The following proposition is a consequence of hyperbolicity.

Proposition 7.1.2. Let $f$ be a Hölder function on $\mathbb{H}^2$. Let $x_1, x_2, X_1, X_2$ be pairwise distinct points in $\partial_\infty \mathbb{H}^2$. Let $(x^n_{i,j})_{n \in \mathbb{N}}$ and $(X^n_{i,j})_{n \in \mathbb{N}}$ be points in $\mathbb{H}^2$, $i, j \in \{1, 2\}$. Assume that 

1. $\lim_{n \to \infty} (x^n_{i,j}) = x_i$. 
2. $\lim_{n \to \infty} (X^n_{i,j}) = X_j$. 
3. $\lim_{n \to \infty} d(x^n_{i,j}, x^n_{j,i}) = 0$. 

Then $\text{Rh}(x_1, x_2, X_1, X_2; f)$ converges to a limit as $t \to +\infty$. 

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\[ \lim_{n \to \infty} d(X^n_{i,j}, X^n'_{i,j}) = 0. \]

Then
\[ \text{Rh}(x_1, x_2, X_1, X_2; f) = \lim_{n \to \infty} \sum_{i,j \in \{1,2\}} (-1)^{i+j} \int_{X^n_{i,j}} f \left( \eta^n_{i,j} \right) \, ds, \]

where we consider the geodesics \( \eta^n_{i,j} \) joining \( x^n_{i,j} \) to \( X^n_{i,j} \) as curves in \( \mathbf{U} \mathbf{H}^2 \).

**Proof.** We may as well assume that
\[ h_{x_i} (x^n_{i,j}) = h_{x_i} (x^n_{i,k}) =: t^n_i, \]
and
\[ h_{X_i} (X^n_{i,j}) = h_{X_i} (X^n_{i,j}) =: T^n_j. \]

Let then \( z^n_{i,j} \) and \( Z^n_{i,j} \) be the points in the geodesic \( \gamma_{i,j} := \gamma_{x_i, X_i} \) so that
\[ h_{X_i} (Z^n_{i,j}) =: T^n_j, \quad h_{x_i} (z^n_{i,j}) =: t^n_i. \]

Elementary hyperbolic geometry implies that the sequences \( \{d(z^n_{i,j}, x^n_{i,j})\}_{n \in \mathbb{N}} \) and \( \{d(Z^n_{i,j}, X^n_{i,j})\}_{n \in \mathbb{N}} \) all converge to zero. For a positive number \( \alpha \), there exists a positive \( \epsilon \) such that if \( d(x, z) \leq \epsilon \) and \( d(X, Z) \leq \epsilon \), then
\[ \left| \int_x^X f (\gamma_{x,X}) \, ds - \int_z^Z f (\gamma_{z,Z}) \, ds \right| \leq \alpha. \]

It follows that
\[ \lim_{n \to \infty} \left| \int_{x^n_{i,j}}^{X^n_{i,j}} f (\eta^n_{i,j}) \, ds - \int_{z^n_{i,j}}^{Z^n_{i,j}} f (\gamma^n_{i,j}) \, ds \right| = 0. \]

The proposition follows from this last assertion. \( \square \)

**7.2. Statement of the results.** We can now state our two main results on the rhombus function.
7.2.1. The period formulation for variations of cross ratio. The following will be proved in Section 7.4.

Theorem 7.2.1. Let \((x, y, X, Y) \in \partial_\infty \pi_1(\Sigma)^*\). Let \(\delta_0\) be a Fuchsian representation giving rise to a hyperbolic structure on \(\Sigma\), and let \(\delta_t\) be a smooth family of representations of \(\pi_1(\Sigma) \to \text{SL}(2, \mathbb{C})\) so that the tangent vector at \(\delta_0\) is a standard oper variation given by a holomorphic \(k\)-differential \(q\). Let \(b^0_t\) be the corresponding family of cross ratios on \(\partial_\infty \pi_1(\Sigma)\). Then

\[
\frac{d}{dt} \bigg|_{t=0} \log b^0_t(x, y, X, Y) = \left(\frac{(-1)^k (n-1)!}{2^k 1^1(n-k)!}\right) \cdot \text{Rh}(x, y, X, Y; \dot{q}) ,
\]

where \(\dot{q}\) is the function defined in Section 2.1.2.

Similarly, assume that the tangent vector at \(\delta_0\) is a standard Hitchin variation given by a holomorphic \(k\) differential \(q\). Let \(b^1_t\) be the corresponding family of cross ratios on \(\partial_\infty \pi_1(\Sigma)\). Then

\[
\frac{d}{dt} \bigg|_{t=0} \log b^1_t(x, y, X, Y) = \left(\frac{(-1)^k (n-1)!}{2^k 2^2(n-k)!}\right) \cdot \Re \left(\text{Rh}(x, y, X, Y; \dot{q})\right) .
\]

7.2.2. The automorphic form formulation. An alternative formula for the rhombus function, and hence for the variation of cross ratios, is given in the following result. We prove this theorem in Section 7.5.

Theorem 7.2.2. Let \(q\) be a \(\pi_1(\Sigma)\)-automorphic holomorphic \(k\)-differential on \(\mathbb{H}^2\), and let \(x, y, X, Y \in \mathbb{H}^2\) be pairwise distinct points. Then

\[
\text{Rh}(x_1, x_2, X_1, X_2; \dot{q}) = r_k \cdot \int_{\mathbb{H}^2} \left( q(z), \sum_{i,j \in \{1,2\}} (-1)^{i+j} \left( \frac{(X_j - x_i) \, dz}{(z-x_i)(z-X_j)} \right)^k \right) \, d\sigma(z) ,
\]

where \(r_k\) is defined in (7).

7.3. Oper, Frenet immersions and cross ratios. Let \(D\) be a nonslipping connection on \(P\) equipped with the flag structure \(\mathcal{F}_1 \subset \mathcal{F}_2 \subset \cdots \subset \mathcal{F}_n = P\). Let \(D\) be the universal cover of \(X\). Let \(\delta\) be the holonomy of this connection. Then \(\mathcal{F}_1\), viewed as a line subbundle over \(D\), defines a \(\delta\)-equivariant holomorphic map \(D \to \text{CP}(\mathcal{P})\). Moreover, this map is a Frenet immersion: for every \(k\), the derivatives up to order \(k\) generate a \(k\)-plane (called the \(k\)th-osculating plane of \(\mathcal{F}\)) which is actually identified to the projective subspace generated by \(\mathcal{F}_1\). Conversely, any such Frenet immersion uniquely defines an oper, and we have thus described an isomorphism between opers up to gauge equivalence, and Frenet immersions up to precomposition by projective transformations.

To an oper we can also associate a cross ratio, which in this case is a function of four generic points of \(D\). The construction goes as follows: let \(x, y, X, Y\) be four points on \(D\). We trivialise \(P\) over \(D\) as \(D \times E\) and consider \(\mathcal{F}_1 =: \xi\) and \(\mathcal{F}_n =: \xi^*\) as maps from \(D\) to \(\text{CP}(E)\) and \(\text{CP}(E^*)\), respectively. Then by definition the cross ratio of four generic points is

\[
b^0(x, y, X, Y) = \frac{\langle \hat{\xi}(x) | \hat{\xi}^*(X) \rangle \langle \hat{\xi}(y) | \hat{\xi}^*(Y) \rangle}{\langle \hat{\xi}(x) | \hat{\xi}^*(Y) \rangle \langle \hat{\xi}(y) | \hat{\xi}^*(X) \rangle} ,
\]

where \(\hat{\xi}(z), \hat{\xi}^*(z)\) is a nonzero vector in \(\xi(z)\) and \(\xi^*(z)\), respectively. The generic condition is that \(\mathcal{F}_1(x) \subset \mathcal{F}^0_{n-1}(Y)\) and \(\mathcal{F}_1(y) \subset \mathcal{F}^0_{n-1}(X)\). For the Veronese oper, this
condition is always satisfied provided $x \neq Y$ and $y \neq X$. Thus given $x, y, X, Y$, there exists an open set of opers for which the cross ratio is defined at $x, y, X, Y$.

