

SORBONNE UNIVERSITÉ

MASTER THESIS

Teichmuller theory and Thurston Earthquake flow

Bonnafoux Etienne

supervised by Carlos Matheus

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1 Introduction

Since Bernhard Riemann, mathematicians knew that a geometric hyperbolic surface can be described with only a finite number of parameters. Knowing this, for a given surface, we can be interested in the set of all geometries we can give it, modulo composition by map isotopic to the identity map, this set is called the Teichmuller space. Other problems rise shortly after this definition. How can we deform in a natural way a surface's geometry into another? What does it mean that a geometry is close to another? What are the natural boundaries of the Teichmuller space?

Oswald Teichmuller, a german mathematician, studied and answered this question in the year preceding the second World War. He created the first metric on this space by finding a solution to an extremal problem: between two hyperbolic geometry on the same surface is there a function which minimize the deformation? The resulting theorem [24] proves not only the existence, but also the unicity of this function. It naturally induced a distance in the now called Teichmuller space by considering the logarithm of the deformation of the extremal function.

Thurston then added other important steps to this theory. He underlined the role of laminations. A lamination is a generalisation of simple closed curve. And he introduced the earthquake flow, in a course in 1976-1977 at Priceton, which plays an important role in Teichmuller theory . Kerckhoff used this tool to show the Nielsen realisation conjecture in 1983 [13] which states that every finite subgroup of the mapping class group have a fixed



Figure 1: The mathematician Oswald Teichmuller

point in the Teichmuller space. This theory have been actively studied and therefore generated a lot of literature and contributions. One can refer to the following books [6], [11], [12] or course notes [18] [19] for a more detailed study of the theory.

Still, some questions remain open, raising further interest of the community in the topic of hyperbolic geometry. Among the main questions encountered, figures the asymptotic number of closed geodesics. To begin with, one can wonder what is the number $\pi(X, L)$ of closed geodesic on a hyperbolic surface X of length less than L. The answer was found in the 1940s and 1950s by Delsarte, Huber and Selberg and -due to its resemblance to the prime number theorem- is called the prime number theorem for hyperbolic surfaces. It states that:

$$\pi(X,L) \simeq e^L/L$$

as $L \to \infty$. The reader can refer to [4] for more details.

A much harder problem was to find the number, $\sigma(X, L)$, of simple (which don't intersect themselves) closed geodesic of length less than L on a hyperbolic surface X. It was found years later, in Mirzakhani's PhD [22], and we have

$$\sigma(X,L) \simeq C_X L^{6g-6}$$

As $L \to \infty$ where g is the genus of the surface X and C_X is a constant which depend of the geometry X.

To obtain such a result, Maryam Mirzakhani conjugated the earthquake flow to the horocycle

flow. This step provides that the Earthquake flow is ergodic and allows us to use Birkoff theorem to understand asymptotic quantities.

The question is now to give error terms to this quantity. In order to do that, we need to understand better the mixing rate of the earthquake flow. This flow is conjugate to the horocyclic flow which has a polynomial mixing rate. But as the conjugacy is only a measurable map it transport the rate of mixing on a subset of the L^2 of the moduli space which could be diffult to use. We should analyse the rate of mixing of the earthquake flow by other means. One research direction would be to consider natural functions on Teichmuller space such as the systole. The systole is the lenth of the shortest simple closed curve on the surface. This function behaves nicely along earthquake path, it is continuous, convex and we know its first derivative at the origin. Moreover in the case of the once punctured torus we can give a frame determined by the continued fraction of the slope.

In this master thesis, I will first give an introduction to Teichmuller theory and some useful and classical tools in this theory such as the collar lemma or a surface's decomposition in pair of pants. Then, we will review the proof of the Mirzakhani's conjugacy between the horocyclic flow and the earthquake flow. After doing so, we will focus on the mixing properties of the geodesic and horocyclic flow, and discuss their mixing rate. Finally we will discuss a special case: the once punctured torus, as it is one of the simplest example of hyperbolic surface.

2 Notations

 \mathcal{ML} measured lamination

 $\mathcal{MF}(S)$ space of all equivalence classes of measured foliations.

 \mathcal{QD} bundle of nonzero quadratic differential. \mathcal{T}_g Teichmuller space of surface of genus g Mod(S) The modular group of a surface S

 $\mathcal{P}^1\mathcal{M}_g$ Product of \mathcal{T}_g and $\mathcal{MF}(S)$ quotiented by the modular group.

 $\mathcal{GC}(S)$ the space of geodesic currents

 \mathbb{D} the Poincaré disk

 \mathbb{S}_{∞} the unit circle seen as the boundary of \mathbb{D}

 $\bar{\mathbb{D}}$ $\mathbb{D} \cup \mathbb{S}$

S set of all simple closed curve of a surface

 $D_{\gamma}(\alpha)$ Dehn's twist of α around γ

3 Teichmuller Theory

3.1 First definitions

To begin we will in this section give the very definition to introduce the theory of Teichmuller space

Definition 3.1. Let S be a surface of genus g, a marking of S is a couple (X, f) made of a closed Riemann surface X and of one homeomorphism $f: S \to X$ which preserve the orientation. On the set of the marking S, we have a equivalence relation, $(X_1, f_1) \sim (X_2, f_2)$ if there exist $\alpha: X_1 \to X_2$ such that $f_2 \circ \alpha \circ f_1^{-1}$ be an homeomorphism of S preserving the orientation and isotope to the identity map. The set of the marking quotient by the previous relation is the *Teichmuller space* and is written \mathcal{T}_g .

Remark. If $g \geq 2$, for every closed curve α of S, there is only one closed geodesic of X freely isotope to $f(\alpha)$. Moreover if α is simple, so is the geodesic. We will note $l_{\alpha}(X)$ its hyperbolic length and we will take the weakest topologie on \mathcal{T}_g which make this map continuous.

This give a map $L: \mathcal{T}_g \to \mathbb{R}^{\mathbf{S}}$, where \mathbf{S} is the set of all simple closed curve of a given surface. We can ask if this map si injective, i.e. that a geometry is given by the length of the set of simple closed curve. The answer is yes and more precisely one can choose only 9g - 9 + 3n curve so that this map is injective, as stated in Farb and Margalit (2011) Theorem 10.7. [6] This give the intuition that Teichmuller space can be describe only by using a finite set of parameter. We will see after that other coordinates which have nice propeties for other use.

An other curious but important fact about simple closed geodesic is that the sum of a function of theirs length is equal to a constant which does not depend on the genus or the geometry of the surface.

Theorem 3.1 (McShane's identity). Let X be a hyperbolic surface then

$$\sum_{\gamma} \frac{1}{1 + e^{l_{\gamma}(X)}} = \frac{1}{2}$$

Where the sum is taken over all simple closed geodesic.

This theorem was proven by McShane and therefore bears his name [20]. A good introduction with a proof using probablistic method on trees can be found in a paper by F. Labourie and S.P. Tan [14].

Definition 3.2. The modular group is the group of homeomorphisms of S which respect orientation quotiented by the subgroup of homeomorphisms isotopic to the identity map. We will call this group Mod_g It act discretely on the Teichmuller space \mathcal{T}_g and the quotient is called the moduli space and is written \mathcal{M}_g .

The Teichmuller space is the universal covering of the modular space. In fact the Teichmuller space is homeomorphic to a sphere and so is contractible and in particular, connected and simply connected. It is the universal covering of the moduli space.

We can ask ourselves how these spaces look and if we can give an easy representation of them. The fact is that the modular space is not manifold but an orbifold, that is a space which locally looks like a ball of a vector space quotiented by a finite group.

We will now describe objects that are natural to understand the Teichmuller space.

Definition 3.3. A foliation on a surface S is the collection of the following data: a finite set of points $P = (p_1, p_2, ...)$, given a open covering U_i on S - P, a collection of C^1 real function ν_i such than $\|d\nu_j\| = \|d\nu_i\|$ on $U_i \cap U_j$, and near each singular point p_s a coordinate neighborhood V with complex coordinate z such that $\|d\nu\| = \|Im(z^{\frac{k}{2}}dz)\|$ for some positive integer k called the degree of the singular point.

Leaves of the foliations are the graphs immersed in S along which dv is constant. In addition if the surface S have a boundary it is required that it is contained in a singular leaf.

The height $h_{\gamma}(\|d\nu\|)$ (of a free homotopy class) of a loop γ on S is the infimum in the homotopy class of the integral by $\|d\nu\|$

$$h_{\gamma}(\|d\nu\|) = inf_{\gamma \sim \gamma'} \int_{\gamma'} \|d\nu\|$$

The topology on the measured lamination that we will use is the weakest that makes the height functions continuous.

Remark. We won't actually study the set of measured foliation but the equivalence class of

$$h_{\gamma}(\|d\nu\|) = h_{\gamma}(\|d\mu\|), \text{ for each loop } \gamma \in S$$

. We can equivalently use Whitehead equivalence relation on singular foliations by collapsing critical intervals to points and taking isotopy of foliation.

Let $\mathcal{MF}(S)$ be the space of all equivalence classes of measured foliations.

Definition 3.4. A lamination λ is a closed set made of an union (non necessarily finite) of geodesics. For each point x in λ , passes only one geodesic of the lamination. We will write this space $\mathcal{ML}(x)$.

To understand the space of geodesic laminations it is sometimes easier to understand the "trace" of this space in the boundary circle of the universal cover of the surface. This observation gives the following definition.

Definition 3.5. Let \mathcal{M}_{∞} be the space of unordered pairs of distinct points in \mathbb{S}^1

$$\mathcal{M}_{\infty} := \{(z, w) \in \mathbb{S}^1 \times \mathbb{S}^1, z \neq w\} / / (z, w) \equiv (w, z)$$

Let G be a discrete torsion-free group in $PSL(2,\mathbb{R})$ such that $\mathbb{H}//G = S$ is a hyperbolic surface. A geodesic current μ on S is a G-invariant Radon measure on \mathcal{M}_{∞} . We will note $\mathcal{GC}(S)$ the space of geodesic currents, see [12]

Remark. $\mathcal{GC}(S)$ have a natural topology which is the weak * convergence on continuous functions.

