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# Finitely-additive measures on the asymptotic foliations of a Markov compactum. 

Alexander I. Bufetov

## 1 Introduction.

### 1.1 Hölder cocycles over translation flows.

Let $\rho \geq 2$ be an integer, let $M$ be a compact orientable surface of genus $\rho$, and let $\omega$ be a holomorphic one-form on $M$. Denote by $\mathfrak{m}=(\omega \wedge \bar{\omega}) / 2 i$ the area form induced by $\omega$ and assume that $\mathfrak{m}(M)=1$.

Let $h_{t}^{+}$be the vertical flow on $M$ (i.e., the flow corresponding to $\Re(\omega)$ ); let $h_{t}^{-}$be the horizontal flow on $M$ (i.e., the flow corresponding to $\Im(\omega)$ ). The flows $h_{t}^{+}, h_{t}^{-}$preserve the area $\mathfrak{m}$ and are uniquely ergodic.

Take $x \in M, t_{1}, t_{2} \in \mathbb{R}_{+}$and assume that the closure of the set

$$
\begin{equation*}
\left\{h_{\tau_{1}}^{+} h_{\tau_{2}}^{-} x, 0 \leq \tau_{1}<t_{1}, 0 \leq \tau_{2}<t_{2}\right\} \tag{1}
\end{equation*}
$$

does not contain zeros of the form $\omega$. Then the set (1) is called an admissible rectangle and denoted $\Pi\left(x, t_{1}, t_{2}\right)$. Let $\overline{\mathfrak{C}}$ be the semi-ring of admissible rectangles.

Consider the linear space $\mathcal{Y}^{+}$of Hölder cocyles $\Phi^{+}(x, t)$ over the vertical flow $h_{t}^{+}$which are invariant under horizontal holonomy. More precisely, a function $\Phi^{+}(x, t): M \times \mathbb{R} \rightarrow \mathbb{C}$ belongs to the space $\mathcal{Y}^{+}$if it satisfies:

1. $\Phi^{+}(x, t+s)=\Phi^{+}(x, t)+\Phi^{+}\left(h_{t}^{+} x, s\right)$;
2. There exists $t_{0}>0, \theta>0$ such that $\left|\Phi^{+}(x, t)\right| \leq t^{\theta}$ for all $x \in M$ and all $t \in \mathbb{R}$ satisfying $|t|<t_{0}$;
3. If $\Pi\left(x, t_{1}, t_{2}\right)$ is an admissible rectangle, then $\Phi^{+}\left(x, t_{1}\right)=\Phi^{+}\left(h_{t_{2}}^{-} x, t_{1}\right)$.

For example, if a cocycle $\Phi_{1}^{+}$is defined by $\Phi_{1}^{+}(x, t)=t$, then clearly $\Phi_{1}^{+} \in \mathcal{Y}^{+}$.
In the same way define the space of $\mathcal{Y}^{-}$of Hölder cocyles $\Phi^{-}(x, t)$ over the horizontal flow $h_{t}^{-}$which are invariant under vertical holonomy, and set $\Phi_{1}^{-}(x, t)=t$.

Given $\Phi^{+} \in \mathcal{Y}^{+}, \Phi^{-} \in \mathcal{Y}^{-}$, a finitely additive measure $\Phi^{+} \times \Phi^{-}$on the semi-ring $\overline{\mathfrak{C}}$ of admissible rectangles is introduced by the formula

$$
\begin{equation*}
\Phi^{+} \times \Phi^{-}\left(\Pi\left(x, t_{1}, t_{2}\right)\right)=\Phi^{+}\left(x, t_{1}\right) \cdot \Phi^{-}\left(x, t_{2}\right) . \tag{2}
\end{equation*}
$$

In particular, for $\Phi^{-} \in \mathcal{Y}^{-}$, set $m_{\Phi^{-}}=\Phi_{1}^{+} \times \Phi^{-}$:

$$
\begin{equation*}
m_{\Phi^{-}}\left(\Pi\left(x, t_{1}, t_{2}\right)\right)=t_{1} \Phi^{-}\left(x, t_{2}\right) \tag{3}
\end{equation*}
$$

For any $\Phi^{-} \in \mathcal{Y}^{-}$the measure $m_{\Phi^{-}}$satisfies $\left(h_{t}^{+}\right)_{*} m_{\Phi^{-}}=m_{\Phi^{-}}$and is an invariant distribution in the sense of G. Forni [5], 6]. For instance, $m_{\Phi_{1}^{-}}=\mathfrak{m}$.