It then makes sense to compute the variation of $b^0(x, y, X, Y)$ along an oper deformation. The main result of this section is an explicit formula for this variation (for points inside the disk $D$) at the Veronese oper. More precisely, we shall prove

**Proposition 7.3.1.** Let $\{x_1^n, x_2^n, X_1^n, X_2^n\}_{n \in \mathbb{N}}$ be a sequence of quadruples in $D$ converging to $x_1, x_2, X_1, X_2$ in $\partial D$. Then along an oper deformation associated to a holomorphic $k$-differential $q$, \[ \lim_{n \to \infty} \frac{b^0(x_1^n, x_2^n, X_1^n, X_2^n)}{b^0(x_1^n, x_2^n, X_1^n, X_2^n)} = A(k, n) \cdot \text{Rh}(x_1, x_2, X_1, X_2; \hat{q}) , \] where $A(k, n) = \frac{(-1)^k (n-k)!}{n!}$.

We will prove this proposition in Section 7.3.3.

7.3.1. *Lifting and elementary functions.* Let us choose points $z$ and $Z$ in $D$ and a path $\gamma : [0, 1] \to D$ so that $\gamma(0) = z$ and $\gamma(1) = Z$. Let $\zeta_0 \in \xi(z)$ and $\zeta_1^* \in \xi^*(Z)$ be nonzero vectors and covectors. Consider the elementary function $G_{z, Z}$ on the space of connections on $E$ defined by \[ G_{z, Z}(V) = \langle \zeta_1^*, \zeta_1(V) \rangle , \] where $t \to \zeta_t(V)$ is the $V$-parallel section along $\gamma$ starting at $\zeta_0$. Let also $t \to \zeta_t^*(V)$ be the $V$-parallel section so that $\zeta_t^*(V) = \zeta_t^*$. From the definition of the cross ratio we immediately get

**Proposition 7.3.2.** For $x_1, x_2, X_1, X_2$ in $D$, and an oper variation associated to a $k$-holomorphic differential $q$, one has \[ \lim_{n \to \infty} \frac{\partial^0(x_1^n, x_2^n, X_1^n, X_2^n)}{\partial^0(x_1^n, x_2^n, X_1^n, X_2^n)} = \sum_{(i, j) \in \{1, 2\}} (-1)^{j+i} \frac{dV G_{z, X}(V)}{G_{z, X}(V)} , \] \[ \text{Observe that, fixing } z \text{ and } Z, G_{z, Z} \text{ is well defined up to a multiplicative constant. Then let } \] \[ p_{z, Z}(s) = \frac{\zeta_1^* \otimes \zeta_0(V)}{\langle \zeta_1^* \mid \zeta_0(V) \rangle} , \] be the projection on the $V$-parallel line $L_z$ generated by $\zeta_0$, along the $V$-parallel hyperplane $P_Z = \ker(\zeta_0)$. Then the variation of parameters method gives

**Proposition 7.3.3.** We have \[ \frac{dV G_{z, Z}(\hat{V})}{G_{z, Z}(V)} = \int_0^1 \text{Tr} \left( p_{z, Z}(\gamma) \cdot \hat{V} \right) \, ds . \] **Proof.** Let $V_\delta$ be a 1-parameter family of connections so that $\frac{d}{ds}_{s=0} V_\delta = \hat{V}$. Let $\beta$ be the section of $\text{End}(P)$ along $\gamma$ so that $\beta(\gamma(0)) = 0$ and $\hat{V} \cdot \gamma(0) = \hat{V} \cdot \gamma(1)$. Let $G^\delta$ be the family of sections of $\text{End}(P)$ along $\gamma$ so that $G^\delta(z) = \text{Id}$ and $(G^\delta)^* V = V^\delta$. Then by construction \[ \frac{d}{ds}_{s=0} G^\delta = \beta . \]
Moreover,

\[ G_{z,L}(\mathcal{V}^*) = \langle \zeta_1^*(\mathcal{V}) | G \zeta_1(\mathcal{V}) \rangle . \]

It follows that

\[ \text{d}_V G_{z,L}(\mathcal{V}) = \langle \zeta_1^*(\mathcal{V}) | \beta (\zeta_1(\mathcal{V})) \rangle . \]

Let now

\[ f(s) = \langle \zeta_1^*(\mathcal{V}) | \beta (\zeta_1(\mathcal{V})) \rangle , \]

so that

\[
\text{d}_V G_{z,L}(\mathcal{V}) = f(1) - f(0) = \int_0^1 \frac{df}{ds} ds = \int_0^1 \langle \zeta_1^*(\mathcal{V}) | \beta (\zeta_1(\mathcal{V})) \rangle ds .
\]

Since \( G_{z,L}(\mathcal{V}) = \text{Tr}(\zeta_1^*(\mathcal{V}) \otimes \zeta_1(\mathcal{V})) \), this completes the proof of the proposition. \( \Box \)

7.3.2. The Fuchsian bundle along a geodesic again. Let \( \gamma \) be a geodesic, we need the following

Lemma 7.3.4. There exist positive number \( \alpha \) and \( t_0 \) only depending on \( n \), so that for any \( k \in \{2, \ldots, n\} \), for any \( s, t > t_0 \)

\[
|\text{Tr} (E^0_{k-1} \cdot p_{\gamma(-s),\gamma(u)}) - A(k,n)| \leq e^{-as} + e^{-as} .
\]

Proof. We use the notation of Section 4.2. Let \( \gamma \) be a geodesic and let \( s \) consider the bundle \( \mathcal{E} \) in (25). As noted previously, this bundle inherits a (split) connection \( \mathcal{V} \) from the hyperbolic structure, as well as a (nonsplit) connection \( D \) from the flat connection. From the description of the Veronese oper in Section 2.4.1 we have \( \xi(z) = S_{n+1} \) and \( \xi^*(z) = S_{n+1}^* \otimes \cdots \otimes S_{n+1}^* \). Finally we have a \( \mathcal{V} \)-parallel frame \( \{u_p\} \)

defined in eq. (48) and which satisfies by eq. (49)

\[ D_{\gamma} u_p = \alpha_p u_p , \]

where \( \alpha_n > \cdots \alpha_2 > \alpha_1 > 0 \). Then \( L_p \) is the (parallel) complex line generated by \( u_p \) and let \( \pi_p \) be the orthogonal projection on \( L_p \) for all \( p \).

By construction, \( \pi_1(\xi(z)) \) is always nonzero. Then, since \( \xi(z) \) is \( \mathcal{V} \) parallel, it follows that there exist positive \( \alpha, K \) and \( t_0 \) so that for \( s > t_0 \)

\[ d(\xi(\gamma(s)), L_1) \leq K e^{-as} , \]

where \( d \) is the metric on \( \mathbb{C}P(\mathcal{E}) \) inherited from the metric on \( \mathcal{E} \). It then follows (using a similar argument on \( E \)) that we have a constant \( K \) so that for \( s, u > t_0 \)

\[ d(\pi_1, p_{\gamma(-s),\gamma(u)}) \leq K(e^{-as} + e^{-as}) . \]

The result now follows from eq. (97). \( \Box \)

As a corollary of Proposition 7.3.3 and Lemma 7.3.4 we get

Corollary 7.3.5. There exist positive number \( \alpha \) and \( t_0 \) only depending on \( n \), so that for any \( k \in \{2, \ldots, n\} \), for any \( s, u > t_0 \)

\[
\left| \frac{\text{d}_V G_{\gamma(-s),\gamma(u)}(\mathcal{V})}{G_{\gamma(-s),\gamma(u)}(\mathcal{V})} - A(k,n) \cdot \int_s^u q(\gamma_1, \ldots, \gamma_k) ds \right| \leq e^{-as} + e^{-as} .
\]
The corollary now follows.

The second assertion is an immediate consequence of the first one using the variation of cross ratios: Theorem 7.2.1. The statement consists of two assertions. The result now follows from Proposition 7.1.2.

7.3.3. Proof of Proposition 7.3.1. Let \([x_1^n, x_2^n, X_1^n, X_2^n]_{n \in \mathbb{N}}\) be a sequence of quadruples in \(\mathbb{D}\) converging to \(x_1, x_2, X_1, X_2 \in \partial \mathbb{D}\). Let \(\gamma_{ij}^n\) be the geodesics joining \(x_i^n\) to \(X_j^n\). Since \(x_1, x_2, X_1, X_2\) are all pairwise distinct, we may choose a parametrization so that \(\gamma_{ij}^n(0)\) stays in a compact region. Thus let \(\gamma_{ij}^n(s_{ij}^n) = x_i\) and \(\gamma_{ij}^n(s_{ij}^n) = X_j\).

From Corollary 7.3.5 and Proposition 7.3.2, it follows that

\[
\left| b^O(x_1^n, x_2^n, X_1^n, X_2^n) - A(k, n) \left( \sum_{i,j \in [1,2]} (-1)^{i+j} \int_{\gamma_{ij}^n} \frac{ds}{g_{ij}} \right) \right| \leq \sum_{i,j} e^{-\alpha s_{ij}^n} + e^{-\alpha s_{ij}^n}.
\]

The result now follows from Proposition 7.1.2.