Example 3.1. Let $g \in G$ be an element that represent a simple closed geodesic on the surface S. The G-orbit of the fix point of g is a discrete set on which we can put Dirac measure. This links the geodesic currents to the measured lamination.

Remark. A multicurve is a formal sum of geodesics $\gamma = \sum a_i \gamma_i$. The space of lamination is in some aspect the closure of the set of all multicurve

Definition 3.6. Consider the square $\mathcal{M}_{\infty}^2 := \mathcal{M}_{\infty} \times \mathcal{M}_{\infty}$. In this space we can consider the open subset \mathcal{IM}_{∞}^2 corresponding to pairs of geodesics which have transversal intersections in \mathbb{H} . G acts on \mathcal{IM}_{∞}^2 . If μ and ν are geodesic currents in $\mathcal{GC}(S)$, the product $\mu \times \nu$ defines a G-invariant measure on \mathcal{IM}_{∞}^2 . Finally if we take the mass of the total space $\mathcal{IM}_{\infty}^2//G$, the result is called the *intersection number*, $i(\mu, \nu)$

Proposition 3.2.

$$i: \mathcal{GC}(S) \times \mathcal{GC}(S) \to \mathbb{R}_+$$

is continuous and bilinear [3].

Remark. If α and β are simple closed geodesics (Dirac measure in $\mathcal{GC}(S)$), then the intersection number is the number of intersection between α and β . Actually, one can define intersection in this way, first on simple closed geodesic, then by bi-linearity on multi-curves and finally by continuity on geodesic current.

Remark. The topology on \mathcal{ML} is the weakest that make i(.,.) a continuous function.

Definition 3.7. A quadratic differential is a section of the square of the tangent space to X. It is locally as $\phi = \phi(z)dz^2$.

Remark. If $\phi(p) \neq 0$ we can find a map including p in which $\phi = dz^2$. Hence ϕ determine a flat metric on X and a foliation \mathcal{F} corresponding to horizontal lines.

A quadratic differential is said to be integral if

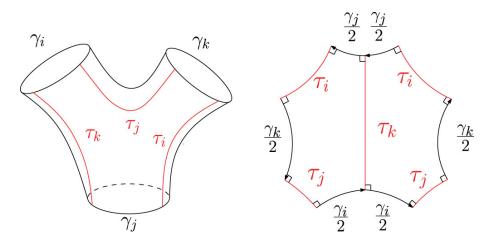
$$\|\phi\| = \int_X |\phi| < \infty$$

We will write Q(x) the Banach space of integral quadratic differentials.

3.2 Decomposition of hyperbolic surface

One way to construct all hyperbolic surface is to decompose them in elementary piece, that we will call pair of pant.

A hyperbolic geometric exercise shows that a hexagon which side are geodesics and with right angles is determined by the length of three sides which are not consecutive.



On the image above γ_i , γ_j and γ_k determined the hexagon. Then we can glue them to have a pair of pant.

Definition 3.8. A pair of pant is a hyperbolic surface with three geodesic boundaries and no punctures.

Remark. The pair of pant is uniquely determined by the length of the three boundary geodesics.

Proof. As a hyperbolic polygon is uniquely determined by three non consecutive sides, the surface with boundaries that we obtain by gluing two together is determined by the length of its three boundary geodesics. \Box

Remark. The length of one or more geodesic can go to zero and the boundaries become a puncture.

We can now decompose, with the following theorem, all hyperbolic surfaces in a collection of pair of pants.

Theorem 3.3. Let S be a surface of genus g and with n punctures. There is a set of 3g-3+n simple closed curves $(\gamma_1, ..., \gamma_{3g-3+n})$ such that $S \cup \gamma_i$ is a disjoint collection of pair of pants.

Definition 3.9. Given a surface S and a pant decomposition $\gamma_1, ..., \gamma_{3g-3+n}$, we have a map

$$(S) \to (\mathbb{R}^{+3g-3+n}, \mathbb{R}^{3g-3+n})X \mapsto (l_{\gamma_1}(X), ..., l_{\gamma_{3g-3+n}}(X), \tau_{\gamma_1}(X), ..., \tau_{\gamma_{3g-3+n}}(X))$$

This map is injective and is call the Fenchel-Nielsen coordinates.

Definition 3.10 (Dehn twist). Let γ be a simple closed curve. There is a tubular neighbourhood of γ called A homeomorph to $[0;1] \times S^1$. A *Dehn's twist* around γ is the homeomorphism which is the identity out of A and is $(t,s) \mapsto (t,e^{2i\pi t}s)$ on A.

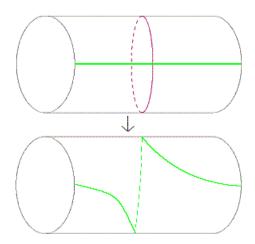


Figure 2: A Dehn twist

We can now give a set of generator of the modular group.

Remark. The Lickorisk theorem states that the modular group is generated by the Dehn's twist and more precisely that one can choose only 2g + 1 generators [16]. We will give here a easier version of this theorem.

Theorem 3.4. There is a collection $\delta_1, ..., \delta_{9g-9}$ of simple closed curves such that $\mathcal{T}_g \to \mathbb{R}^{9g-9}$ is injective.

Proof. Let us take $(\gamma_1, ..., \gamma_{9g-9})$ a decomposition in pair of pants, let $(\alpha_1, ..., \alpha_{9g-9})$ be a collection of simple closed curves such that $i(\gamma_i, \alpha_i) > 0$ and $i(\gamma_i, \alpha_j) = 0$ for $i \neq j$, finally we take $\beta_i = D_{\gamma_i}(\alpha_i)$. We want to show that the length of this collection of 9g - 9 curves determined the hyperbolic structure $X \in \mathcal{T}_g$. X already has the Fenchel-Nielsen coordinate, of the pant decomposition $(\gamma_1, ..., \gamma_{3g-3})$, $(l_{\gamma_1}(X), \tau_{\gamma_1}(X), ..., l_{\gamma_{3g-3}}(X), \tau_{\gamma_{3g-3}}(X))$ so we need only to show that the parameters $\tau_{\gamma_i}(X)$ are determined by the length of the collection. Up to a re-normalisation we can take $\tau_{\gamma_i}(X) = 0$ for every i. Now let's take $t = (t_1, ..., t_{3g-3}) \in \mathbb{R}^{3g-3} \setminus 0$. We will note X_t the hyperbolic geometry which has the same length as X and twist parameters t in the Fenchel-Nielsen coordinate of the pair of pants. So $X_0 = X$. We will consider the function $A(t) = l_{\alpha_1}(X_t)$ and $B(t) = l_{\beta_1}(X_t)$. This function depend only of t_1 as $i(\gamma_i, \alpha_j) = 0$

for $i \neq j$. Moreover they are strictly convex and by definition we have $A(t_1 + l_{\gamma_1}(X)) = B(t_1)$. We will show that there is no $t_1 \neq 0$ such that $A(t_1) = A(0)$ and $B(t_1) = B(0)$ that is $A(t_1 + l_{\gamma_1}(X)) = A(l_{\gamma_1}(X))$. We will note $s = t_1$ and $L = l_{\gamma_1}(X)$ Suppose there is $s \neq 0$ such that A(s) = A(0), we can take s > 0, the other case is symmetric. If s < L, then by convexity for every $t \in]0$; s[, A(t) < A(0) = A(s) and A is strictly increasing after s so as s < L < L + s we have A(L) < A(L + s). If s > L, then L < s < L + s and A(L) < A(L + s). The final case s = L is also impossible since we would have A(0) = A(L) = A(2L). Finally we can make the same argument for the other twist parameters which conclude the proof.

3.3 Flow on Teichmüller space

We will define the main object of this document, earthquake flow.

Definition 3.11. The *earthquake flow* is family of maps defined for $t \in \mathbb{R}$

$$E_t: \mathcal{ML} \times \mathcal{T}_g \to \mathcal{ML} \times \mathcal{T}_g$$

$$(\lambda, X) \mapsto (\lambda, E_{t\lambda}X)$$

where $E_{t\lambda}$ is first defined if λ is a simple closed curve. In this case, we open the surface along lambda, then twist the left part of t unit and put it together again. Then if γ_1 and γ_2 are two curves that don't intersect, E_{γ_1} and E_{γ_2} commute. So we can define the earthquake map on weighted multicurves by twisting one curve after the other by amounts proportional to the weight. Finally as multicurves are dense in the set of lamination we can extend this map for any lamination [13].

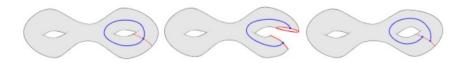


Figure 3: Effect of a twist on a transverse curve, image from [26]

Remark. If the lamination is just a simple closed curve γ then $E_{l_{\gamma}(X)}(X,\gamma)$ is just a Dehn twist around γ . Moreover if we take a decomposition in a pair of pant that contains γ , it is just a translation in the coordinate of the twist of γ .

One could give a precise proof that the extension of the earthquake map from the multicurves to the lamination is rigorous, that is if we take two sequence α_n and α'_n of multicurves which converge to the same lamination, then the sequences of earthquake map along this multicurves converge to the same map.

We will work on universal covering of the surface, the half-plane \mathbb{H} If $v \in T^1\mathbb{H}$ we will call $\gamma(v)$ the geodesic passing by the base-point of v and with v as tangent vector at this point.

We have two useful lemmas to control the distortion of the earthquake map. We will not prove them, but their proofs are in [13].