A $\mathbb{C}$-linear pairing between $\mathcal{Y}^{+}$and $\mathcal{Y}^{-}$is given, for $\Phi^{+} \in \mathcal{Y}^{+}, \Phi^{-} \in \mathcal{Y}^{-}$, by the formula

$$
\begin{equation*}
<\Phi^{+}, \Phi^{-}>=\Phi^{+} \times \Phi^{-}(M) \tag{4}
\end{equation*}
$$

The space of Lipschitz functions is not invariant under $h_{t}^{+}$, and a larger function space $\operatorname{Lip}_{w}^{+}(M, \omega)$ of weakly Lipschitz functions is introduced as follows. A bounded measurable function $f$ belongs to $\operatorname{Lip}_{w}^{+}(M, \omega)$ if there exists a constant $C$, depending only on $f$, such that for any admissible rectangle $\Pi\left(x, t_{1}, t_{2}\right)$ we have

$$
\begin{equation*}
\mid \int_{0}^{t_{1}} f\left(h_{t}^{+} x\right) d t-\int_{0}^{t_{1}} f\left(h_{t}^{+}\left(h_{t_{2}}^{-} x\right) d t \mid \leq C\right. \tag{5}
\end{equation*}
$$

Let $C_{f}$ be the infimum of all $C$ satisfying (5). We norm $\operatorname{Lip}_{w}^{+}(X)$ by setting

$$
\|f\|_{L i p_{w}^{+}}=\sup _{X} f+C_{f} .
$$

By definition, the space $\operatorname{Lip}_{w}^{+}(M, \omega)$ contains all Lipschitz functions on $M$ and is invariant under $h_{t}^{+}$. We denote by $L i p_{w, 0}^{+}(M, \omega)$ the subspace of $\operatorname{Lip}_{w}^{+}(M, \omega)$ of functions whose integral with respect to $\mathfrak{m}$ is 0 .

### 1.2 Flows along the stable foliation of a pseudo-Anosov diffeomorphism.

Assume that $\theta_{1}>0$ and a diffeomorphism $g: M \rightarrow M$ are such that

$$
\begin{equation*}
g^{*}(\Re(\omega))=\exp \left(\theta_{1}\right) \Re(\omega) ; g^{*}(\Im(\omega))=\exp \left(-\theta_{1}\right) \Im(\omega) \tag{6}
\end{equation*}
$$

The diffeomorphism $g$ induces a linear automorphism $g^{*}$ of the cohomology space $H^{1}(M, \mathbb{C})$. Denote by $E^{+}$the expanding subspace of $g^{*}$ (in other words, $E^{+}$is the subspace spanned by vectors corresponding to Jordan cells of $g^{*}$ with eigenvalues exceeding 1 in absolute value). The action of $g$ on $\mathcal{Y}^{+}$is given by $g^{*} \Phi^{+}(x, t)=\Phi^{+}\left(g x, \exp \left(\theta_{1}\right) t\right)$.

Proposition 1 There exists a $g^{*}$-equivariant isomorphism between $E^{+}$and $\mathcal{Y}^{+}$.
Theorem 1 There exists a continuous mapping $\Xi^{+}: \operatorname{Lip}_{w}^{+}(M, \omega) \rightarrow \mathcal{Y}^{+}$such that for any $f \in \operatorname{Lip}_{w}^{+}(M, \omega)$, any $x \in X$ and any $T>0$ we have

$$
\left|\int_{0}^{T} f \circ h_{t}^{+}(x) d t-\Xi^{+}(f)(x, T)\right|<C_{\varepsilon}\|f\|_{L i p_{w}^{+}}(1+\log (1+T))^{2 \rho+1}
$$

The mapping $\Xi^{+}$satisfies $\Xi^{+}\left(f \circ h_{t}^{+}\right)=\Xi^{+}(f)$ and $\Xi^{+}(f \circ g)=g^{*} \Xi^{+}(f)$.

The mapping $\Xi^{+}$is constructed as follows. By Proposition 1 applied to the flow $h_{t}^{-}$, there exists a $g$-equivariant isomorphism between $\mathcal{Y}^{-}$and the contracting space for the action of $g^{*}$ on $H^{1}(M, \mathbb{C}$ ) (in other words, the subspace spanned by vectors corresponding to Jordan cells with eigenvalues strictly less than 1 in absolute value).
Proposition 2 The pairing $<,>$ given by (4) is nondegenerate and $g^{*}$-invariant.
Remark. Under the identification of $\mathcal{Y}^{+}$and $\mathcal{Y}^{-}$with respective subspaces of $H^{1}(M, \mathbb{C})$, the pairing $<,>$ is taken to the cup-product on $H^{1}(M, \mathbb{C})$ (see Proposition 4.19 in Veech [14]).

If $f \in \operatorname{Lip}_{w}^{+}(M, \omega)$, then $f$ is Riemann-integrable with respect to $m_{\Phi^{-}}$for any $\Phi^{-} \in \mathcal{Y}^{-}$(see (30) for a precise definition of the integral). Assign to $f$ a cocycle $\Phi_{f}^{+}$in such a way that for all $\Phi^{-} \in \mathcal{Y}^{-}$we have

$$
\begin{equation*}
<\Phi_{f}^{+}, \Phi^{-}>=\int_{M} f d m_{\Phi^{-}} \tag{7}
\end{equation*}
$$

By definition, $\Phi_{f \circ h_{t}^{+}}^{+}=\Phi_{f}^{+}$. The mapping $\Xi^{+}$of Theorem 1 is given by the formula

$$
\begin{equation*}
\Xi^{+}(f)=\Phi_{f}^{+} \tag{8}
\end{equation*}
$$

The first eigenvalue for the action of $g^{*}$ on $E^{+}$is $\exp \left(\theta_{1}\right)$ and is always simple. If its second eigenvalue has the form $\exp \left(\theta_{2}\right)$, where $\theta_{2}>0$, and is simple as well, then the following limit theorem holds for $h_{t}^{+}$.