7.4. Extension to the boundary at infinity. We can now prove our first result on variation of cross ratios: Theorem 7.2.1. The statement consists of two assertions. The second assertion is an immediate consequence of the first one using the description of the tangent space to the Fuchsian point of Section 3.1.

The proof of the theorem therefore follows from Proposition 7.3.1 and the following result, whose proof occupies the rest of this section.

**Proposition 7.4.1.** Assume that the function defined for every quadruple of pairwise distinct points in \(\mathbb{D}\) by

\[
(x, y, X, Y) \mapsto \frac{d}{dt}igr|_{t=0} b^O(x, y, X, Y),
\]

extends continuously to a function \(f\) defined on quadruple of pairwise distinct points of \(\partial \mathbb{D}\). Then, for every \((x, y, X, Y)\) pairwise distinct in \(\partial \mathbb{D}\),

\[
f(x, y, X, Y) = \frac{d}{dt}igr|_{t=0} b^O(x, y, X, Y).
\]

7.4.1. Operas and the boundary at infinity. For the moment we prove the following result. Let us first notice that, by the openness property of Anosov representations there exists a neighborhood \(U_0\) of the Veronese oper, so that if an oper \(O\) belongs to \(U_0\) its monodromy is also Anosov. In particular, one can associate to \(O\) a cross ratio on \(\partial \mathbb{D}\) (identified with \(\partial_\infty \pi_1(\Sigma)\)), which will also be denoted by \(b^O\).

Although it is quite likely that on a smaller neighborhood, the cross ratio is actually defined for all quadruple of distinct points and extends continuously to the boundary at infinity, we content ourselves with a weaker and easier result.

In this section and the next one, \(\xi_0, \xi^*_0\) are the Frenet immersions from \(\mathbb{D}\) to \(\mathbb{CP}(E)\) and \(\mathbb{CP}(E^*)\), respectively, associated to the oper \(O\), and \(\xi, \xi^*\) those associated to the Veronese oper as in Section 7.3. Let also denote by \(\delta_0\) (and similarly \(\delta\)) the monodromy of the oper \(O\) and the Veronese oper, respectively. Finally the contragredient of a representation \(\eta\) is denoted \(\eta^\cdot\).
Proposition 7.4.2. Let \( x \) be a point in \( \mathcal{D} \), let \( \alpha_1, \alpha_2, \alpha_3 \) and \( \alpha_4 \) be four nontrivial elements of \( \pi_1(\Sigma) \). Let \( \alpha^+_1 \) be the attracting point of \( \alpha_1 \) in \( \partial \mathcal{D} \). Assume that \( \alpha^+_1, \alpha^+_2, \alpha^+_3 \) and \( \alpha^+_4 \) are all pairwise distinct. Then there exists \( n_0 \) and a neighborhood \( U \) of the Veronese oper, such that if \( O \in U \), then

- First, for \( n > n_0 \), \( b^O(\alpha^+_1(x), \alpha^+_2(x), \alpha^+_3(x), \alpha^+_4(x)) \) is well defined,
- Moreover, \( \lim_{n \to \infty} b^O(\alpha^+_1(x), \alpha^+_2(x), \alpha^+_3(x), \alpha^+_4(x)) = b^O(\alpha^+_1, \alpha^+_2, \alpha^+_3, \alpha^+_4) \), and the limit is uniform in \( O \).

Proof. Observe that \( \xi(x) \) and \( \xi^*(x) \) both lie in the basin of attraction of the unique fixed point of \( \delta(\alpha_i) \) in either \( \mathbb{C}P(E) \) or \( \mathbb{C}P(E') \). By continuity the same holds for \( \xi_O \) an \( \xi^*_O \) for an oper \( O \) close enough to the Veronese oper. It then follows that \( \xi_O(\alpha^+_i(x)) \) converges to the attracting fixed point of \( \delta_O(\alpha_i) \), and similarly \( \xi^*_O(\alpha^+_i(x)) \) converges to the attracting fixed point of \( \delta^*_O(\alpha_i) \). The results follow immediately, and the uniform convergence in \( O \) is similarly obtained.

7.4.2. Analyticity. We note the following

Proposition 7.4.3. Let \( x, y, X, Y \) in \( \mathcal{D} \) (or all in \( \partial \mathcal{D} \)). Then the function \( O \mapsto b_O(x, y, X, Y) \) is complex analytic in a neighborhood of the Veronese oper.

Proof. For points in \( \partial \mathcal{D} \), this follows from [8, Theorem 6.1]. For points in \( \mathcal{D} \) this follows from the fact that for a differential equation whose parameters vary complex analytically, the solution also varies complex analytically.

7.4.3. Proof of Proposition 7.4.1. We remark that by continuity it is enough to prove that on a dense set of points in \( \partial \mathcal{D} \),

\[
\phi(x, y, X, Y) = b^O(x, y, X, Y).
\]

We thus will only prove this proposition when \( x, y, X, Y \) are the attracting points of elements \( \alpha_i \) in \( \pi_1(\Sigma) \). But by Proposition 7.4.2, there exists sequences \( \{x_n\}_{n \in \mathbb{N}}, \{y_n\}_{n \in \mathbb{N}} \), \( \{X_n\}_{n \in \mathbb{N}} \) and \( \{Y_n\}_{n \in \mathbb{N}} \) converging respectively to \( x, y, X \) and \( Y \) so that the function

\[
O \mapsto b^O(x_n, y_n, X_n, Y_n),
\]

converges uniformly to \( O \mapsto b^O(x, y, X, Y) \). Since all the functions involved are analytic in \( O \) by Proposition 7.4.3, it follows that their derivatives in \( O \) also converges. This is what we needed to prove, and the proof of Proposition 7.4.1 is now complete.

7.5. The rhombus function and automorphic forms. In this section, we prove Theorem 7.2.2 to give an alternative formula for the rhombus function, and hence also for the variation of cross ratios.

7.5.1. Slabs and cusps. For \( u \) and \( v \) in \( \mathbb{H}^2 \), let \( c \) be the geodesic arc between \( u \) and \( v \). The slab \( S(u,v) \) defined by \( u \) and \( v \) is the region bounded by the two orthogonal geodesics \( c_x \) and \( c_y \) to \( c \) at \( u \) and \( v \), respectively. We extend the notion of slabs to \( u \) and being possibly at infinity.

We need the following two propositions, proved in the next two sections, and their corollaries. Our first proposition and corollary deal with integration over slabs. Recall (4).
Proposition 7.5.1. \textbf{[Integration over slabs]} There exists a constant $K$, such that such that for any bounded $k$-differential $q$, and $u, v \in \mathbb{H}^2$, we have
\[
\left| \int_{S(u, v)} \langle q, \Theta^k_{u,v} \rangle \, d\sigma \right| \leq K \cdot \|q\|_{\infty} \cdot d(u, v) .
\]

\textbf{Proof.} We can work in the upper half plane model and assume that $u$ and $v$ belong to the imaginary axis so that
\[
\Theta_{u,v} = \frac{dz}{z}.
\]
Then the slab $S(u, v)$ is defined in polar coordinates as
\[
S(u, v) = \{(r, \theta) | \theta \in [-\frac{\pi}{2}, \frac{\pi}{2}], |u| \leq r |v| \}.
\]
Then
\[
d\sigma = \frac{1}{r \sin(\theta)^2} \, dr \, d\theta, \quad \|\Theta^k_{u,v}\| = \sin^2(\theta).
\]
Thus the result follows from a simple integration with $K = \int_{-\frac{\pi}{2}}^{\frac{\pi}{2}} \sin^{k-2}(\theta) \, d\theta$. \hfill \Box

It will be useful to note the following.