Lemma 3.5. Let l and l' be two geodesics and $x \in l$, $y \in l'$ two points at most ϵ apart. Then if v and v' are the two geodesic tangents to l and l' respectively at points x and y, we have $d(v,v') < C\epsilon$ for a universal constant C.

Lemma 3.6. Let v and v' be two vectors such that $d(v,v') < \epsilon$, let denote $\gamma = \gamma(v)$ and $\gamma' = \gamma(v')$. Let $w \in T^1\mathbb{H}$, then

- 1. $d(E_{t\gamma}w, E_{t\gamma'}w) \leq Kt\epsilon$
- 2. $d(E_{t\gamma}w, w) \leq Kt$

for all $t \leq T$ and for a constant K depending on T and on the distance between the base-point of v and w.

With these two lemma, we can describe what happens if we change a discrete lamination by a simple closed curve that average it.

Lemma 3.7. Let $x, y \in \mathbb{H}$, $v \in T_1\mathbb{H}$ based at y, \bar{A} the geodesic from x to y. Suppose γ is a discrete lamination with equal measure on each leaf whose intersection with \bar{A} is included in a subarc A. Let l be a single geodesic intersecting A with angle equal to the average angle of the intersections of γ and A and with mass equal to $\mu = i(A, \gamma)$.

If A has length less than δ then for every $T \in \mathbb{R}^+$ the distance between $E_{t\gamma}v$ and $E_{tl}v$ is less than $Kt\mu\delta$, for all $t \leq T$ and K is a constant depending only of $T\mu$ and d(x,y).

Proof. Let denote $l_1, l_2, ..., l_n$ the leaves of γ . The image of \bar{A} is a disconnected arc. The point x is connected to $E_{t\gamma}y$ by a staircase path going along component of $\bar{A}\setminus\gamma$ and subarc of l_i . Let denote the successive components $A_0, A_1, ..., A_n$ and δ_i the length of A_i . So the staircase path is moving by δ_0 along A_0 then by μ/n along l_1 , and so on.

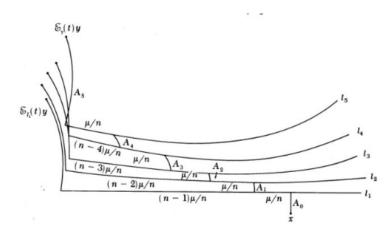


Figure 4: Image of the construction described, image from [13]

We now alter the path by replacing the shearing along l_n distance μ/n by a shearing along l_{n-1} distance $2\mu/n$. The change is less than $Kt^{\mu}_{n}C\delta_{n-1}$ by lemma 3.5 and lemma 3.6. Then we change the shearing by a shearing along l_{n-2} of distance $3\mu/n$ and we continue until we shear a distance μ along l_1 . The total change is less than $KtC\sum_{i=1}^{n-1}\frac{i\mu}{n}\delta_{n-i}$ which is less than $KCt\mu\delta$. We now pass from l_1 to l with lemma 3.5 and 3.6, and we obtain the new lemma.

We need a final lemma to conclude of the well-founded definition of the earthquake flow. In this case we control the difference after earthquaking by two simple geodesic going through the same point of the arc.

Lemma 3.8. Let x, y, v, and \bar{A} be as above, if l and l' are geodesics of \mathbb{H} with measure μ and μ' such that $l \cup \bar{A} = l' \cup \bar{A} = p \notin \{x, y\}$ and the difference between the vectors tangent of l and l' at p of length μ and μ' is less than ϵ .

Then for any T, $d(E_{tl}v, E_{tl'}v) < Kt\epsilon$, for $t \leq T$ and K a constant which depends only on d(x,y) and $T\mu$

Finally this give the theorem which control the distance between two earthquake paths.

Proposition 3.9. Let $\nu \in \mathcal{ML}$ be a lamination and let x, y be in \mathbb{H} , A be the geodesic from x to y and $v \in T^1\mathbb{H}$ be based at y, and x and y do not lie on the atomic part of ν . Then for any ϵ , T, there is a neighbourhood U of ν in \mathcal{ML} such that for all γ , $\bar{\gamma}$ weighted multicurve in U, $d(E_{t\gamma}v, E_{t\bar{\gamma}}v) < Kt\epsilon$, for all $t \leq T$, K a constant depending only on d(x, y) and $Ti(\nu, A)$

Corollary. The earthquake flow is well defined along any lamination $\nu \in \mathcal{ML}$ and for all time t.

Remark. The earthquake flow is an isometry outside the support of the lamination and is continuous outside the atomic part, i.e. the simple closed geodesics of the lamination.

Thurston show that given two point in the Teichmüller space, there is a lamination λ such that the earthquake flow from one point with respect to λ reach the other point (also in [13]). We can ask ourselves what is an invariant measure of this flow.

Definition 3.12. The Weil-Peterson form is the form

$$\omega_{WP} = \sum dl_i \wedge d\tau_i$$

Where $(l_1, ..., \tau_1)$ are the Fenchel-Nielsen according to a pant decomposition. This induces a measure μ_{WP} .

There is a finite measure ν_g in the Lebesgue measure class on $\mathcal{P}^1\mathcal{M}_g$ which is invariant under the earthquake flow. This measure projects to the volume form given by $B(X) \times \mu_{WP}$ on \mathcal{M}_g , where

$$B(X) = \mu_{Th}(\lambda \in \mathcal{ML}, l_{\lambda}(X) \le 1)$$

There are two other important flows, the geodesic flow and the horocyclic flow. First there is a natural homeomorphism between $T^1\mathbb{H} \simeq SL_2(\mathbb{R})/SO_2(\mathbb{R})$, since $SL_2(\mathbb{R})/SO_2(\mathbb{R})$ act simply transitively on it. This morphism can be chosen up to a conjugation via an other element of $SL_2(\mathbb{R})/SO_2(\mathbb{R})$. We we will be interested in a special kind of subgroup.

Definition 3.13. A fuchsian group Γ is a finitely generated and discrete subgroup of $SL_2(\mathbb{R})$. Then Γ acts discontinuously on \mathbb{H} .

A Hyperbolic surface can be represented as $\mathbb{H}/\Gamma = SL_2(\mathbb{R})/SO_2(\mathbb{R})/\Gamma$. If U is a one parameter subgroup of $SL_2\mathbb{R}$ it acts on the quotient.

There are two important example:

Definition 3.14. The *geodesic flow* is a flow on the Teichmuller space given by the action of the diagonal matrices

$$u_t = \begin{pmatrix} e^t & 0\\ 0 & e^{-t} \end{pmatrix}$$

Definition 3.15. The horocycle flow is a flow on the bundle of non-zero quadratic differential, QD, of the Teichmuller space given by the unipotent action of

$$u_t = \begin{pmatrix} 1 & t \\ 0 & 1 \end{pmatrix}$$

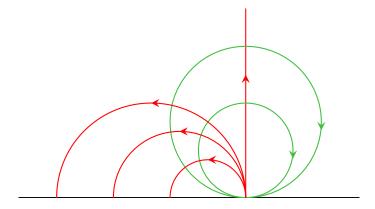


Figure 5: Representation of the horocycle flow, in green, and the geodesic flow, in red

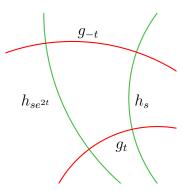


Figure 6: The conjugaison of the horocycle flow by the geodesic one.

The geodesic flow is also the flow that we obtain by following geodesic line on \mathbb{H} and the horocycle is the flow we obtain by following curves which are everywhere orthogonal to the geodesic, which is the horizontal line and the circle tangent to the real line.

An important relation is how this two flows interact between each other

$$\begin{pmatrix} e^t & 0 \\ 0 & e^{-t} \end{pmatrix} u_t = \begin{pmatrix} 1 & s \\ 0 & 1 \end{pmatrix} \begin{pmatrix} e^{-t} & 0 \\ 0 & e^t \end{pmatrix} = \begin{pmatrix} 1 & se^{2t} \\ 0 & 1 \end{pmatrix}$$

So the the conjugation of the horocyclic flow by the geodesic one is still the horocyclic flow.

We give the definition of a third flow that we will use to demonstrate the ergodicity of the two previous one.

Definition 3.16. The rotational flow is a flow on the bundle of nonzero quadratic differential, QD, of the Teichmuller space given by the action of

$$e_{\theta} = \begin{pmatrix} \cos(\theta) & \sin(\theta) \\ -\sin(\theta) & \cos(\theta) \end{pmatrix}$$

We have the relation

$$h_s = e_{\theta} g_t e_{\pi + \theta}$$

3.4 The collaring theorem

We will now give a useful tool to give necessary condition on length of two intersecting geodesics.

The collar function $\eta:]0;\infty[\to]0;\infty[$ is defined as follow. We draw a segment of length l>0 on a geodesic γ , then we project perpendicularly to the geodesic the end of this segment to infinity and draw the geodesic δ which have this endpoint. So we have $\eta(l)=d(\gamma,\delta)$.

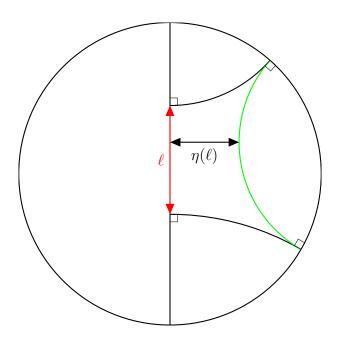


Figure 7: In red the segment of length l and in green the geodesic whose endpoint are orthogonal projection of the end of the segment.

It is an exercise to show that:

$$\eta(l) = \frac{1}{2} ln(\frac{\cosh(l/2) + 1}{\cosh(l/2) - 1})$$

This quantity will be the side of a long "tube" generated by a simple closed geodesic in the hyperbolic surface. We give a definition to make this a little more precise.