Given a bounded measurable function $f: X \rightarrow \mathbb{R}$ and $x \in X$, introduce a continuous function $\mathfrak{S}_{n}[f, x]$ on the unit interval by the formula

$$
\begin{equation*}
\mathfrak{S}_{n}[f, x](\tau)=\int_{0}^{\tau \exp \left(n \theta_{1}\right)} f \circ h_{t}^{+}(x) d t \tag{9}
\end{equation*}
$$

The functions $\mathfrak{S}_{n}[f, x]$ are $C[0,1]$-valued random variables on the probability space $(M, \mathfrak{m})$.

Theorem 2 If $\left.g^{*}\right|_{E^{+}}$has a simple, real second eigenvalue $\exp \left(\theta_{2}\right), \theta_{2}>0$, then there exists a continuous functional $\alpha: \operatorname{Lip}_{w}^{+}(M, \omega) \rightarrow \mathbb{R}$ and a compactly supported non-degenerate measure $\eta$ on $C[0,1]$ such that for any $f \in \operatorname{Lip}_{w, 0}^{+}(M, \omega)$ satisfying $\alpha(f) \neq 0$ the sequence of random variables

$$
\frac{\mathfrak{S}_{n}[f, x]}{\alpha(f) \exp \left(n \theta_{2}\right)}
$$

converges in distribution to $\eta$ as $n \rightarrow \infty$.
The functional $\alpha$ is constructed explicitly as follows. Under the assumptions of the theorem the action of $g^{*}$ on $E^{-}$has a simple eigenvalue $\exp \left(-\theta_{2}\right)$; let $v(2)$ be the eigenvector with eigenvalue $\exp \left(-\theta_{2}\right)$, let $\Phi_{2}^{-} \in \mathcal{Y}^{-}$correspond to $v(2)$ by Proposition 1 and $m_{\Phi_{2}^{-}}$be given by (3); then

$$
\alpha(f)=\int f d m_{\Phi_{2}^{-}}
$$

### 1.3 Generic translation flows.

Let $\rho \geq 2$ and let $\kappa=\left(\kappa_{1}, \ldots, \kappa_{\sigma}\right)$ be a nonnegative integer vector such that $\kappa_{1}+\cdots+\kappa_{\sigma}=2 \rho-2$. Denote by $\mathcal{M}_{\kappa}$ the moduli space of Riemann surfaces of genus $\rho$ endowed with a holomorphic differential of area 1 with singularities of orders $k_{1}, \ldots, k_{\sigma}$ (the stratum in the moduli space of holomorphic differentials), and let $\mathcal{H}$ be a connected component of $\mathcal{M}_{\kappa}$. Denote by $g_{t}$ the Teichmüller flow on $\mathcal{H}$ (see [6], [8]), and let $\mathbb{A}(t, X)$ be the Kontsevich-Zorich cocycle over $g_{t}$ [8].

Let $\mathbb{P}$ be a $g_{t}$-invariant ergodic probability measure on $\mathcal{H}$. For $X \in \mathcal{H}$, $X=(M, \omega)$, let $\mathcal{Y}_{X}^{+}, \mathcal{Y}_{X}^{-}$be the corresponding spaces of Hölder cocycles. Denote by $E_{X}^{+}$the space spanned by the positive Lyapunov exponents of the KontsevichZorich cocycle.

Proposition 3 For $\mathbb{P}$-almost all $X \in \mathcal{H}$, we have $\operatorname{dim} \mathcal{Y}_{X}^{+}=\operatorname{dim} \mathcal{Y}_{X}^{-}=\operatorname{dim} E_{X}^{+}$, and the pairing $<,>$ between $\mathcal{Y}_{X}^{+}$and $\mathcal{Y}_{X}^{-}$is non-degenerate.

Remark. In particular, if $\mathbb{P}$ is the Masur-Veech "smooth" measure 10, 12, then $\operatorname{dim} \mathcal{Y}_{X}^{+}=\operatorname{dim} \mathcal{Y}_{X}^{-}=\rho$.

Assign to $f \in \operatorname{Lip}_{w}^{+}(M, \omega)$ a cocycle $\Phi_{f}^{+}$by (7).
Theorem 3 For any $\varepsilon>0$ there exists a constant $C_{\varepsilon}$ depending only on $\mathbb{P}$ such that for $\mathbb{P}$-almost every $X \in \mathcal{H}$, any $f \in \operatorname{Lip}_{w}^{+}(X)$, any $x \in X$ and any $T>0$ we have

$$
\left|\int_{0}^{T} f \circ h_{t}^{+}(x) d t-\Phi_{f}^{+}(x, T)\right|<C_{\varepsilon}\|f\|_{L i p_{w}^{+}}\left(1+T^{\varepsilon}\right) .
$$

If both the first and the second Lyapunov exponent of the measure $\mathbb{P}$ are positive and simple (as, by the Avila-Viana Theorem [2], is the case with the Masur-Veech "smooth" measure on $\mathcal{H}$ ), then the following limit theorem holds.