\textbf{Corollary 7.5.2.} Let $q$ be a bounded $k$-differential. Then there exist constants $K$ and $\varepsilon_0$ such that if $\varepsilon < \varepsilon_0$ and $d(u_0, v_0) + d(u_1, v_1) \leq \varepsilon$ and $d(u_0, u_1) \geq K$ then
\[
\left| \int_{S(u_0, u_1)} \langle q, \Theta^k_{u_0, u_1} \rangle \, d\sigma - \int_{S(v_0, v_1)} \langle q, \Theta^k_{v_0, v_1} \rangle \, d\sigma \right| \leq K \cdot \varepsilon .
\]

\textbf{Proof.} Let $\gamma$ be the geodesic passing through $u_0$ and $u_1$ so that $\gamma(0) = u_0$ and $\gamma(\varepsilon) = u_1$. Then, by continuity and elementary hyperbolic geometry, there exist constants $A$ and $\varepsilon_0$ such that if $\varepsilon < \varepsilon_0$ and $d(u_0, v_0) + d(u_1, v_1) \leq \varepsilon$ and $d(u_0, u_1) \geq K$ then, writing $u_0^+ = \gamma(\pm A \varepsilon)$ and $u_1^+ = \gamma(A \varepsilon, \pm \varepsilon))$, we have
\[
S(u_0^+, u_1^+) \subset S(v_0, v_1) \subset S(u_0^+, u_1^+) .
\]
Now let $\chi$ be the characteristic function of $S(v_0, v_1)$. Then using the previous proposition twice yields
\[
\left| \int_{S(u_0, u_1)} \langle q, \Theta^k_{u_0, u_1} \rangle \, d\sigma - \int_{S(v_0, v_1)} \langle q, \Theta^k_{v_0, v_1} \rangle \, d\sigma \right| \leq 2AK \cdot \varepsilon \|q\|_{\infty} + \left| \int_{S(u_0^+, u_1^+)} \langle q, \Theta^k_{u_0, u_1} \rangle \, d\sigma - \int_{S(v_0, v_1)} \langle q, \Theta^k_{v_0, v_1} \rangle \, d\sigma \right| \leq 2AK \cdot \varepsilon \|q\|_{\infty} + \int_{S(u_0^+, u_1^+)} \langle (1 - \chi)q, \Theta^k_{u_0, u_1} \rangle \, d\sigma \leq 2AK \cdot \varepsilon \|q\|_{\infty} + \int_{S(u_0^+, u_1^+)} \langle (1 - \chi)q, \Theta^k_{u_0, u_1} \rangle \, d\sigma \leq 6AK \cdot \varepsilon \|q\|_{\infty} .
\]
The result follows from this last inequality. \hfill \Box
Proposition 7.5.3. Let \( q \) be a Lipschitz differential. Let \( \gamma \) be a geodesic from \( x \) to \( X \) in \( \partial_\infty \mathbb{H}^2 \). Let \( \tau \) be a parabolic transformation fixing \( X \) then
\[
\lim_{t \to \infty} \int_{S(\gamma(t),X)} \langle q - \tau^* q, \theta_k^{\gamma,\gamma^+} \rangle = 0.
\]

Proof. We use the upper half plane model and the geodesic joining 0 to \( \infty \). Let also \( \tau : z \mapsto \frac{az}{z + \alpha} \).

Let \( \chi_t \) be the indicatrix of the set \( S(\gamma(t),X) \). We want to prove that
\[
\langle q - \tau^* q, \frac{dz}{z} \rangle \cdot \chi_t \overset{L^1}{\longrightarrow} 0,
\]
where the convergence in \( L^1(\mathbb{H}^2,\sigma) \). Observe that since \( q \) is Lipschitz, there exists a constant \( K \) so that
\[
\|q - \tau^* q\|_\alpha \leq K \cdot d(z, \tau(z)).
\]
Moreover, there exists a constant \( K_2 \), so that \( d_H(z, \tau(z)) \leq K_2 J(z) \). It follows that we have a constant \( K_3 \) so that, using polar coordinates,
\[
\left| \langle q - \tau^* q, \frac{dz}{z} \rangle \right| \leq K_3 J(z)^{k+1} = K_3 \cdot r \sin^{k+1}(\theta).
\]
Integrating along the area form of the hyperbolic space
\[
d\sigma = \frac{drd\theta}{r \sin(\theta)^2},
\]
we get
\[
\left| \langle q - \tau^* q, \frac{dz}{z} \rangle \right| d\sigma \leq K_3 \sin^{k-1}(\theta) drd\theta.
\]
The result follows. \( \square \)

As a corollary of Proposition 7.5.3, we have

Corollary 7.5.4. Let \( x, X_1, X_2 \in \partial_\infty \mathbb{H}^2 \). Let \( h \) be a Busemann function associated to \( x \), so that \( h(-\infty) = x \). Let \( \gamma_i \) be geodesics from \( x \) to \( X_i \) for \( i \in \{1,2\} \). We parametrise the geodesics so that \( h \circ \gamma_i(t) = t \). Then for any sequence \( \{s_n\}_{n \in \mathbb{N}}, \{t_n\}_{n \in \mathbb{N}} \) going to \( +\infty \) so that \( t_n \geq s_n \), for any automorphic form \( q \), we have
\[
\lim_{n \to \infty} \left( \sum_{i \in \{1,2\}} (-1)^i \int_{S(\gamma_i(t_n),X_i)} \langle q, \theta_k^{\gamma_i,\gamma} \rangle \right) d\sigma = 0.
\]

7.5.2. Conclusion of the proof of Theorem 7.2.2. Let \( \Omega \) be the set of those points in \( \partial_\infty \mathbb{H}^2 \) which are end points of geodesics whose projection on \( UX \) have a dense \( \omega \)-limit set. Observe then that any geodesic ending on \( \Omega \) satisfies the latter property. By ergodicity and Poincaré recurrence, \( \Omega \) is nonempty, and hence, by \( \pi_1(\Sigma) \) invariance, dense.

By density, it is enough to prove the result when \((x_1, x_2, X_1, X_2)\) belongs to \( \Omega \). Let \( \gamma_{ij} \) be the geodesic joining \( x_i \) to \( X_j \). Then by hypothesis, for all \( i, j \in \{1,2\} \), there exist
sequences of points \( \{ u^n_{i,j,k} \}_{n \in \mathbb{N}} \) and \( \{ v^n_{i,j,k} \}_{n \in \mathbb{N}} \) of \( \gamma_{i,j} \) converging to \( x_i \) and \( x_j \), respectively, as well as elements \( \eta_{i,j}^n \) in \( \pi_1(\Sigma) \) so that for all \( i, j, k \in \{1, 2\} \),

\[
\begin{align*}
  d(u^n_{i,j,k}, u^n_{i,j,k}) & \xrightarrow{n \to \infty} 0, \\
  d(v^n_{i,j,k}, v^n_{i,j,k}) & \xrightarrow{n \to \infty} 0, \\
  d(u^n_{i,j,k}, \eta^n_{i,j}(v^n_{i,j,k})) & \xrightarrow{n \to \infty} 0 .
\end{align*}
\]

Applying the closing lemma, we deduce from (81) that there exist \( \hat{u}^n_{i,j} \) and \( \hat{v}^n_{i,j} \) on the axes of \( \eta^n_{i,j} \) with \( \eta^n_{i,j}(\hat{u}^n_{i,j}) = \hat{v}^n_{i,j} \) so that

\[
\begin{align*}
  d(u^n_{i,j,k}, \hat{u}^n_{i,j}) & \xrightarrow{n \to \infty} 0 ,

(82)
\end{align*}
\]

\[
\begin{align*}
  d(v^n_{i,j,k}, \hat{v}^n_{i,j}) & \xrightarrow{n \to \infty} 0 .

(83)
\end{align*}
\]

Let us denote by \( \chi(u, v) \) the characteristic function of \( S(u, v) \). By Corollary 7.5.4, as \( n \to \infty \),

\[
\sum_{i,j} (-1)^{i+j} \chi(u^n_{i,j,k}, v^n_{i,j,k}) \theta^k_{x_i x_j} \xrightarrow{\text{cor}} \sum_{i,j} (-1)^{i+j} \theta^k_{x_i x_j} ,
\]

where \( \ast \) denote the convergence in the dual of the space \( C^{0,1}(\mathbb{H}^2) \) of Lipschitz function on \( \mathbb{H}^2 \). Then by (82) and (83), as well as Corollary 7.5.2, we obtain,

\[
\begin{align*}
  \sum_{i,j} (-1)^{i+j} \chi(\hat{u}^n_{i,j,k}, \hat{v}^n_{i,j,k}) \theta^k_{\hat{u}^n_{i,j} \hat{v}^n_{i,j,k}} \xrightarrow{\text{cor}} \sum_{i,j} (-1)^{i+j} \theta^k_{x_i x_j} .