Definition 3.17. Let γ be a simple closed geodesic of length l on a hyperbolic surface X. If the δ -neighbourhood

$$A_{\delta}(\gamma) := \{x \in X | d(x, \gamma) < \delta\}$$

is isometric to the δ -neighbourhood of the unique simple closed geodesic on the cylinder of modulus $\frac{\pi}{l}$, we say that γ admit a δ -collar, or that $A_{\delta}(\gamma)$ is the δ -collar of γ .

We can now state a useful theorem.

Theorem 3.10. Let X be a complete hyperbolic surface, and let $\Gamma := \gamma_1, ...$ be a collection of disjoint simple closed geodesic, each γ_i of length l_i . Then $A_{\eta(l_i)}(\gamma_i)$ are collars around the γ_i , and they are disjoint.

Proof. Choose γ_1 and γ_2 and add other simple closed curve to have a maximal multicurve that includes both. Now cutting along this curve we have a set of pair of pants so we only have to show that the $\eta(l_i)$ neighbourhood of γ_i the boundaries of the pair of pants do not intersect each other. We cut the pair of pants along geodesics coming from a boundary C and meeting the two other boundaries A and B. We unfold this figure in the hyperbolic plane and name the side of the octagon following the figure below.

Since a' and b' have the common perpendicular C', they do not intersect and similarly for a'' and b''. The theorem follow easily by the definition of the function η .

There are some corollaries which follow from this theorem and are ready to use in many occasions.

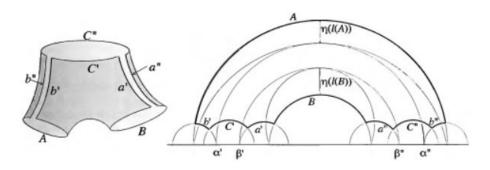


Figure 8: image from [11]

Corollary. Let X be a hyperbolic surface, and γ_1 , γ_2 two simple closed geodesics on X of lengths l_1 and l_2 . If $l_2 < 2\eta(l_1)$, then either $\gamma_1 = \gamma_2$ or $\gamma \cap \gamma_2 = \emptyset$

Proof. If $\gamma_1 \neq \gamma_2$ and $\gamma \cap \gamma_2 \neq \emptyset$ then γ_2 must cross the collar neighbourhood of γ_1 from one boundary to the other and so have length strictly superior than $2\eta(l_1)$.

Corollary. Let X be a hyperbolic surface, and let γ_1 , γ_2 be two simple closed geodesics with lengths $< ln(3 + 2\sqrt{2})$. Then either $\gamma_1 = \gamma_2$ or $\gamma_1 \cap \gamma_2 = \emptyset$.

Corollary. Let X be a complete hyperbolic surface, γ a simple closed geodesic on X of length l, and A_{γ} the collar around γ . Then any simple geodesic δ on X that enter A_{γ} either intersect γ or spirals towards γ .

Proof. Suppose the geodesic δ enters A_{γ} . We can lift the situation in the universal cover of the hyperbolic disc, where $\tilde{\gamma}$ is a lift of γ can be a diameter of the circle. Then if $\tilde{\delta}$ do not intersect $\tilde{\gamma}$ and do not have the same point at infinity, then its two endpoint are on the same side of $\tilde{\gamma}$ in the disc. Now the translation along $\tilde{\gamma}$ of length $l_{\gamma}(X)$ is in the representation of $\pi_1(X)$. If $\tilde{\delta}$ intersect A_{γ} then by the definition of δ it will intersect with its image by the translation cited before and hence is not simple in X.

4 Mirzakhani's Isomorphism

The aim of this part is to demonstrate the following statement

Theorem 4.1. There is a measurable conjugacy F between the earthquake flow $(\lambda, X) \mapsto (\lambda, E_{t\lambda}(X))$ on $\mathcal{ML} \times \mathcal{T}_g$ and the Teichmüller unipotent flow action of

$$u_t = \begin{pmatrix} 1 & t \\ 0 & 1 \end{pmatrix}$$

on the bundle QD of non-zero quadratic differentials over Teichmüller space \mathcal{T}_g .

$$\mathcal{ML} \times \mathcal{T}_g \xrightarrow{E_t} \mathcal{ML} \times \mathcal{T}_g$$

$$\downarrow F \qquad \qquad \downarrow F$$

$$\mathcal{QD} \xrightarrow{\mathcal{U}} \mathcal{QD}$$

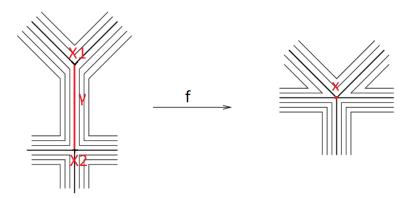
To do this we will decompose F between intermediates maps.

4.1 Tightening map

A first correspondence, found by Thurston, exist between measured foliations and measured laminations. We will mostly follow the paper of Levitt [15].

Definition 4.1. We say that two foliations are equivalent if we can pass from one to the other by Whithead (see definition below) moves or isotopy (homeomorphism isotopic to the identity).

Definition 4.2. Given a measure foliation, an critical segment γ is an arc between two singularities along a leaf which is not a simple closed curve. There is a map f homotopic to the identity that collapse γ to a point x and is identity outside a neighbourhood of γ which contain no other singularity. Doing so we reduce the number of singularities of the foliation and if the extremities of γ are singularities of order k_1 and k_2 , x is now a singularity of the new foliation of order $k_1 + k_2 - 2$. This action is called a Whithead move.



Theorem 4.2. Let X be a closed orientable hyperbolic surface and \mathcal{F} a foliation. There is a canonical geodesic lamination $\gamma(\mathcal{F})$ associated to \mathcal{F} . If \mathcal{F} and \mathcal{F}' are associated foliation then $\gamma(\mathcal{F}) = \gamma(\mathcal{F}')$. In the opposite direction given a geodesic lamination γ , one can find a foliation \mathcal{F} such that $\gamma(\mathcal{F}) = \gamma$ and it's unique up to equivalence.

Definition 4.3. A transverse curve is a simple closed curve C which is never tangent to \mathcal{F} and contain no singularity of \mathcal{F} .

Remark. Since \mathcal{F} contain only saddle singularities, C cannot be contractible, therefore C is isotopic to a simple closed geodesic.

We will work on the universal cover of X, which is the Poincaré disc \mathbb{D} with circle "at infinity" \mathbb{S}_{∞} . We call $p: X \to \mathbb{D}$ the universal projection. $\tilde{\mathcal{F}}$ is $p^{-1}(\mathcal{F})$.

We will say that a foliation follow the (*) condition if the following is true:

If f_1 and f_2 are two compact homotopic leaves then all leaf in the open annulus between them is also compact.

Lemma 4.3. Let h be a leaf of $\tilde{\mathcal{F}}$. Each end of h converge to a point of S; the two point at infinity cannot be the same.

Proof. First, we should notice that the behaviour of leaf at infinity do not change if we take an equivalent foliation. Indeed a homeomorphism ϕ on a compact fundamental domain isotopic to the identity can be extend to an homeomorphism $\tilde{\phi}$ on \mathbb{D} such that $dist(x, \tilde{\phi}(x)) \leq K$. This implies that $\tilde{\phi}$ extend as the identity on the boundary \mathbb{S}_{∞} .

Then given a leaf h of $\tilde{\mathcal{F}}$, we take a half leaf h_0 . If $p(h_0)$ is compact or spiral toward a compact leaf of \mathcal{F} then the first part of the lemma is immediate.

Otherwise, p(h) meet a transverse curve C infinitely often. With an isotopy we can take C to be a geodesic. Now h_0 can meet a connect component of $\tilde{C} = p^{-1}(C)$ only one time. Otherwise there will be a disk bound by an arc of \tilde{C} and an arc of h_0 , which is impossible considering the transversity of C and that \mathcal{F} have no 1-type singularities.

Now every compact of \mathbb{D} meet a finite number of connected components of \tilde{C} so the limit set of h_0 must be on \mathbb{S}_{∞} . This limit set is connected and non empty. Moreover it should not contain any end of a connected component of \tilde{C} . But the ends of connected components of \tilde{C} are dense in \mathbb{S}_{∞} as \tilde{C} is the image of a geodesic by $\pi_1(X)$. This show the first point of the lemma.

The second assertion is clear if p(h) is compact or if it meet a transverse curve C at least twice since then every connected components of \tilde{C} separate the end of h.

Otherwise p(h) spiral toward two compact leaf f_1 and f_2 . If $f_1 = f_2$ and the two endpoints of h are the same then there will be a singularity that would not be a saddle. $f_1 \neq f_2$ is impossible since \mathcal{F} follow the condition (*).

We can now associate to every leaf h a geodesic $\gamma(g)$ by joining the endpoints. Then $\gamma(\tilde{\mathcal{F}}) = \bigcup_{h \in \mathcal{F}} \gamma(h)$ is a disjoint union of geodesic invariant by $\pi_1(X)$. We have to show that this set is closed to conclude that we have a lamination.

Lemma 4.4. $\gamma(\tilde{\mathcal{F}})$ is closed in $\bar{\mathbb{D}}$

Proof. Let $g_n = \gamma(h_n)$ be a sequence of geodesics in $\gamma(\tilde{\mathcal{F}})$ converging to a geodesic g. We want to show $g \in \gamma(\tilde{\mathcal{F}})$. We can suppose that all the g_n are distinct of g and are all on the same side.

Let L be the limit set in $\bar{\mathbb{D}}$. For all leaf m in $\gamma(\tilde{\mathcal{F}})$, we call \bar{m} the closure of m by adding the two end point in \mathbb{S}_{∞} . Then L meet at least one connected component of $\bar{\mathbb{D}}$ \bar{m} . As the end points of all leaf of $\gamma(\tilde{\mathcal{F}})$ is a dense subset of \mathbb{S}_{∞} , L contain a leaf h. Taking a half-leaf h_0 , we want to show that the end point is the same as one of g.