As before, consider a $C[0,1]$-valued random variable $\mathfrak{S}_{t}[f, x]$ on $(M, \mathfrak{m})$ defined by the formula

$$
\mathfrak{S}_{s}[f, x](\tau)=\int_{0}^{\tau \exp (s)} f \circ h_{t}^{+}(x) d t
$$

Let $\|v\|$ be the Hodge norm in $H^{1}(M, \mathbb{R})$. Let $\theta_{2}>0$ be the second Lyapunov exponent of the Kontsevich-Zorich cocycle and let $v_{2}(X)$ be a Lyapunov vector corresponding to $\theta_{2}$ (by our assumption, such a vector is unique up to scalar multiplication). Introduce a real-valued multiplicative cocycle $H_{2}(t, X)$ over $g_{t}$ by the formula

$$
\begin{equation*}
H_{2}(t, X)=\frac{\left\|A(t, X) v_{2}(X)\right\|}{\left\|v_{2}(X)\right\|} . \tag{10}
\end{equation*}
$$

Theorem 4 Assume that both the first and the second Lyapunov exponent of the Kontsevich-Zorich cocycle with respect to the measure $\mathbb{P}$ are positive and simple. Then for $\mathbb{P}$-almost any $X^{\prime} \in \mathcal{H}$ there exists a non-degenerate compactly supported measure $\eta_{X^{\prime}}$ on $C[0,1]$ and, for $\mathbb{P}$-almost all $X, X^{\prime} \in \mathcal{H}$, there exists a
sequence of moments $s_{n}=s_{n}\left(X, X^{\prime}\right)$ such that the following holds. For $\mathbb{P}$-almost every $X \in \mathcal{H}$ there exists a continuous functional

$$
\mathfrak{a}^{(X)}: \operatorname{Lip}_{w}^{+}(X) \rightarrow \mathbb{R}
$$

such that for $\mathbb{P}$-almost every $X^{\prime}$ and for any real-valued $f \in L i p_{w, 0}^{+}(X)$ satisfying $\mathfrak{a}^{(X)}(f) \neq 0$, the sequence of $C[0,1]$-valued random variables

$$
\frac{\mathfrak{S}_{s_{n}}[f, x](\tau)}{\left(\mathfrak{a}^{(X)}(f)\right) H_{2}\left(s_{n}, X\right)}
$$

converges in distribution to $\eta_{X^{\prime}}$ as $n \rightarrow \infty$.
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## 2 Asymptotic foliations of a Markov compactum.

### 2.1 Definitions and notation.

Let $m \in \mathbb{N}$ and let $\Gamma$ be an oriented graph with $m$ vertices $\{1, \ldots, m\}$ and possibly multiple edges. We assume that that for each vertex there is an edge starting from it and an edge ending in it.

Let $\mathcal{E}(\Gamma)$ be the set of edges of $\Gamma$. For $e \in \mathcal{E}(\Gamma)$ we denote by $I(e)$ its initial vertex and by $F(e)$ its terminal vertex. Let $Q$ be the incidence matrix of $\Gamma$ defined by the formula

$$
Q_{i j}=\#\{e \in \mathcal{E}(\Gamma): I(e)=i, F(e)=j\}
$$

By assumption, all entries of the matrix $Q$ are positive. A finite word $e_{1} \ldots e_{k}$, $e_{i} \in \mathcal{E}(\Gamma)$, will be called admissible if $F\left(e_{i+1}\right)=I\left(e_{i}\right), i=1, \ldots, k$.

To the graph $\Gamma$ we assign a Markov compactum $X_{\Gamma}$, the space of bi-infinite paths along the edges:

$$
X_{\Gamma}=\left\{x=\ldots x_{-n} \ldots x_{0} \ldots x_{n} \ldots, x_{n} \in \mathcal{E}(\Gamma), F\left(x_{n+1}\right)=I\left(x_{n}\right)\right\}
$$

Remark. As $\Gamma$ will be fixed throughout this section, we shall often omit the subscript $\Gamma$ from notation and only insert it when the dependence on $\Gamma$ is underlined.

Cylinders in $X_{\Gamma}$ are subsets of the form $\left\{x: x_{n+1}=e_{1}, \ldots, x_{n+k}=e_{k}\right\}$, where $n \in \mathbb{Z}, k \in \mathbb{N}$ and $e_{1} \ldots e_{k}$ is an admissible word. The family of all cylinders forms a semi-ring which we denote by $\mathfrak{C}$.

For $x \in X, n \in \mathbb{Z}$, introduce the sets

$$
\begin{aligned}
\gamma_{n}^{+}(x)=\left\{x^{\prime} \in X_{\Gamma}: x_{t}^{\prime}\right. & \left.=x_{t}, t \geq n\right\} ; \gamma_{n}^{-}(x)
\end{aligned}=\left\{x^{\prime} \in X_{\Gamma}: x_{t}^{\prime}=x_{t}, t \leq n\right\} ; 又 \bigcup_{n \in \mathbb{Z}} \gamma_{n}^{+}(x) ; \gamma_{\infty}^{-}(x)=\bigcup_{n \in \mathbb{Z}} \gamma_{n}^{-}(x) .
$$

The sets $\gamma_{\infty}^{+}(x)$ are leaves of the asymptotic foliation $\mathcal{F}^{+}$on the space $X_{\Gamma}$; the sets $\gamma_{\infty}^{+}(x)$ are leaves of the asymptotic foliation $\mathcal{F}^{-}$on $X_{\Gamma}$.