(84)
\end{align*}
\]

As in Proposition 2.1.1, since \( \hat{v}^n_{i,j} = \eta^n_{i,j}(\hat{u}^n_{i,j}) \), for any automorphic form \( q \),

\[
\begin{align*}
  r_k \cdot \int_{S([\hat{v}^n_{i,j}, v^n_{i,j,k}])} \langle q, \theta^k_{\hat{v}^n_{i,j} v^n_{i,j,k}} \rangle \, d\sigma = \int_{[\hat{u}^n_{i,j}, v^n_{i,j,k}]} \hat{q} \, ds ,
\end{align*}
\]

where \( [u, v] \) denotes the geodesic between \( u \) and \( v \). Observe now that

\[
\begin{align*}
  d(\hat{u}^n_{i,j,k}, \hat{v}^n_{i,j,k}) & \xrightarrow{n \to \infty} 0 ,

(85)
\end{align*}
\]

\[
\begin{align*}
  d(\hat{v}^n_{i,j,k}, \hat{v}^n_{i,j,k}) & \xrightarrow{n \to \infty} 0 .

(86)
\end{align*}
\]

Thus we get from Proposition 7.1.2, that

\[
\begin{align*}
  r_k \cdot \sum_{i,j \in \{1, 2\}} (-1)^{i+j} \int_{S([\hat{v}^n_{i,j,k}, v^n_{i,j,k}])} \langle q, \theta^k_{\hat{v}^n_{i,j,k} v^n_{i,j,k}} \rangle \, d\sigma \xrightarrow{n \to \infty} R_h(x_1, x_1, X_2; \hat{q}) .
\end{align*}
\]

Combining this last equation with (84) concludes the proof of the theorem.

7.6. A remark on triple ratios for \( \text{SL}(3, \mathbb{R}) \). For representations into \( \text{SL}(3, \mathbb{R}) \), the 
triple ratio of three points is defined as

\[
T(x, y, w) = b^H(x, y, w, m) \cdot b^H(w, x, y, m) \cdot b^H(y, w, x, m) ,
\]

for all auxiliary points \( m \). It follows from the formulae in Theorem 7.2.2 that the variation of the triple ratio along a cubic differential \( q \) is given by

\[
\frac{d}{dt} \log(T(x_1, x_2, x_3)) = r_3 \int_{\mathbb{H}^2} \langle q(z), \sum_{i \neq j} \theta^3_{x_i x_j} \rangle \, d\sigma(z) .
\]

For \( \text{SL}(n, \mathbb{R}) \) there are \( \frac{(n-1)(n-2)}{2} \) cross ratios, and similar formulas give only a small fraction of all triple ratios.
8. Large \(n\) Asymptotics and Applications

In this section, we describe two phenomena related to the large \(n\) asymptotics of our formula. In particular, we shall see that asymptotically the relation between the pressure and \(L^2\)-metrics becomes remarkably simpler, once we perform a natural renormalization.

8.1. Large \(n\) asymptotics of the pressure metric. For large \(n\) asymptotics, it is more natural to consider the renormalized highest eigenvalue

\[
\mu_{\gamma} = (\lambda_{\gamma})^{1/n}.
\]

and the associated renormalized pressure metric. The reason for this normalization is that then, by (41), the highest eigenvalue does not depend on \(n\) along the Fuchsian locus.

We prove the following asymptotics

**Theorem 8.1.1.** [Large \(n\)-asymptotics] We have the large \(n\)-asymptotics for the renormalized highest eigenvalue given by

\[
\mu_{\gamma} \sim \mu_{\gamma} \cdot (-1)^k \frac{(2k-1)!}{2^{k-2}(k-1)!} \cdot \int_{\gamma} \Re(\bar{q}_k) \, dt.
\]

Finally, the large \(n\) asymptotics for the renormalized pressure metric for the renormalized deformation \(\psi(q)\) associated to a holomorphic \(k\)-differential \(q\) is given by

\[
P(\psi(q), \psi(q)) \sim \frac{(2k-1)!}{2^{k-1} \cdot 3 \cdot \pi |\chi(X)|} \int_X ||q||^2 \, d\sigma.
\]

We need the following consequence of Lemma A.1.

**Lemma 8.1.2.** As \(n \to \infty\) we have the asymptotic expression,

\[
-\text{Tr} \left( E_k^0 F_k^0 \right) \sim n^{2k+1} \cdot \frac{(k!)^2}{(2k+1)!}.
\]

**Proof of Theorem 8.1.1.** Since \(q_k E_{k-1}^0 = q_k \eta_{k-1} E_{k-1}\), where \(\eta_{k-1}\) is as in (16), we apply Theorem 4.0.2 to \(q_k \eta_{k-1}\) to get

\[
\frac{\dot{\mu}_{\gamma}}{\mu_{\gamma}} = \frac{1}{n-1} \frac{\dot{\lambda}_{\gamma}}{\lambda_{\gamma}} = (-1)^k \frac{n-2}{2^{k-2}} \int_{\gamma} \Re \left( q_k^{(\gamma, \ldots, \gamma)} \right) \, dt.
\]

We have \(d(n) \sim n^2/6\). Using (16), and combining with Lemma 8.1.2, this implies

\[
\eta_{k-1} \sim \frac{1}{n^{k-2} (k-1)!} \left( \frac{(2k-1)!}{6} \right)^{1/2}.
\]

The result then follows from the fact that

\[
\frac{(n-2)!}{(n-k)!} \sim n^{k-2}.
\]
8.2. Asymptotic freedom of eigenvalues. Fix $p \in \mathbb{N}$. As previously, let $\lambda_{\gamma}$ denote the largest eigenvalue. Let $\bar{\gamma} = (\gamma_1, \ldots, \gamma_p)$ be pairwise distinct conjugacy classes of primitive elements of $\pi_1(\Sigma)$. We can then construct a map

$$\Lambda_{\bar{\gamma}}^{(n)} : \mathcal{H}(\Sigma, n) \longrightarrow \mathbb{R}^p,$$

$$\delta \mapsto \left( \lambda_{\delta(\gamma_1)}, \ldots, \lambda_{\delta(\gamma_p)} \right).$$

We then have

**Theorem 8.2.1.** For $n$ sufficiently large (depending on $\bar{\gamma}$) the image of $\Lambda_{\bar{\gamma}}^{(n)}$ contains an open set.

This theorem will be an easy consequence of the following, which is of independent interest:

**Proposition 8.2.2.** Let $C$ be the set of conjugacy classes of primitive elements of $\pi_1(\Sigma)$. Let $k_0 \geq 2$ be some integer. Let $\mathbb{R}^C$ be the real vector space freely spanned by elements of $C$. Then the real linear map defined by

$$\mathbb{R}^C \rightarrow \bigoplus_{k=k_0}^{\infty} H^0(X, K^k),$$

$$\gamma \mapsto (\Theta^{(k)}_\gamma \ldots, \Theta^{(k)}_\gamma \ldots). \quad (85)$$

is injective.

8.2.1. **Proof of Proposition 8.2.2.** The proof relies on the following

**Lemma 8.2.3.** Let $\{f_i : i \in \mathbb{N}\}$ be nonzero holomorphic functions defined on the unit disk $D$, and $\{\alpha_i : i \in \mathbb{N}\}$ be a bounded infinite sequence of complex numbers. Assume that:

1. the real analytic functions $|f_i|^2$ are all pairwise distinct;
2. the series $\sum_{i=1}^{\infty} f_i(z)$ is pointwise absolutely convergent for every $z \in D$;
3. there exists a sequence $\{N_m\}_{m \in \mathbb{N}}$, $N_m \rightarrow \infty$, such that for all $m$ and all $z \in D$,

$$\sum_{i=1}^{\infty} \alpha_j \cdot f_i^{N_m}(z) = 0.$$

Then, for all $i$, $\alpha_i = 0$.

**Proof.** Since $\sum_{i=1}^{\infty} f_i(z)$ is pointwise absolutely convergent, it follows that given $z_0 \in D$, there exists:

- an integer $i_0$;
- a finite set $I_0 \subset \mathbb{N}$ containing $i_0$;
- a neighborhood $U_0$ of $z_0$;
- a real number $K_0 > 1$,

such that

\[ (*) \quad \begin{cases} |f_i(z_0)| = |f_{i_0}(z_0)| & \text{for all } i \in I_0 \\ |f_i(z)| \geq K_0 |f_{i_0}(z)| & \text{for all } i \notin I_0, \text{ and all } z \in U_0. \end{cases} \]

If $\# I_0 > 1$, then using the fact that the functions $|f_i|$ are distinct one can find $z_1$ near $z_0$, satisfying conditions $(*)$ for $z_1, I_1, U_1, K_1$ (instead of $(z_0, I_0, U_0, K_0)$) where

$$\# I_1 < \# I_0.$$

In other words, by induction, we may as well assume that $I_0 = \{i_0\}$. 