A first case is if there is a simple closed curve C transverse to \mathcal{F} which meet $p(h_0)$ infinitely often. If h_0 does not converge to the corresponding point at infinity then there would be a connected component of $p^{-1}(C)$ that contains the point of infinity of h_0 but does not contain the point of infinity of h_n which is impossible for large n.

A second case is if $p(h_0)$ spirals toward a compact leaf, then closed leaf close to $p(h_0)$ also spirals toward the same compact leaf. Then h_0 converges to one of the points at infinity of g which is a point at infinity of h_n for n large.

Finally if p(h) is compact then $p(h_n)$ spirals toward it for large n, therefore $\gamma(h)$ and g have one point in common at infinity. If the second was different, by applying a transformation leaving $\gamma(h)$ invariant (but no g), we would separate h from the leaves h_n , and it is a contradiction.

Now we want to exhibit an inverse construction which take a lamination λ and give a foliation μ . To do this we still consider $\tilde{\lambda}$ in the universal cover. We will suppose that every complementary region is an ideal polygon.

We can build a skeleton that it composed of edges between vertices and a chosen point in the center. After building the skeleton for every polygon we fill the complementary region. These regions are quadrangles delimited by edges previously built, two of its vertices are on the disk, the remaining ones lye in the center of two adjacent triangles.

This map is often called the "collapsing" map and its inverse the "tightening" map.

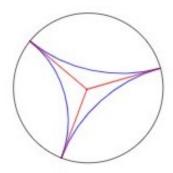


Figure 9: In red the skeleton of the ideal triangle in blue, Image from [25]

The measure we put on this foliation is uniform on every region between four edges of the skeleton.

4.2 Correspondance between foliations and quadratic differentials

For a quadratic differential q, one can define two measured foliations, the horizontal h(q) and the vertical v(q) corresponding in local coordinates to Re(z) and Im(z).

Moreover the measures are

$$h_{\nu} = \int Im(\sqrt{q(z)}dz)$$

and

$$h\mu = \int Re(\sqrt{q(z)}dz)$$

This gives a map to the pair of foliations but it is not surjective, we should restrict it to the image. Define $\Delta = \{(\alpha, \beta) : i(\alpha, \gamma) = i(\beta, \gamma) = 0, \text{ for some } \gamma \in \mathcal{MF}\}$. Δ contains the diagonal (α, α) and is kind of "fat" diagonal.

Lemma 4.5. For any $q \in \mathcal{QD}$, $(h(q), v(q)) \notin \Delta$

Proof. Suppose that there is γ such as $i(h(q), \gamma) = i(v(q), \gamma) = 0$ for some $\gamma \in \mathcal{ML}$. Let's take a sequence of simple closed weighted curves γ_i converging to γ . By continuity of the intersection number we have that $i(h(q), \gamma_i) \to 0$. So there is a sequence of saddle connection in the same homotopy class as γ_i whose x-component is very small. We can have the same argument in the vertical direction and find a contradiction.

Theorem 4.6. The map $q \mapsto (h(q), v(q))$ defines a homeomorphism $\mathcal{QD} \to \mathcal{MF} \times \mathcal{MF} \setminus \Delta$

Proof. We can describe the inverse map. If we take two measured foliation h and v we can tighten them into laminations (which we also call h and v) as in the previous section. These two laminations do not share any leaf, otherwise we would have $(h, v) \in \Delta$ by considering the leaf as a geodesic. The complementary region of $h \cup v$ are compact polygons, i.e. they do not have a vertex on the boundary of the disk. Now we can fill the polygon to obtain the quadratic differential with a singularity in each complementary region of order the number of side of the polygon.

We will note the inverse map q(.,.)

More details on the proof can be found in [5] Lemma 6.2 and [10]. The injectivity is discussed with more details in [8] section 3.

4.3 Shear Coordinate

Finally there is a map that, given a hyperbolic structure X and a lamination λ creates a measured foliation which is transverse to λ . We will note this map F_{λ}

To simplify we will impose for λ to be a maximal lamination i.e. if $\tilde{\lambda}$ is the pre-image of λ in the universal cover \mathbb{D} , $\mathbb{D}\setminus\tilde{\lambda}$ is made of ideal triangle. We will first work in this triangle minus a region on the center, then give a measure in this foliation, and finally show that it is a homeomorphism.

So in one ideal triangle, given two sides we can draw an arc perpendicular with these two sides, which is the intersection of the ideal triangle with the circle tangent to the boundary of \mathbb{D} in the common extremity the two arc chosen. Then by rescaling by a factor $r \in [0;1]$ and doing the same procedure for the two other pairs we get a foliation in the ideal triangle minus a locus in the center.

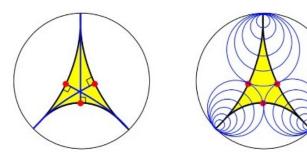


Figure 10: Image from [17]

Then we can find a full foliation by pitching the center which is not cover.

We have a natural transverse measure to this foliation. For an arc in one ideal triangle of the lamination we can project it along the leaf of the foliation to a segment in the edge of the ideal triangle, the length of the arc will be the length of the segment. As the leaves are horocycle circles based on the same point, it will not depend of which side we choose to project. Then given an arbitrarily arc we decompose it along the ideal triangle it meets.

We want to show that this construction is reversible, which means that given $\mu \in \mathcal{MF}_{\lambda}$, the set of foliation transverse to the lamination, we can construct $X \in \mathcal{T}_g$ whose horocycle foliation is μ .

The idea is that, if we already know X, the lamination λ can be lifted to $\tilde{\lambda}$ which is invariant of Γ the Fushian group of X. But we can build $\tilde{\lambda}$ only with the information given by μ .

We will note $\tilde{\mu}$ the pre-image of μ in \mathbb{D} . We consider two triangles T_1 and T_2 that are complementary regions of $\tilde{\lambda}$ and we take a segment A in a leaf of $\tilde{\mu}$ that goes to an edge of T_1 to T_2 . We name v_1 and v_2 the two vectors with footprint in the edge and tangent to them. Then there is a Moebius transformation S which take one to the other. With one more information, the "shear", we can place T_2 on \mathbb{D} , according to the position of T_1

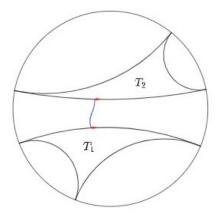


Figure 11: Image from [25]

Indeed given only a Moebius transformation we still have a one parameter family of triangles T_2 with an edge generated by v_2 . To fix this we trace two orthogeodesics coming from the vertex of the ideal triangles facing the considered edges and from the point of intersection in T_1 we follow a leaf of the foliation, then, when we meet T_2 we have to move along the geodesic edge to find the other point of intersection. This length is the shear between T_1 and T_2

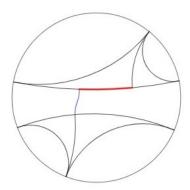


Figure 12: Image from [25]

Let I be the set of all triangles in \mathbb{D} that A meets. For each $i \in I$ we can define v_i^+ and v_i^- the vectors tangent to the edge of the corresponding triangle at the intersection of the edges and A. Note that I is a countable totally ordered but non well ordered set. So if we take S_i the Moebius transformation which takes v_i^- to v_i^+ , we have to give a meaning to the expression

$$\prod_{i \in I} S_i$$

Definition 4.4. Given a countable totally ordered set of indices I and elements S_i in a Banach algebra, we say that $\prod_i S_i$ is well defined and equal to S if for any increasing chain

$$I_0 \subset I_1 \subset ... \subset I$$

with $\bigcup_k I_k = I$ we have $\lim_{k \to \infty} \prod_{i \in I_k} S_i = S$.

Lemma 4.7. For element s_i in a Banach algebra indexed by a countable totally ordered set, if $\sum ||s_i|| < \infty$, then $\prod (1 + s_i)$ is well-defined.

Proof. For $1 \leq m \leq n$, we have

$$\left\| \prod_{i=1}^{n} (1+s_i) - \prod_{i=1, i \neq m}^{n} (1+s_i) \right\| \le \|s_m\| \left\| \prod_{i=1, i \neq m}^{n} (1+s_i) \right\| \le \|s_m\| \prod_{i=1}^{n} (1+\|s_i\|)$$

Or with the assumption $\sum ||s_i|| < \infty$ we have that $\prod (1 + ||s_i||) \le C < \infty$, so removing or adding $1 + s_m$ produce a change bound by $||s_m||C$.

Now we want to apply this lemma to $S_i - Id$, with Id being the identity matrix.

Lemma 4.8. For the previous S_i , if we note $s_i = S_i - Id$ we have $\sum ||s_i|| < \infty$.

Proof. Each S_i is conjugate to a horocycle transformation of time one. The conjugacy is made by a geodesic flow along the edges of the triangle. We can compute

$$\begin{pmatrix} e^{-t/2} & 0 \\ 0 & e^{t/2} \end{pmatrix} \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix} \begin{pmatrix} e^{t/2} & 0 \\ 0 & e^{-t/2} \end{pmatrix} = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} + \begin{pmatrix} 0 & e^{-t} \\ 0 & 0 \end{pmatrix}$$

So the norm of s_i is inversely correlated to the amount of geodesic flow used in the conjugation. Now we can partition the indices set I into finitely many subsets (I_k) according to which spike of the lamination the arc of A crosses. Then for a spike the sum of $||s_i||$ with $i \in I_k$ is finite, indeed the distance between two neighbouring crossing is bounded below by a constant and so the amount of time we follow the geodesic flow increases at least linearly. Finally the norm of the s_i should decrease geometrically.

So we can conclude that there is an unique Moebius transformation S equal to the meaningful expression $\prod_i S_i$.

Now we can conclude the proof. There exists, without any hyperbolic structure X given, topological classes for $\tilde{\mu}$ and $\tilde{\lambda}$. We choose one arbitrary ideal triangle T_1 in the lamination. For every other triangle T_2 , the Moebius transformation and the shear are data that can be computed only using the transverse measure of $\tilde{\mu}$. So we can place T_2 , and the other triangle. The closure of this set gives the lamination $\tilde{\lambda}$. $\tilde{\lambda}$ will be preserved by a Fushian group Γ and we will have $X = \mathbb{D}/\Gamma$.