For $n \in \mathbb{Z}$ let $\mathfrak{C}_{n}^{+}$be the collection of all subsets of $X_{\Gamma}$ of the form $\gamma_{n}^{+}(x)$, $n \in \mathbb{Z}, x \in X$; similarly, $\mathfrak{C}_{n}^{-}$is the collection of all subsets of the form $\gamma_{n}^{-}(x)$. Set

$$
\begin{equation*}
\mathfrak{C}^{+}=\bigcup_{n \in \mathbb{Z}} \mathfrak{C}_{n}^{+} ; \mathfrak{C}^{-}=\bigcup_{n \in \mathbb{Z}} \mathfrak{C}_{n}^{-} \tag{11}
\end{equation*}
$$

The collection $\mathfrak{C}_{n}^{+}$is a semi-ring for any $n \in \mathbb{Z}$. Since every element of $\mathfrak{C}_{n}^{+}$ is a disjoint union of elements of $\mathfrak{C}_{n+1}^{+}$, the collection $\mathfrak{C}^{+}$is a semi-ring as well. The same statements hold for $\mathfrak{C}_{n}^{-}$and $\mathfrak{C}^{-}$.

Let $\exp \left(\theta_{1}\right)$ be the spectral radius of the matrix $Q$, and let $h=\left(h_{1}, \ldots, h_{m}\right)$ be the unique positive eigenvector of $Q$ : we thus have $Q h=\exp \left(\theta_{1}\right) h$. Let $\lambda=\left(\lambda_{1}, \ldots, \lambda_{m}\right)$ be the positive eigenvector of the transpose matrix $Q^{t}$ : we have $Q^{t} \lambda=\exp \left(\theta_{1}\right) \lambda$. The vectors $\lambda, h$ are normalized as follows:

$$
\begin{equation*}
\sum_{i=1}^{m} \lambda_{i}=1 ; \sum_{i=1}^{m} \lambda_{i} h_{i}=1 \tag{12}
\end{equation*}
$$

Introduce a sigma-additive positive measure $\Phi_{1}^{+}$on the semi-ring $\mathfrak{C}^{+}$by the formula

$$
\begin{equation*}
\Phi_{1}^{+}\left(\gamma_{n}^{+}(x)\right)=h_{F\left(x_{n}\right)} \exp \left((n-1) \theta_{1}\right) \tag{13}
\end{equation*}
$$

and a sigma-additive positive measure $\Phi_{1}^{-}$on the semi-ring $\mathfrak{C}^{-}$by the formula

$$
\begin{equation*}
\Phi_{1}^{-}\left(\gamma_{n}^{-}(x)\right)=\lambda_{I\left(x_{n}\right)} \exp \left(-n \theta_{1}\right) \tag{14}
\end{equation*}
$$

Let $n \in \mathbb{Z}, k \in \mathbb{N}$, and let $e_{1} \ldots e_{k}$ be an admissible word. The Parry measure $\nu$ on $X_{\Gamma}$ is defined by the formula

$$
\begin{equation*}
\nu\left(\left\{x: x_{n+1}=e_{1}, \ldots, x_{n+k}=e_{k}\right\}\right)=\lambda_{I\left(e_{k}\right)} h_{F\left(e_{1}\right)} \exp \left(-k \theta_{1}\right) \tag{15}
\end{equation*}
$$

The measures $\Phi_{1}^{+}, \Phi_{1}^{-}$are conditional measures of the Parry measure $\nu$ in the following sense. If $C \in \mathfrak{C}$, then $\gamma_{\infty}^{+}(x) \cap C \in \mathfrak{C}^{+}, \gamma_{\infty}^{-}(x) \cap C \in \mathfrak{C}^{-}$for any $x \in C$, and we have

$$
\begin{equation*}
\nu(C)=\Phi_{1}^{+}\left(\gamma_{\infty}^{+}(x) \cap C\right) \cdot \Phi_{1}^{-}\left(\gamma_{\infty}^{-}(x) \cap C\right) \tag{16}
\end{equation*}
$$

### 2.2 Finitely-additive measures on leaves of asymptotic foliations.

Given $v \in \mathbb{C}^{m}$, write

$$
\begin{equation*}
|v|=\sum_{i=1}^{m}\left|v_{i}\right| . \tag{17}
\end{equation*}
$$