We now proceed by contradiction. Suppose some \( \alpha_i \neq 0 \). Then after relabeling, we may assume that all \( \alpha_i \neq 0 \). Applying the previous argument we may find \( z_0, K_0 > 1 \), and \( i_0, j_0 \), so that
\[
|f_{i_0}(z_0)| \geq K_0 |f_i(z_0)|
\]  
(86)
for all \( i \neq i_0 \). Then, the equation
\[
\sum_{i=1}^{\infty} \alpha_i \cdot f_i^{N_m}(z_0) = 0 
\]
yields
\[
\alpha_{i_0} = - \sum_{i \neq i_0} \alpha_i \left( \frac{f_i(z_0)}{f_{i_0}(z_0)} \right)^{N_m},
\]
where the right end term converges. Then since
\[
\left| \alpha_i \left( \frac{f_i(z_0)}{f_{i_0}(z_0)} \right)^{N_m} \right| \leq \left| \left( \frac{f_i(z_0)}{f_{i_0}(z_0)} \right) \right| \leq A \cdot |f_i(z_0)|,
\]
where \( A = |f_{i_0}(z_0)|^{-1} \sup_{i \in \mathbb{N}} |\alpha_i| \), Lebesgue dominated convergence theorem yields that
\[
\lim_{m \to \infty} \sum_{i \neq i_0} \alpha_i \left( \frac{f_i(z_0)}{f_{i_0}(z_0)} \right)^{N_m} = 0.
\]
Thus \( \alpha_{i_0} = 0 \); a contradiction. \( \square \)

We can now pass to the proof of Proposition 8.2.2. Let \( \gamma_1, \ldots, \gamma_p \) be different conjugacy classes. We may as well assume that we have an involution \( \sigma \) of \( \{1, \ldots, p\} \) so that \( \gamma_i = \gamma_{i_0}^{-1} \). Assume that there exist \( \alpha_1, \ldots, \alpha_p \) real numbers so that for all \( k \),
\[
\sum_{i=1}^{p} \alpha_i \cdot \Theta_{\gamma_i}^{(k)} = 0.
\]
Let \((x_i, X_i)\) be the end points at infinity of \( \gamma_i \), then we can write
\[
0 = \sum_{i=1}^{p} \alpha_i \cdot \Theta_{\gamma_i}^{(k)} = \left( \sum_{i=1}^{p} \alpha_i \sum_{\eta \in \pi_1(\Sigma)/\langle \gamma \rangle} \eta^* \frac{dz}{(z-x_i)(z-X_i)^k} \right).
\]
(87)
Defining \( g_{i,\eta} \) by
\[
g_{i,\eta}(z) \cdot dz = \eta^* \left( \frac{dz}{(z-x_i)(z-X_i)^k} \right),
\]
we get
\[
0 = \sum_{i=1}^{p} \sum_{\eta \in \pi_1(\Sigma)/\langle \gamma \rangle} \alpha_i \cdot (g_{i,\eta}(z))^k.
\]
Hence, for all even \( k \),
\[
0 = \sum_{i=1}^{p} \sum_{\eta \in \pi_1(\Sigma)/\langle \gamma \rangle} (\alpha_i + \alpha_{i_0}) \cdot (g_{i,\eta}(z))^k.
\]
(88)
Similarly, for all odd \( k \),

\[
0 = \sum_{i=1}^{p} \sum_{\eta \in \pi_1(\Sigma)/(\gamma_i)} (\alpha_i - \alpha_{i(0)}) \cdot (g_{i\eta}(z))^k .
\] (89)

Next we observe that if for all \( z \), \(|g_{i\eta}(z)| = |g_{j\zeta}(z)|\), then \( \eta = \zeta \), and \( i = j \) or \( i = \sigma(j) \). Recall that \( g_{\sigma(i)} = -g_{i} \). Therefore, applying Lemma 8.2.3 to the sequence of even integers, we get from (88), that \( \alpha_i + \alpha_{i(0)} = 0 \) for all \( i \). Similarly, applying Lemma 8.2.3 to the sequence of odd integers, we get from (89), that \( \alpha_i - \alpha_{i(0)} = 0 \) for all \( i \). As a consequence, \( \alpha_i \) vanishes for all \( i \). This concludes the proof of the proposition.

8.2.2. Proof of Theorem 8.2.1. The proof is an immediate consequence of the following

Proposition 8.2.4. Let \( \delta \) be a Fuchsian representation. For \( n \) large enough, \( \Lambda_n^{(n)} \) is a submersion at \( \delta \).

Proof. Let \( \delta \) be a Fuchsian representation with values in \( \text{SL}(n, \mathbb{R}) \), \( \gamma \) be a nontrivial element of \( \pi_1(\Sigma) \), let \( \lambda_{\gamma}(\delta) \) be the highest eigenvalue \( \delta(\gamma) \). Then by Theorem 5.4.1, the Hamiltonian vector field \( H_{\gamma} \) of \( \log \lambda_{\gamma} \) at \( \delta \) is given by

\[
H_{\gamma} = \sum_{k=2}^{n} C(n,k) \cdot \psi \left( i \cdot \Theta_{\gamma}^{(k)} \right),
\] (90)

where \( C(n,k) \) are all nonzero (see Remark 5.4.2), and are independent of \( \gamma \). Now if \( \gamma_1, \ldots, \gamma_p \) are primitive elements of \( \pi_1(\Sigma) \), the proof of the proposition will be complete if we can show that the set of vector fields \( \{H_{\gamma_i}, i=1, \ldots, p\} \) is linearly independent. Let \( W \) be the finite dimensional subspace of \( \mathbb{R}^j \) spanned by the \( \gamma_i, i = 1, \ldots, p \). Then from the finite dimensionality of \( W \) and Proposition 8.2.2, there exists some \( k_1 \) so that the map

\[
\phi_W : W \rightarrow \sum_{k=2}^{k_1} H^0(X, K^k), \quad \gamma \mapsto (\Theta_{\gamma}^{(2)}, \ldots, \Theta_{\gamma}^{(k_1)}),
\]

is injective. Suppose that we have \( \alpha_i \in \mathbb{R} \), so that

\[
\sum_{i=1}^{p} \alpha_i \cdot H_{\gamma_i} = 0 .
\]

Since the \( \psi \left( i \cdot \Theta_{\gamma_i}^{(k)} \right) \) are orthogonal for different values of \( k \), from (90) we obtain that

\[
C(n,k) \sum_{i=1}^{p} \alpha_i \cdot \psi \left( i \cdot \Theta_{\gamma_i}^{(k)} \right) = 0 ,
\]

for all \( k \in \{2, \ldots, k_1\} \). Since \( C(n,k) \neq 0 \), and \( \psi \) is real linear and injective, we find

\[
\sum_{i=1}^{p} \alpha_i \cdot \Theta_{\gamma_i}^{(k)} = 0 ,
\]

and thus

\[
\phi_W \left( \sum_{i=1}^{p} \alpha_i \cdot \gamma_i \right) = 0 .
\]
Finally, by the injectivity of $\phi|_W$, we conclude that $\alpha_i = 0$ for all $i$. This shows that $\{H_{i,\gamma}\}_{i=1}^p$ is linearly independent and finishes the proof. \hfill \Box

**Appendices**

A. **Computation of traces.** For the irreducible embedding $\kappa_n : \mathfrak{sl}(2) \to \mathfrak{sl}(n)$, set $X = \kappa_n(x)$, $Y = \kappa_n(y)$ (see (9)). Then the action of $X, Y$ in the standard basis $\{\hat{w}_p\}$ of $\mathbb{C}^n$ is given by

\[
(-X)(\hat{w}_p) = c_p \cdot \hat{w}_{p+1},
\]

\[
Y(\hat{w}_p) = c_{p-1} \cdot \hat{w}_{p-1}.
\]

where $c_p = (p(n - p)/2)^{1/2}$ (see [23, p. 978]). If we make a change of basis by rescaling $w_p = ((p - 1)!/(n - p)!)^{1/2} \cdot \hat{w}_p$ then

\[
(-X)(w_p) = \frac{1}{\sqrt{2}}(n - p) \cdot w_{p+1},
\]

\[
Y(w_p) = \frac{1}{\sqrt{2}}(p - 1) \cdot w_{p-1}.
\]

The action of $X$ and $Y$ is then identified with that of first order differential operators on homogeneous polynomials of degree $n - 1$ in the variables $z = x + iy$ and $\bar{z} = x - iy$. Namely, identifying $w_p = z^{p-1} \bar{z}^{n-p}$ we have

\[
-\sqrt{2}X = z \frac{\partial}{\partial z}, \quad \sqrt{2}Y = \bar{z} \frac{\partial}{\partial \bar{z}}.
\]

For $k = 1, \ldots, n - 1$, let $E_k^0 = (-\sqrt{2}X)^k$, $F_k^0 = -(\sqrt{2}Y)^k$.