4.4 Conjugation between the horocyclic flow and the earthquake flow

We now have to prove that the map between $\mathcal{ML} \times \mathcal{T}_g$ and \mathcal{QD} conjugate the earthquake flow and the horocyclic flow.

Lemma 4.9. Denote by $Shear_X(T_1, T_2)$ the shear for two triangles joined by an arc A of the horocyclic foliation on the hyperbolic surface X. Then

$$Shear_{E_{t\lambda}(X)}(T_1, T_2) = Shear_X(T_1, T_2) + t\lambda(A)$$

where $\lambda(A)$ denote the transverse measure of A and t is sufficiently small.

Proof. T_1 and T_2 are separated by infinitely many leaves of $\tilde{\lambda}$. We want to understand how T_2 moved relatively to T_1 by the action of the earthquake $E_{t\lambda}$. We can approximate the measured lamination between the two triangles by a discrete one. If we follow the earthquake path along a leaf γ of $\tilde{\lambda}$ between T_1 and T_2 by an amount t, then this changes the shear coordinate between T_1 and T_2 by t. Indeed each arc of the horocycle foliation which have endpoint on γ see their endpoints translated by t. Similarly if the earthquake moves finitely many leaves of λ with measure a_i , the shear changes by precisely $t \sum a_i$. So taking a limit, we have the lemma for an arbitrary measured foliation.

We will now show that the Mirzakhani's map conjugates the earthquake flow to the horocyclic flow. Let \mathcal{ML}_0 denote the set of maximal lamination and \mathcal{QD}_0 the locus of quadratic differentials with simple zeros and no horizontal saddle connection. We begin with an easy lemma.

Lemma 4.10. An arc joining two singularities of a quadratic differential have in its isotopy class a path made of horizontal and vertical arcs between two singularities of the quadratic differential.

Corollary. Suppose q_t is a path in the space of quadratic differential such that for every t_0 and every path γ on q_{t_0} joining two singularities, the period $x_t + iy_t$ of γ satisfies

$$\frac{d}{dt}|_{t=t_0}x_t = y_{t_0}, \frac{d}{dt}|_{t=t_0}y_t = 0$$

Then q_t is an orbit of the horocyclic flow.

Proof. We have

$$\begin{pmatrix} 1 & t \\ 0 & 1 \end{pmatrix} \begin{pmatrix} x \\ y \end{pmatrix} = \begin{pmatrix} x + ty \\ y \end{pmatrix}$$

And we can integrate the two equations to find the linear action.

We can now prove theorem 4.1.

Proof. We want to show that

$$q(\lambda, F_{\lambda}(E_{t\lambda}(X)))$$

is a horocyclic flow path. Here λ is firstly seen as a foliation, and then as a lamination.

We pick an arbitrary time t_0 and look at the derivative of the path, called γ . The coordinate y_t is constant equal to $\lambda(\gamma)$ and the derivative of x_t is by lemma 4.9 equal to y_t . We conclude using the previous lemma.

5 Mixing rate

5.1 Mixing properties of the geodesic and horocyclic flows

The geodesic flow and the horocyclic flow can act in two situation. In the first one they belong to the action of $SL_2(\mathbb{R})$ on the quotient of the upper half plane by a lattice seen as a finite volume hyperbolic surface. A second manner is in the Teichmuller space and then in the modular space of a given surface.

We will first give the behaviours of two flows we described before.

Theorem 5.1. The geodesic flow and the horocycle flow are mixing as seen as flow on surface and on modular space.

We will just give the proof in the case of a finite volume hyperbolic surface. The detail of the proof can be found in this course [18].

Proof. The proof will be in four step, first we will show that the geodesic flow is ergodic, the horocyclic flow is ergodic, the geodesic flow is mixing, finally the horocycle flow is mixing.

Geodesic flow is ergodic

The following proof is called the Hopf's argument [9] and is now common in ergodic theorie. We will look at the time average of a function $f \in L^2$ which is continuous and compactly support. We will study only the case where the surface is compact, hence f is uniformusly continuous. It will be sufficient since this space is dense. We note:

$$F(x) = \lim_{T \to \infty} \frac{1}{T} \int_0^T f(g_t x) dt$$

We want to show that F is almost everywhere constant.

For a point x in the surface, F(x) is constant a.e. on the stable foliation $W^s(x)$ since if two geodesic with the same forward endpoint are asymptotic we can control the difference of the integral with the uniformous continuity of f. By the same argument F is constant a.e. on the unstable foliation $W^u(x)$. By Birkhoff's theorem this two constant are the same almost everywhere.

We should now add that the measure is absolutly continuous with respect to the stable foliation and the unstable foliation, the definition is given below.

Definition 5.1. Let X be a metric space, f_t of borelian flow whith an invariant borelian measure μ of finite mass. We say that μ is absolutly continous with respect to the stable foliation and unstable foliation if for every $v \in X$ we can find a open neighborhood U of v, a $\delta > 0$ and a homeomerphism $\phi : \mathbb{R} \times \mathbb{R} \to U$ such has for every $(x, y) \in \mathbb{R} \times \mathbb{R}$,

- 1. $\phi(\{x\} \times \mathbb{R}) \subset W^s(\phi(x,y))$
- 2. $\phi(\mathbb{R} \times \{y\}) \subset f_{|]-\delta,\delta[}W^u(\phi(x,y))$
- 3. $\phi^*(\mu_{|U})$ is equivalent to a product measure

We will admit that this is the case in our situation. Moreover we can conclude with the following lemma.

Lemma 5.2. Let (X, \mathcal{T}, μ) and (Y, \mathcal{S}, ν) be two probability spaces and $f: X \times Y \to \mathbb{R}$ a function in L^2 . We suppose that there is $\phi_1: X \to \mathbb{R}$ and $\phi_2: Y \to \mathbb{R}$ two measurable function, such has for $Z \subset X \times Y$ a subset of full $\mu \otimes \nu$ we have:

$$\forall (x,y) \in Z, f(x,y) = \phi_1(x), f(x,y) = \phi_2(y)$$

Then f is constant almost everywhere.

Proof. According to Fubini's theorem, there is $Y_0 \subset Y$ of full measure and $x_0 \in X$ such has $\{x_0\} \times Y_0 \subset Z$. For all $(x,y) \in Z \cap (X \times Y_0)$, the point (x_0,y) is in Z so we have

$$\phi_1(x_0) = f(x_0, y) = \phi_2(y) = f(x, y)$$

which finish the lemma.

Putting the pieces together we have that the geodesic flow is ergodic.

Ergodicity of the horocycle flow

Now if we take $f \in L^2(X)$ a function invariant under the horocycle flow and of mean zero, we want to show that f = 0 almost everywhere. Let G^t , H^s and E^r correspond to the operators of the different flows. We have the relation

$$H^s = E^r G^t E^{\pi + r}$$

where $r \to 0$ as $t \to \infty$. Since $H^s f = f$, we have for any T > 0,

$$f = \frac{1}{T} \int_0^T E^r G^t E^{\pi + r} f dt$$

Thus for every $g \in L^2(X)$ we have

$$< g, f > = \int_{X} g \frac{1}{T} \int_{0}^{T} E^{r} G^{t} E^{\pi + r} f dt$$

As $r \to 0$ then $t \to \infty$ we can show by controlling the difference that

$$\langle g, f \rangle = \lim_{t \to \infty} \int_X g \frac{1}{T} \int_0^T G^t E^{\pi} f dt$$

Then as we have shown that the geodesic flow is ergodic, we have by Von Neumann ergodic theorem

$$\langle g, f \rangle = \langle g, \int_{X} E^{\pi} f \rangle = 0$$

This conclude the proof that the horocycle flow is ergodic.

Mixing of the geodesic flow

We have the relation

$$h_s g_t = g_t h_{se^{2t}}$$

Let us take $f_0, f_1 \in C_0(X)$, we have for small s

$$< f_0, g_t f_1 > \simeq < h_{-s} f_0, g_t f_1 > = < f_0, h_s g_t f_1 > = < f_0, g_t h_{se^{2t}} f_1 > = < g_{-t} f_0, h_{se^{2t}} f_1 > = < g_{-t} f_0, h_{se^{2t}} f_1 > = < f_0, g_t h_{se^$$

Now we have for small s,

$$h_{Se^{2t}}f_1 \simeq \frac{1}{S} \int_0^S h_{e^{2t}s} f_1 ds = F_t$$

and for large t, as the horoclic flow is ergodic

$$F_t = \frac{1}{S} \int_0^S h_{e^{2t}s} f_1 ds \simeq \int_X f_1 = \langle f_1, 1 \rangle$$

So we have

$$< f_0, g_t f_1 > \simeq < g_{-t} f_0, F_t > \simeq < g_{-t} f_0, 1 > < f_1, 1 > = < f_0, 1 > < f_1, 1 >$$

Which is the mixing of the geodesic flow.

Mixing of the horocyclic flow

We use again the relation $h^s = e_r g_t e_{\pi+r}$.

$$< h_s f_0, f_1 > = < g_t e_{\pi+r} f_0, e_{-r} f_1 > \simeq < g_t e_{\pi} f_0, f_1 >$$

for t large (and so r small). By the mixing property of the geodesic flow, this quantity converges to $< f_0, 1 > < f_1, 1 >$, and the horocyclic flow is mixing.

Then we have this elementary corollary

Corollary. The Earthquake flow is also mixing with respect to the Liouville measure.

Proof. With Mirzakhani's conjugacy, the earthquake flow is conjugated to the horocycle flow which is mixing. This property is carried over. \Box

5.2 Rate of mixing of this flows

Then we want to know at which rate the mixing of the geodesic and horocyclic flow happens. Ratner showed in 1986 [23] that the geodesic flow is exponentially mixing and the horocyclic flow has a polynomial rate of mixing. She used representation theory.