The norms of all matrices in this paper are understood with respect to this norm. Consider the direct-sum decomposition

$$
\mathbb{C}^{m}=E^{+} \oplus E^{-},
$$

where $E^{+}$is spanned by Jordan cells of eigenvalues of $Q$ with absolute value exceeding 1 , and $E^{-}$is spanned by Jordan cells corresponding to eigenvalues of $Q$ with absolute value at most 1 . Let $v \in E^{+}$and for all $n \in \mathbb{Z}$ set $v^{(n)}=$ $Q^{n} v$ (note that $\left.Q\right|_{E^{+}}$is by definition invertible). Introduce a finitely-additive complex-valued measure $\Phi_{v}^{+}$on the semi-ring $\mathfrak{C}^{+}$(defined in (11)) by the formula

$$
\begin{equation*}
\Phi_{v}^{+}\left(\gamma_{n+1}^{+}(x)\right)=\left(v^{(n)}\right)_{F\left(x_{n+1}\right)} . \tag{18}
\end{equation*}
$$

The measure $\Phi_{v}^{+}$is invariant under holonomy along $\mathcal{F}^{-}$: by definition, we have the following

Proposition 4 If $F\left(x_{n}\right)=F\left(x_{n}^{\prime}\right)$, then $\Phi_{v}^{+}\left(\gamma_{n}^{+}(x)\right)=\Phi_{v}^{+}\left(\gamma_{n}^{+}\left(x^{\prime}\right)\right)$.
The measures $\Phi_{v}^{+}$span a complex linear space, which we denote $\mathcal{Y}^{+}$(or, sometimes, $\mathcal{Y}_{\Gamma}^{+}$, when dependence on $\Gamma$ is stressed.) The map

$$
\begin{equation*}
\mathcal{I}: v \rightarrow \Phi_{v}^{+} \tag{19}
\end{equation*}
$$

is an isomorphism between $E^{+}$and $\mathcal{Y}_{\Gamma}^{+}$.
For $Q^{t}$, we have the direct-sum decomposition

$$
\mathbb{C}^{m}=\tilde{E}^{+} \oplus \tilde{E}^{-},
$$

where $\tilde{E}^{+}$is spanned by Jordan cells of eigenvalues of $Q^{t}$ with absolute value exceeding 1 , and $\tilde{E}^{-}$is spanned by Jordan cells corresponding to eigenvalues of $Q^{t}$ with absolute value at most 1. As before, for $\tilde{v} \in \tilde{E}^{+}$set $\tilde{v}^{(n)}=\left(Q^{t}\right)^{n} \tilde{v}$ for all $n \in \mathbb{Z}$, and introduce a finitely-additive complex-valued measure $\Phi_{\tilde{v}}^{\bar{v}}$ on the semi-ring $\mathfrak{C}^{-}$(defined in (11)) by the formula

$$
\begin{equation*}
\Phi_{\tilde{v}}^{-}\left(\gamma_{n}^{-}(x)\right)=\left(\tilde{v}^{(-n)}\right)_{I\left(x_{n}\right)} . \tag{20}
\end{equation*}
$$

By definition, the measure $\Phi_{\overline{\tilde{v}}}^{-}$is invariant under holonomy along $\mathcal{F}^{+}$: more precisely, we have the following

Proposition 5 If $I\left(x_{n}\right)=I\left(x_{n}^{\prime}\right)$, then $\Phi_{\tilde{v}}^{-}\left(\gamma_{n}^{-}(x)\right)=\Phi_{\tilde{\tilde{v}}}^{-}\left(\gamma_{n}^{-}\left(x^{\prime}\right)\right)$.

Let $\mathcal{Y}_{\Gamma}^{-}$be the space spanned by the measures $\Phi_{v}^{-}, v \in \tilde{E}^{+}$. The map

$$
\begin{equation*}
\tilde{\mathcal{I}}: v \rightarrow \Phi_{v}^{-} \tag{21}
\end{equation*}
$$

is an isomorphism between $\tilde{E}^{+}$and $\mathcal{Y}_{\Gamma}^{-}$.
Let $\sigma: X_{\Gamma} \rightarrow X_{\Gamma}$ be the shift defined by $(\sigma x)_{i}=x_{i+1}$. The shift $\sigma$ naturally acts on the spaces $\mathcal{Y}_{\Gamma}^{+}, \mathcal{Y}_{\Gamma}^{-}$: given $\Phi \in \mathcal{Y}_{\Gamma}^{+}\left(\right.$or $\left.\mathcal{Y}_{\Gamma}^{-}\right)$, the measure $\sigma_{*} \Phi$ is defined, for $\gamma \in \mathfrak{C}^{+}$, by the formula

$$
\sigma_{*} \Phi(\gamma)=\Phi(\sigma \gamma)
$$

From the definitions we obtain
Proposition 6 The following diagrams are commutative:


### 2.3 Pairings.

Given $\Phi^{+} \in \mathcal{Y}^{+}, \Phi^{-} \in \mathcal{Y}^{-}$, introduce, in analogy with (16), a finitely additive measure $\Phi^{+} \times \Phi^{-}$on the semi-ring $\mathfrak{C}$ of cylinders in $X_{\Gamma}$ : for any $C \in \mathfrak{C}$ and $x \in C$, set

$$
\begin{equation*}
\Phi^{+} \times \Phi^{-}(C)=\Phi^{+}\left(\gamma_{\infty}^{+}(x) \cap C\right) \cdot \Phi^{-}\left(\gamma_{\infty}^{-}(x) \cap C\right) \tag{22}
\end{equation*}
$$

Note that by Propositions 4. 5, the right-hand side in (22) does not depend on $x \in C$.