We first compute the following

**Lemma A.1.**

\[
-\text{Tr} \left( E_k^0 F_k^0 \right) = (k!)^2 \binom{n + k}{2k + 1}.
\]

In particular,

\[
-\text{Tr} \left( E_{n-1}^0 F_{n-1}^0 \right) = ((n - 1)!)^2,
\]

and

\[
-\text{Tr} \left( E_1^0 F_1^0 \right) = d(n),
\]

for all $n$ (see (17)).

**Corollary A.2.** Recall the definition (18). Then,

\[
\eta_k = \frac{1}{k!} \left( \frac{n + 1}{3} \right)^{1/2} \left( \frac{n + k}{2k + 1} \right)^{-1/2}.
\]

**Proof of Lemma A.1.** From (91) and (92),

\[
E_k^0 \left( \frac{w_p}{(n-p)!} \right) = \frac{w_{p+k}}{(n-p-k)!},
\]

\[
F_k^0 \left( \frac{w_p}{(p-1)!} \right) = -\frac{w_{p-k}}{(p-k-1)!}.
\]

Hence, for $p \geq k + 1$,

\[
-E_k^0 F_k^0 w_p = E_k^0 \left( \frac{(p-1)!}{(p-k-1)!} \right) w_{p-k} = \frac{(p-1)!}{(p-k-1)!} \frac{(n-k)!}{(n-p)!} w_p.
\]

\[
-\text{Tr} \left( E_k^0 F_k^0 \right) = (k!)^2 \binom{n + k}{2k + 1}.
\]
By summing over $p$, we get

$$-\text{Tr}(E_k^0 P_k^0) = (kt)^2 \sum_{p=k+1}^{n} \binom{p-1}{k} \binom{n-p+k}{k}.$$  

We claim that

$$\sum_{p=k+1}^{n} \binom{p-1}{k} \binom{n-p+k}{k} = (n+k) \binom{2k+1}{k}.$$  \hspace{1cm} (96)

This is actually clear from the combinatorial interpretation. Alternatively, define a generating function

$$f_k(x) := \sum_{m=k}^{\infty} x^m \binom{m}{k}.$$  

Then the left hand side of (96) is coefficient of $x^{n+k-1}$ in $f_k^2(x)$. On the other hand, by (80),

$$f_k(x) = \frac{x^k}{(1-x)^{k+1}},$$  

from which $xf_k^2(x) = f_{2k+1}(x)$. Hence, the coefficient of $x^{n+k-1}$ in $f_k^2(x)$ is the coefficient of $x^n$ in $f_{2k+1}(x)$, which is the right hand side of (96). This proves (96) and hence also the lemma.  

Let $u_p = x^{p-1} u^{n-p}$, and let $\pi_p$ denote the orthogonal projection to the 1-dimensional space spanned by $u_p$. Note that by the hermiticity mentioned above, $\text{Tr}(F_k^0 \pi_p) = -\text{Tr}(E_k^0 \pi_p)$. We then have the following proposition

**Proposition A.3.**

$$\text{Tr}(E_k^0 \pi_1) = \frac{(-1)^k (n-1)!}{2^k (n-k-1)!},$$  \hspace{1cm} (97)

$$\text{Tr}(E_k^0 \pi_p) = \frac{(p-1)! (n-p)!}{2^k (n-k)!} \sum_{j=\max(0, k+p-n)}^{\min(k, p-1)} \binom{n-k-1}{j} \binom{k}{j} (-1)^{i+j}.$$  \hspace{1cm} (98)

**Proof.** Set $a_{n,p} = 2^{1-p} (2i)^{n-p}$, so that $u_p = a_{n,p} (z + \bar{z})^{n-1} (z - \bar{z})^{n-p}$. As a warmup, consider the easy case when $p = 1$ and prove (97). Then using (93),

$$E_k^0 \cdot u_1 = a_{n,1} \left( z \frac{\partial}{\partial z} \right)^k (z - \bar{z})^{n-1}$$

$$= (-1)^k a_{n,1} z^k \frac{(n-1)!}{(n-1-k)!} (z - \bar{z})^{n-k-1}$$

$$= \frac{a_{n,1}}{(-2)^k} \frac{(n-1)!}{(n-1-k)!} (z + \bar{z})^{n-1-k} ((z + \bar{z}) + (z - \bar{z}))^k$$

$$= \frac{a_{n,1}}{(-2)^k} \frac{(n-1)!}{(n-1-k)!} (z - \bar{z})^{n-1} + \sum_{i=2}^{n} A_i \cdot u_i$$

$$= (-2)^k \frac{(n-1)!}{(n-1-k)!} u_1 + \sum_{i=2}^{n} A_i \cdot u_i,$$  \hspace{1cm} (99)
where $A_i$ are some complex numbers. Eq. (97) follows. Let us move to the general case and prove (98).

\[ E^0_a u_p = a_n p z^k \frac{\partial^j}{\partial z^j} \left( (z + \bar{z})^{p-j} (z - \bar{z})^{n-j} \right) \]

\[ = a_n p z^k \sum_{j=0}^{\min(k,p-1)} \binom{k}{j} \frac{\partial^j}{\partial z^j} (z + \bar{z})^{p-j} \left( \frac{(n-p)!}{(n+j-k-p)!} (z - \bar{z})^{n-j} \right) \]

\[ = a_n p z^k \sum_{j=0}^{\min(k,p-1)} (-1)^j \binom{k}{j} \left( \frac{(p-1)!}{(p-j)!} (z + \bar{z})^{p-j} \right) \left( \frac{(n-p)!}{(n+j-k-p)!} (z - \bar{z})^{n-j} \right) \]

\[ = a_n p z^k \sum_{j=0}^{\min(k,p-1)} (-1)^j \binom{k}{j} \left( \frac{(p-1)!}{(p-j)!} (z + \bar{z})^{p-j} \right) \left( \frac{(n-p)!}{(n+j-k-p)!} (z - \bar{z})^{n-j} \right) \]

\[ = k \lambda n_p \sum_{j=0}^{\min(k,p-1)} (-1)^j \binom{k}{j} \left( \frac{(p-1)!}{(p-j)!} (z + \bar{z})^{p-j} \right) \left( \frac{(n-p)!}{(n+j-k-p)!} (z - \bar{z})^{n-j} \right) \]

\[ = k \lambda n_p \sum_{j=0}^{\min(k,p-1)} (-1)^j \binom{k}{j} \left( \frac{(p-1)!}{(p-j)!} (z + \bar{z})^{p-j} \right) \left( \frac{(n-p)!}{(n+j-k-p)!} (z - \bar{z})^{n-j} \right) \]

\[ = k \lambda n_p \sum_{j=0}^{\min(k,p-1)} (-1)^j \binom{k}{j} \left( \frac{(p-1)!}{(p-j)!} (z + \bar{z})^{p-j} \right) \left( \frac{(n-p)!}{(n+j-k-p)!} (z - \bar{z})^{n-j} \right) \]

\[ = k \lambda n_p \sum_{j=0}^{\min(k,p-1)} (-1)^j \binom{k}{j} \left( \frac{(p-1)!}{(p-j)!} (z + \bar{z})^{p-j} \right) \left( \frac{(n-p)!}{(n+j-k-p)!} (z - \bar{z})^{n-j} \right) \]

for some coefficients $A_i$ that we do not need to make explicit. Thus

\[ \text{Tr}(E^0_a u_p) = \frac{k!}{2^k} \sum_{j=\max(0,k+n-p)}^{\min(p-1)} \left( \frac{(p-1)!}{(p-j)!} \left( \frac{(n-p)!}{(n+j-k-p)!} \right) \right) \]

\[ = \frac{(p-1)!}{(p-j)!} \left( \frac{(n-p)!}{(n+j-k-p)!} \right) \]

\[ = \frac{(p-1)!}{(p-j)!} \left( \frac{(n-p)!}{(n+j-k-p)!} \right) \]

\[ = \frac{(p-1)!}{(p-j)!} \left( \frac{(n-p)!}{(n+j-k-p)!} \right) \]

\[ = \frac{(p-1)!}{(p-j)!} \left( \frac{(n-p)!}{(n+j-k-p)!} \right) \]