We will first give some definitions before getting to the main theorem.

Definition 5.2. Let H be a complex separable Hilbert space, U(H) the group of all unitary transformation of H onto itself, G a Lie group and $T: G \to U(H)$ a unitary representation. We note $T(g) = T_g \in U(H)$.

An element $v \in H$ is called a C^k -vector for T, $k = 0, 1, ..., \infty$ if $g \mapsto T_g(v)$ is a C^k -map from G to the Hilbert space. We will note K(T, p) the space of C^k -vector

Remark. The space of C^{∞} vector is dense in H.

Definition 5.3. If v is a C^1 -vector of H and X is in the Lie algebra of G, the Lie derivative $L_X v$ is

$$L_X v = \lim_{t \to 0} \frac{T(exptX)v - v}{t}$$

Now if Γ is a Fushian group, $SL_2(\mathbb{R})$ actss on the hyperbolic surface with its measure (X, μ) as seen as \mathbb{H}/Γ . Then it will also act by a representation ρ on $L_0^2(X, \mu)$ the space of zero average function on the surface.

We want to study the decorrelation induced by some subgroup of $SL_2(\mathbb{R})$. It will link with the following quantity.

Definition 5.4. Let ϕ and ψ be zero-mean functions in L^2 , the matrix coefficient $C_{\phi,\psi}$ is

$$g \mapsto | <\phi, \rho(g)\psi > |$$

Remark. As the horocyclic flow and the geodesic flow are mixing we have $C_{\phi,\psi}(e_t) \to 0$ and $C_{\phi,\psi}(h_s) \to 0$ for all functions ϕ and ψ .

We can decompose this representation as an integral of irreducible representations

$$L_0^2(X,\mu) = \int H_{\zeta} d\nu(\zeta)$$

With ρ_{ζ} the representation of H_{ζ} .

This set of representation decompose in three parts.

- 1. The principal serie
- 2. The discrete serie
- 3. The complementary serie

This decomposition is given by the spectrum of the Casimir operator on each H_{ζ} . The discrete part of the spectrum is the discrete series, and if q is the eigenvalue, the part $q \geq 1/4$ is the principal series and q < 1/4 the complementary series. This decomposition was studied by Bargmann in [2].

So we can decompose the measure $\nu = \nu_p + \nu_d + \nu_c$.

Moreover we know that if ρ_{ζ} is in the complementary series, there exist $s = s(\zeta) \in]0;1[$ such that the representation is isomorphic to the representation π_s on the Hilbert space

$$\mathcal{H}_s = \{ f : \mathbb{R} \to \mathbb{C}, \|f\|^2 = \int_{\mathbb{R} \times \mathbb{R}} \frac{f(x)\overline{f(y)}}{|x - y|^{1 - s}} dx dy < \infty \}$$

With the action

$$\pi_s \begin{pmatrix} a & b \\ c & d \end{pmatrix} f(x) = \frac{1}{(cx+d)^{1+s}} f(\frac{ax+b}{cx+d})$$

The representation ρ is isolated from the trivial representation if and only if there exist $\epsilon > 0$ such that $s(\zeta) < 1 - \epsilon$ for ν_c almost every ζ .

So in our case, the Lie algebra of $SL(2,\mathbb{R})$, is the vector space of 2×2 matrix with zero trace. A basis for this space is

$$W = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}, Z = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}, V = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}$$

Which gives the Casimir operator on C^2 vectors of $SL(2,\mathbb{R})$,

$$\Omega_T = (L_v^2 + L_Z^2 - L_W^2)/4$$

On a irreducible part of T, Ω_T is a scalar multiple of the identity, i.e.

$$\Omega_T v = \lambda v$$

for some $\lambda = \lambda(T) \in \mathbb{R}$ and all C^2 -vectors $v \in H(T)$

If Γ is a lattice of $SL(2,\mathbb{R})$, Ω_T create an operator Δ on $\Gamma \backslash SL(2,\mathbb{R})$. We call $\Lambda(\Delta)$ its spectrum and

1.
$$A(\Gamma) = \Lambda(\Gamma) \cup] - 1/4; 0[$$

2.
$$B(\Gamma) = \sup A(\Gamma)$$

3.
$$C(\Gamma) = -1 + \sqrt{1 + 4B(\Gamma)}$$

We can now write the theorem about the decay of correlation.

Theorem 5.3. Let Γ be a lattice in $SL(2,\mathbb{R})$, $M = \Gamma \backslash SL(2,\mathbb{R})$ and T be the regular representation of $SL(2,\mathbb{R})$ on $L^2(M,\mu)$. Let $v,w \in K(T,p)$, p > 0, < w, 1 >= 0 and $B(t) = < v, w \circ g_t >$, $C(t) = < v, w \circ h_t >$. Then there exist $t_0 = t_0(\Gamma) > 0$ such that for all $|t| \ge t_0$ and some E, F > 0

1.
$$|B(t)| < E(b(|t|))^{\alpha(p)}$$

2.
$$|C(t)| \leq F(b(\ln|t|))^{\alpha(p)}$$

where

1.
$$b(t) = e^{\sigma(\Gamma)}$$
 if $A(\Gamma) \neq \emptyset$

2.
$$b(t) = e^{-t}$$
 if $A(\Gamma) = \emptyset$, $\sup(\Lambda \cap]-\infty; -1/4[) < -1/4$ and $-1/4$ is not an eigenvalue of Ω

3.
$$b(t) = te^{-t}$$
 if $A(\Gamma) = \emptyset$, $\sup(\Lambda \cap] - \infty; -1/4[) = -1/4$ or $-1/4$ is an eigenvalue of Ω

and $\alpha(p)$ is

1. 1 if
$$p \ge 3$$

2.
$$\frac{2p}{2p+1}$$
 if $2 \le p < 3$

3.
$$\frac{2p}{2p+3}$$
 if $1 \le p < 2$

4.
$$\frac{p}{n+3}$$
 if 0

Remark. We are only interested in functions in $L^2(X,\mu)$ which depend only on the footprint of the tangent space and are constant on the tangent space. This implies that they are all in $K(T,\infty)$.

On the other we can make the path backward, by learning information on the representation via the mixing rate of the flow it generates. This way is taken in the appendix B of [1].

Definition 5.5. Let G be a locally compact σ -compact group. A continuous unitary representation of G is said to have almost invariant vectors if for every $\epsilon > 0$ and for every compact subset $K \subset G$, there exists a unit vector V such that $||g * v - v|| < \epsilon$ for all $g \in K$.

A unitary action which does not have almost invariant vectors is said to be *isolated from* trivial representation.

If G is a semi-simple Lie group, a representation which is isolated from trivial representation is also said to have a spectral gap.

With this definition, we can write the following theorem.

Proposition 5.4. Let consider a representation of $SL(2,\mathbb{R})$ by a measure preserving action of automorphisms of a probability space. Let ρ be the representation associate on H the space of the L^2 function with zero averages. Assume there is $\delta \in]0;1[$ and a dense subset of the subspace of $SO(2,\mathbb{R})$ -invariant function $H' \subset H$ consisting of functions ϕ for which the correlations $\langle \phi, \rho(g_t) \rangle$, $g_t = \begin{pmatrix} e^t & 0 \\ 0 & e^{-t} \end{pmatrix}$, decays like $O(e^{-\delta t})$. Then ρ is isolated from the trivial representation.

6 Example of the once punctured torus

We will now study the simplest hyperbolic surface: the once punctured torus. A usual torus can not have a hyperbolic structure, it has naturally a flat structure as the quotient $\mathbb{R}^2/\mathbb{Z}^2$. But this changes when one studies the once punctured torus. It is a torus S where we choose a point p and remove it (or just mark it).

The construction of this object can be done in two ways at least. For the first construction, we give a fundamental domain in the upper half plane. It is a hexagon with two vertices "at infinity" on the real line. We glue each arc to its symmetric.

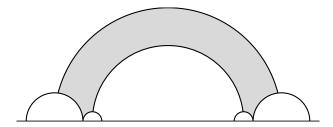


Figure 13: A fundamental domain of a punctured torus

A second construction is given by the representation of the Fushian group. Given two hyperbolic isomorphism of \mathbb{H} A and B with different fixed point on $\delta \mathbb{H}$. We ask in addition that the commutator $H := ABA^{-1}B^{-1}$ should be a parabolic element.

6.1 First properties

The geometry of the one punctured torus makes it simple to discuss closed curves and laminations.

Remark. Given two generators α and β of $\pi_1(S)$, two closed curves non homotopically trivial which intersect once, one can parametrize all other lamination. Indeed a given lamination $\lambda \in \mathcal{ML}$ is determined by the couple $(i(\alpha, \lambda), i(\beta, \lambda))$ where i(., .) is the geometric intersection number.

The collaring is a powerful tool in this situation since two closed curves can not be disjoint. The length of one is bound below by the collar of the other. So we have the following lemma:

Lemma 6.1. If X is a hyperbolic surface homeomorph to the once punctured torus with a simple closed curve of length inferior to $\ln(3+2\sqrt{2})$, then there is only one simple closed curve representing the systole.

Moreover, we have this useful lemma to estimate the length of the systole function in the Teichmuller space.