More explicitly, let $v \in E^{+}, \tilde{v} \in \tilde{E}^{+}, \Phi_{v}^{+}=\mathcal{I}(v), \Phi_{\tilde{v}}^{-}=\tilde{\mathcal{I}}(\tilde{v})$. As above, denote $v^{(n)}=Q^{n} v, \tilde{v}^{(n)}=\left(Q^{t}\right)^{n} v$. Let $n \in \mathbb{Z}, k \in \mathbb{N}$ and let $e_{1} \ldots e_{k}$ be an admissible word. Then

$$
\begin{equation*}
\Phi_{v}^{+} \times \Phi_{\tilde{v}}^{-}\left(\left\{x: x_{n+1}=e_{1}, \ldots, x_{n+k}=e_{k}\right\}\right)=\left(v^{(n)}\right)_{F\left(e_{1}\right)}\left(\tilde{v}^{(-n-k)}\right)_{I\left(e_{n+k}\right)} \tag{23}
\end{equation*}
$$

There is a natural $\mathbb{C}$-linear pairing $<,>$ between the spaces $\mathcal{Y}_{\Gamma}^{+}$and $\mathcal{Y}_{\Gamma}^{-}$: for $\Phi^{+} \in \mathcal{Y}_{\Gamma}^{+}, \Phi^{-} \in \mathcal{Y}_{\Gamma}^{-}$, set

$$
\begin{equation*}
<\Phi^{+}, \Phi^{-}>=\Phi^{+} \times \Phi^{-}\left(X_{\Gamma}\right) \tag{24}
\end{equation*}
$$

From (23) we derive

Proposition 7 Let $v \in E^{+}, \tilde{v} \in \tilde{E}^{+}, \Phi_{v}^{+}=\mathcal{I}_{\Gamma}(v), \Phi_{\tilde{v}}^{-}=\tilde{\mathcal{I}}_{\Gamma}(\tilde{v})$. Then

$$
\begin{equation*}
<\Phi_{v}^{+}, \Phi_{\tilde{v}}^{-}>=\sum_{i=1}^{m} v_{i} \tilde{v}_{i} \tag{25}
\end{equation*}
$$

In particular, the pairing $<,>$ is non-degenerate and $\sigma^{*}$-invariant.
In particular, for $\Phi^{-} \in \mathcal{Y}^{-}$denote

$$
\begin{equation*}
m_{\Phi^{-}}=\Phi_{1}^{+} \times \Phi^{-} \tag{26}
\end{equation*}
$$

### 2.4 Weakly Lipschitz Functions.

Introduce a function space $\operatorname{Lip}_{w}^{+}(X)$ in the following way. A bounded Borelmeasurable function $f: X \rightarrow \mathbb{C}$ belongs to the space $\operatorname{Lip}_{w}^{+}(X)$ if there exists a constant $C>0$ such that for all $n \geq 0$ and any $x, x^{\prime} \in X$ satisfying $F\left(x_{n+1}\right)=$ $F\left(x_{n+1}^{\prime}\right)$, we have

$$
\begin{equation*}
\left|\int_{\gamma_{n}^{+}(x)} f d \Phi_{1}^{+}-\int_{\gamma_{n}^{+}\left(x^{\prime}\right)} f d \Phi_{1}^{+}\right| \leq C \tag{27}
\end{equation*}
$$

If $C_{f}$ be the infimum of all $C$ satisfying (27), then we norm $\operatorname{Lip}_{w}^{+}(X)$ by setting

$$
\|f\|_{L i p_{w}^{+}}=\sup _{X} f+C_{f} .
$$

As before, let $\operatorname{Lip} p_{w, 0}^{+}(X)$ be the subspace of $\operatorname{Lip}_{w}^{+}(X)$ of functions whose integral with respect to $\nu$ is zero.

Take $\Phi^{-} \in \mathcal{Y}^{-}$. Any function $f \in \operatorname{Lip}_{w}^{+}(X)$ is integrable with respect to the measure $m_{\Phi^{-}}$, defined by (26), in the following sense. Let $\tilde{v} \in E^{-}$be the vector corresponding to $\Phi^{-}$by (20) and let $\tilde{v}^{(n)}=\left(Q^{t}\right)^{n} \tilde{v}$. Recall that

$$
\begin{equation*}
\left|\tilde{v}^{(-n)}\right| \rightarrow 0 \text { exponentially fast as } n \rightarrow \infty \tag{28}
\end{equation*}
$$

Take arbitrary points $x_{i}^{(n)} \in X, n \in \mathbb{N}$ satisfying

$$
\begin{equation*}
F\left(\left(x_{i}^{(n)}\right)_{n}\right)=i, i=1, \ldots, m \tag{29}
\end{equation*}
$$

and consider the expression

$$
\begin{equation*}
\sum_{i=1}^{m}\left(\int_{\gamma_{n}^{+}\left(x_{i}^{(n)}\right)} f d \Phi_{1}^{+}\right) \cdot\left(\tilde{v}^{(1-n)}\right)_{i} \tag{30}
\end{equation*}
$$