This concludes the proof of the proposition. \qed

**B. Proof of Theorem 6.4.1.**

**B.1. Closed formula for the integrals.** We shall need the technical formula

**Proposition B.1.** We have

\[ I_{m,d}(R) = \sum_{k=0}^{d} R^{m+2k+1} (1 - R^2)^{d-k} \cdot \frac{(2d)!!(m-1)!!}{(2(d-k))!!(m+2k+1)!!} \cdot (100) \]
Proof. An integration by parts leads to the recursion relation

\[ I_{m,d}(R) = \frac{2d}{m+1} I_{m+2,d-1}(R) + \frac{1}{m+1} R^{m+1} (1-R^2)^d. \]  

(101)

Indeed,

\[
I_{m,d}(R) = \int_0^R t^m (1-t^2)^d \, dt = \frac{2d}{m+1} \int_0^R t^{m+2} (1-t^2)^{d-1} \, dt + \frac{1}{m+1} \left[ t^{m+1} (1-t^2)^d \right]_0^R
\]

\[
= \frac{2d}{m+1} I_{m+2,d-1}(R) + \frac{1}{m+1} R^{m+1} (1-R^2)^d.
\]

Observe that

\[ I_{m,0}(R) = \frac{1}{m+1} R^{m+1}, \]

satisfies equation (100). Moreover, the right hand side of equation (100) satisfies the recursion formula. □

B.2. Asymptotic series. If \( G_m \) and \( H_m \) are positive functions on \([0, 1]\), we write

\[ G_m(R) \sim H_m(R), \]

If

\[
\left\{ \begin{array}{c}
G_m(R) \\
H_m(R)
\end{array} \right\}_{m \in \mathbb{N}} \overset{c^0}{\underset{m \to \infty}{\longrightarrow}} f(R) \quad \text{where} \quad \lim_{R \to 1} f(R) = 1.
\]

We will also use the nonstandard convention that

\[
\left( \begin{array}{c}
n-1 \\
n
\end{array} \right) = \frac{1}{n},
\]

which is coherent with

\[
\left( \begin{array}{c}
n+k \\
n
\end{array} \right) \sim n^k,
\]

for all \( k \geq -1 \). Let also

\[
F(k, p) := \sum_{i+j=k, i,j \leq p-1} \frac{(2p-2)!!}{(2(p-1-i))!!(2(p-1-j))!!} \cdot (2i)!!(m+2i+1)!!(m+2j+1)!!.
\]

(102)

In particular

\[ F(2p-2, p) = (2p-2)!!^2 = 2^{2p-2}(p-1)!^2. \]

We now prove

**Proposition B.2.** We have

\[
\left( \begin{array}{c}
m+2p-1 \\
m
\end{array} \right) I_{m,p-1}^2 = R^{2m} \cdot \sum_{k=0}^{2p-2} R^{2k+2} (1-R^2)^{2p-2-k} \cdot \left( \begin{array}{c}
m+2p-3-k \\
m
\end{array} \right) \cdot F(k, p).
\]

(103)

Proof. From equation (100), we get

\[
I_{m,d}^2(R) = \sum_{k=0}^{2d} \left( R^{2(m+k+1)} (1-R^2)^{2d-k} \cdot \sum_{i+j=k, i,j \leq d} \frac{(2d)!!(m-1)!!^2}{(2i)!!(2(d-i))!!(m+2i+1)!!(m+2j+1)!!} \right).
\]
Since \((m + 2i + 1)!(m + 2j + 1)! \sim m^{k+2}(m - 1)!^2\), we obtain
\[
I_{m,p-1}^2(R) = 2^p \sum_{k=0}^{2p-2} R^{2(m+k+1)} \left( (1 - R^2)^{2p-2-k} \cdot \frac{1}{m^{k+2}} \sum_{i+j=k, i,j \leq p-1} \frac{(2p-2)!^2}{(2(p-1-i)!!(2(p-j-1))!!} \right),
\]
Recall that
\[
b_m = \left( m + \frac{2p - 1}{m} \right) \sim m^{2p-1}.
\]
Then by (102)
\[
b_m \cdot I_{m,p-1}^2(R) = 2^p \sum_{k=0}^{2p-2} R^{2(m+k+1)} \left( (1 - R^2)^{2p-2-k} \cdot m^{2p-3-k} \cdot F(k, p) \right)
\]
This finishes the proof of the proposition. □

B.3. A preliminary lemma. We need an easy lemma.

**Lemma B.3.** Assume that the positive functions \(G_m(R)\) and \(H_m(R)\) satisfy
\[
\begin{align*}
G_m(R) & = H_m(R), \\
\forall m, \lim_{R \to 1} H_m(R) & = \lim_{R \to 1} G_m(R) = 0, \\
\lim_{R \to 1} \sum_{m=0}^{\infty} H_m(R) & = D < \infty,
\end{align*}
\]
then
\[
\lim_{R \to 1} \sum_{m=0}^{\infty} G_m(R) = \lim_{R \to 1} \sum_{m=0}^{\infty} H_m(R).
\]
**Proof.** Let \(\varepsilon > 0\), from the hypothesis (104), there exist \(0 < \alpha < 1\) and \(k_0\) such that for all \(m > k_0\) and \(\alpha < R < 1\)
\[
\left| 1 - \frac{G_m(R)}{H_m(R)} \right| \leq \varepsilon.
\]
Thus
\[
\left| \sum_{m=k_0}^{\infty} (G_m(R) - H_m(R)) \right| \leq 2\varepsilon \cdot D. \quad (105)
\]
Using (104), we now choose \(\alpha\) so that for all \(m \leq k_0\) and all \(R \geq \alpha\),
\[
\left| \sum_{m=0}^{k_0} (G_m(R) - H_m(R)) \right| \leq \varepsilon. \quad (106)
\]
Combining equations (105) et (106) one gets that for all \(R > \alpha\), one has
\[
\left| \sum_{m=0}^{\infty} (G_m(R) - H_m(R)) \right| \leq \varepsilon \cdot (2D + 1).
\]
This last assertion concludes the proof of the lemma. □
B.4. Proof of Theorem 6.4.1. We now prove the theorem.

Proof. Let
\begin{align*}
G_m(R) &:= -\binom{m + 2p - 1}{m} \frac{f^2_{m,p-1}}{\log(1 - R)} \\
H_m(R) &:= -\frac{R^{2m}}{\log(1 - R)} \sum_{k=0}^{2p-2} R^{2k+2}(1 - R^2)^{2p-2-k} \cdot \left(\frac{m + 2p - 3 - k}{m}\right) \cdot F(k, p)
\end{align*}

From proposition B.2, one gets that
\[ G_m(R) \approx H_m(R). \]

Using the asymptotic expansion, for \(0 \leq k \leq 2p - 3\),
\begin{align*}
(1 - R^2)^{2p-2-k} F(k, p) \cdot \sum_{m=0}^{\infty} R^{2m} \left(\frac{m + 2p - 3 - k}{m}\right) &= (1 - R^2)^{2p-2-k} \frac{F(k, p)}{(1 - R^2)^{2p-2-k}} \\
&= F(k, p).
\end{align*}

Similarly for \(k = 2p - 2\), we get that
\begin{align*}
(1 - R^2)^{2p-2-k} F(2p - 2, p) \cdot \sum_{m=1}^{\infty} R^{2m} \left(\frac{m - 1}{m}\right) &= F(2p - 2, p) \cdot \sum_{m=1}^{\infty} R^{2m} \cdot \frac{1}{m} \\
&= F(2p - 2, p)(1 - \log(1 - R)).
\end{align*}

It follows that (with the convention that \(H_0(R) := 0\))
\[ \sum_{m=0}^{\infty} H_m(R) = R^{2k+2} \cdot \left(\frac{F(2p - 2, p)}{\log(1 - R)} - \sum_{k=0}^{2p-2} F(k, p) \cdot R^{2k+2} \right). \]

In particular
\[ \lim_{R \to 1} \sum_{m=0}^{\infty} H_m(R) = F(2p - 2, p) = 2^{2p-2}(p - 1)!^2. \]

We now observe that \(H_m(R)\) and \(G_m(R)\) satisfies the hypothesis \((104)\) of the previous lemma, thus
\[ \lim_{R \to 1} \sum_{m=0}^{\infty} G_m(R) = 2^{2p-2}(p - 1)!^2. \quad (107) \]

Thus we obtain Theorem 6.4.1 as a consequence of \((107)\). \(\square\)

References


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