Lemma 6.2. Pick γ a simple closed geodesic, and $X \in \mathcal{T}(S_{1,1})$, if X has Fenchel-Nielsen coordinate $(L, \frac{p}{q})$ with respect to γ , where gcd(p,q) = 1 and $\frac{p}{q} \in]0; 1[$, then

$$C_1(L)e^{\frac{-L}{2q}} < l_{sus}(X) < C_2(L)e^{\frac{-L}{2q}}$$

where $C_1(L)$, $C_2(L)$ both limits to 4 when p, q are fixed and L goes to ∞

Proof. Let R(L) be the length of the shortest geodesic arc with endpoints on γ . We have

$$R(L) = 2log(coth(L/4)) = 2log(\frac{e^{l/2} + 1}{e^{l/2} - 1})$$

By the collar lemma 3.10. Then if we take α a simple closed curve that intersect γ , q times exactly, we obtain the following inequality with $a = l_{\alpha}(X)$

$$qR(a) < L < qR(a) + \frac{qa}{2}$$

Reorganizing the terms we have

$$e^{-a/4}tanh(a/4) < e^{-L/2q} < tanh(a/4)$$

As $a \to 0$ then $L \to \infty$

$$C_1(L)e^{\frac{-L}{2q}} < a < C_2(L)e^{\frac{-L}{2q}}$$

To complete the proof, we need to show that the length of α is shorter than any other geodesic closed curves, but with the collar lemma there is only one systole whose length goes to 0 and this is the case for α as $L \to \infty$.

6.2 Rate of mixing of the Earthquake flow

If ρ is a positive real number we will note $J(\rho) = \{t \in \mathbb{R}, l_{\gamma}(t) \leq \rho\}$, where $l_{\lambda,\gamma}(t) = l_{\gamma}(E_t(X,\lambda))$. To shorten the notation we will use $\mu = i(\lambda_p, \gamma)$ and $\epsilon_{\gamma} = \min_{t \in \mathbb{R}} l_{\gamma}(t)$. Kerckoff show in [13] the following properties of the earthquake flow.

Proposition 6.3. For every $(x, \lambda) \in \mathcal{PT}^1$ and every $\gamma \in \Gamma^S$, the function $l_{\lambda,\gamma} : t \mapsto l_{\gamma}(E_t(x, \lambda))$ is C^1 and convex.

If $i(\gamma, \lambda) = 0$ then $l_{\lambda,\gamma}$ is constant, and if $i(\gamma, \lambda) > 0$ then $l_{\lambda,\gamma}$ is proper, with derivative strictly increasing and given by

$$\frac{d}{dt}l_{\lambda,\gamma}(t) = \int_{\gamma_t} \cos(\theta_{y,t}) d\lambda(y) \tag{1}$$

Here γ_t is the geodesic representative of γ in $E_t(X,\lambda)$, and $\theta_{\gamma,t}$ is the angle of the intersection between λ and γ_t .

With this proposition we can easily show that

$$-i(\lambda_p, \gamma) < \frac{d}{dt}l_{p,\gamma}(t) < i(\lambda_p, \gamma)$$

The following lemma can be found in [21].

Lemma 6.4. There are constants ρ and C, depending only on S, such that for any $(X, \lambda) \in \mathcal{PT}^1, \gamma \in \Gamma^S$ and all $t \in J(\rho)$,

$$\mu |t-t_\gamma| - C\epsilon_\gamma \le l_{\lambda,\gamma}(t) \le \mu |t-t_\gamma| + \epsilon_\gamma$$

We now fix a ϵ , such has $2\epsilon(2+C) < min(\rho, ln(3+2\sqrt{2}))$ and a μ . Let

$$C_{\epsilon,\mu} = \{(X,\lambda) \in \mathcal{T}(S) \times \mathcal{ML}, l_{sys}(X) \in [\frac{\epsilon}{2}; \epsilon], i(\lambda,\gamma) \in [-\mu,\mu]\}$$

where γ is the closed curve representing the systole. Moreover, let D be the set $\{(X, \lambda) \in \mathcal{T}(S) \times \mathcal{ML}, l_{sys}(X) > ln(3 + 2\sqrt{2})\}.$

So let X be a hyperbolic surface in $C_{\epsilon,\mu}$. The systole of this surface is γ and has length $\tilde{\epsilon}$ When we earthquake X along λ , the length function will reach a minimum value ϵ_{γ} at the time t_{γ} . Using the inequalities of lemma 6.4, we find that

$$\frac{\tilde{\epsilon} - \epsilon_{\gamma}}{\mu} \le |t_{\gamma}| \le \frac{\tilde{\epsilon} + C\epsilon_{\gamma}}{\mu}$$

and as $0 < \epsilon_{\gamma} \le \tilde{\epsilon}$ we have $|t_{\gamma}| \in [0; \frac{\epsilon}{\mu}(1+C)]$.

Injecting this in the right-hand side of inequality of lemma 6.4, we get:

$$l_{\gamma}(E_t(X,\lambda)) \le \tilde{\epsilon}(2+C) + \mu|t|$$

Hence, for all $|t| \leq \frac{\epsilon}{\mu}(2+C) = t_{lim}$ we have that $l_{\gamma}(E_t(X,\lambda)) \leq 2\epsilon(2+C)$.

Let $f_{\epsilon,\mu}$ a Lipschitz positive function with support include in $C_{\epsilon,\mu}$, $\int f_{\epsilon,\mu} > Vol(C_{\epsilon,\mu})/2$. We take also g a positive Lipschitz function with support in D, $K_1 = \int g$ and $||g||_{Lip} = K_2$. So we have:

$$\left| \int f_{\epsilon,\mu} \circ E_{t_{lim}} g - \int f_{\epsilon,\mu} \int g \right| = K_1 \int f_{\epsilon,\mu}$$

If we have exponential mixing, we would have two constants τ' and K_3 such that

$$\frac{Vol(C_{\epsilon,\mu})}{2}K_1 \le K_1 \int f_{\epsilon,\mu} < K_2 K_3 ||f_{\epsilon,\mu}||_{Lip} e^{-\epsilon \frac{2+C}{\mu\tau'}}$$

Or more simply

$$Vol(C_{\epsilon,\mu}) < K' \| f_{\epsilon,\mu} \|_{Lip} e^{-\frac{\epsilon}{\mu\tau}}$$
 (2)

We can also make the hypothesis that the earthquake flow has a polynomial decay rate. We change the inequality 2 in

$$Vol(C_{\epsilon,\mu}) < K' \| f_{\epsilon,\mu} \|_{Lip} \left(\frac{\epsilon}{\mu} \right)^{-\gamma}$$
(3)

Where $\gamma \in \mathbb{R}^{+*}$.

We should now understand the volume of $C_{\epsilon,\mu}$ and the Lipschitz norm of $f_{\epsilon,\mu}$. Note that we are interested in an upper bound of the norm and a lower bound of the volume.

One can choose $f_{\epsilon,\mu}$ to be a product of piecewise linear functions as follow:

$$f_{\epsilon,\mu} = g_{\epsilon} h_{\epsilon} j_{\epsilon,\mu}$$

With

$$g_{\epsilon} = \begin{cases} 0 & \text{if } l_{sys} \notin [\epsilon/2; \epsilon] \\ 1 & \text{if } l_{sys} \in \left[\frac{4\epsilon}{6}; \frac{5\epsilon}{6}\right] \\ \frac{6}{\epsilon}(l_{sys}(X) - \frac{\epsilon}{2}) & \text{if } l_{sys} \in \left[\frac{\epsilon}{2}; \frac{4\epsilon}{6}\right] \\ \frac{-6}{\epsilon}(l_{sys}(X) - \epsilon) & \text{if } l_{sys} \in \left[\frac{5\epsilon}{6}; \epsilon\right] \end{cases}$$

$$h_{\lambda} = \begin{cases} 0 & \text{if } l_{\lambda}(X) > 1 \\ 1 & \text{if } l_{\lambda}(X) \leq \frac{1}{2} \\ -2(l_{\lambda}(X - 1)) & \text{if } l_{\lambda}(X) \in \left[\frac{1}{2}; 1\right] \end{cases}$$

$$j_{\epsilon,\mu} = \begin{cases} 0 & \text{if } i(\lambda, \gamma) > \mu \\ 1 & \text{if } i(\lambda, \gamma) \leq \mu/2 \\ \frac{-2}{\mu}(i(\lambda, \gamma) - \mu) & \text{if } i(\lambda, \gamma) \in \left[\mu/2; \mu\right] \end{cases}$$

Then we need to calculate the different Lipschitz norms of this function.

To estimate the volume of $C_{\epsilon,\mu}$, and following [7], we use the fact that the measure ν restricted to the space of unit-length measured laminations is equal to the Thurston volume of the unit ball.

When we will have the two previous estimation, we would know if the earthquake flow has an exponential mixing rate or a polynomial rate and a bound of the degree of the polynomial in the last case.

7 Conclusion

In this master thesis, I wanted to show how the hyperbolicity provides us with powerful tools. These tools, among which the collaring lemma or the Fenchel-Nielsen coordinates can be cited, are a cornerstone of understanding flows on the underlying surface. The objects introduced (lamination, foliation, quadratic differential) are diverse yet all linked and interconnected. This strong consistency is key to transporting qualitative properties as we did in the case of the ergodicity between the horocyclic flow and the earthquake flow. But some quantitative properties are still unknown and require more efforts. Although we can compute decay rates of geodesic and horocyclic flows, we can not directly find this rate for the earthquake flow. So far, we're still working on deriving a bound in the simplest possible case, which is the once punctured torus.

In the future, efforts will be set on finishing the derivations and computations of the mixing rate of the earthquake flow on the once punctured torus. Then we hope that this simple case will yield enlightening result and allow us to apply similar arguments in higher genus. This will provide us with an upper bound of the rate of mixing. However, we expect that showing that the bounds obtained are tight and examining in which situations they can be reached will require us to use other approaches and tools.

I am very enthusiastic about these perspectives and look forward to working on them in my Ph.D under the supervision of Carlos Matheus.

I would like to thanks Carlos Matheus who shared his knowledge and guided me on this topic. I am also grateful to Ser-Wei Fu, who helped us a lot to understand the bound we can give to the length of systole along the earthquake flow. My thanks also go to Ming Kun Liu who helped me getting a hand on some tricky points, and to Athénaïs Gautier who corrected some of my mistakes.

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