By (27) and (28), as $n \rightarrow \infty$ the expression (30) tends to a limit which does not depend on the particular choice of $x_{i}^{(n)}$ satisfying (29). This limit is denoted

$$
m_{\Phi^{-}}(f)=\int_{X} f d m_{\Phi^{-}}
$$

Introduce a measure $\Phi_{f}^{+} \in \mathcal{Y}^{+}$by requiring that for any $\Phi^{-} \in \mathcal{Y}^{-}$we have

$$
\begin{equation*}
<\Phi_{f}^{+}, \Phi^{-}>=\int_{X} f d m_{\Phi^{-}} \tag{31}
\end{equation*}
$$

Note that the mapping $\Xi^{+}: \operatorname{Lip}_{w}^{+}(X) \rightarrow \mathcal{Y}^{+}$given by $\Xi^{+}(f)=\Phi_{f}^{+}$is continuous by definition and satisfies

$$
\begin{equation*}
\Xi^{+}(f \circ \sigma)=\sigma^{*} \Xi^{+}(f) \tag{32}
\end{equation*}
$$

From the definitions we also have
Proposition 8 Let $\Phi^{+}(1), \ldots, \Phi^{+}(r)$ be a basis in $\mathcal{Y}^{+}$and let $\Phi^{-}(1), \ldots, \Phi^{-}(r)$ be the dual basis in $\mathcal{Y}^{-}$with respect to the pairing $<,>$. Then for any $f \in$ Lip $_{w}^{+}(X)$ we have

$$
\Phi_{f}^{+}=\sum_{i=1}^{r}\left(m_{\Phi^{-}(i)}(f)\right) \Phi^{+}(i) .
$$

### 2.5 Approximation.

Let $\Theta$ be a finitely-additive complex-valued measure on the semi-ring $\mathfrak{C}_{0}^{+}$. Assume that there exists a constant $\delta(\Theta)$ such that for all $x, x^{\prime} \in X$ and all $n \geq 0$ we have

$$
\begin{equation*}
\left|\Theta\left(\gamma_{n}^{+}(x)\right)-\Theta\left(\gamma_{n}^{+}\left(x^{\prime}\right)\right)\right| \leq \delta(\Theta) \text { if } F\left(x_{n+1}\right)=F\left(x_{n+1}^{\prime}\right) \tag{33}
\end{equation*}
$$

In this case $\Theta$ will be called a weakly Lipschitz measure.
Lemma 1 There exists a constant $C_{\Gamma}$ depending only on $\Gamma$ such that the following is true. Let $\Theta$ be a weakly Lipschitz finitely-additive complex-valued measure on the semi-ring $\mathfrak{C}_{0}^{+}$. Then there exists a unique $\Phi^{+} \in \mathcal{Y}_{\Gamma}^{+}$such that for all $x \in X$ and all $n>0$ we have

$$
\begin{equation*}
\left|\Theta\left(\gamma_{n}^{+}(x)\right)-\Phi^{+}\left(\gamma_{n}^{+}(x)\right)\right| \leq C_{\Gamma} \delta(\Theta) n^{m+1} \tag{34}
\end{equation*}
$$

Assign to the graph $\Gamma$ the Markov compactum $Y_{\Gamma}$ of one-sided infinite sequences of edges:

$$
Y=\left\{y=y_{1} \ldots y_{n} \cdots: y_{n} \in \mathcal{E}(\Gamma), F\left(y_{n+1}\right)=I\left(y_{n}\right)\right\}
$$

and, as before, let $\sigma$ be the shift on $Y_{\Gamma}:(\sigma y)_{i}=y_{i+1}$. For $y, y^{\prime} \in Y_{\Gamma}$, write $y^{\prime} \searrow y$ if $\sigma y^{\prime}=y$.

Lemma 1 will be derived from
Lemma 2 There exists a constant $C_{\Gamma}$ depending only on $\Gamma$ such that the following is true. Let $\varphi_{n}$ be a sequence of measurable complex-valued functions on $Y_{\Gamma}$. Assume that there exists a constant $\delta$ such that for all $y \in Y$ and all $n \geq 0$ we have

$$
\begin{equation*}
\left.\mid \varphi_{n+1}(y)\right)-\sum_{y^{\prime} \searrow y} \varphi_{n}\left(y^{\prime}\right) \mid \leq \delta \tag{35}
\end{equation*}
$$

and for all $n \geq 0$ and all $y, \tilde{y} \in Y_{\Gamma}$ satisfying $F\left(y_{1}\right)=F\left(\tilde{y}_{1}\right)$ we have

$$
\begin{equation*}
\left.\mid \varphi_{n}(y)\right)-\varphi_{n}(\tilde{y}) \mid \leq \delta \tag{36}
\end{equation*}
$$

Then there exists a unique $v \in E^{+}$such that for all $y \in Y$ and all $n>0$ we have

$$
\begin{equation*}
\left.\mid \varphi_{n}(y)\right)-\left(Q^{n} v\right)_{F\left(y_{n+1}\right)} \mid \leq C_{\Gamma} \delta n^{m+1} \tag{37}
\end{equation*}
$$

Proof of Lemma 2. Take arbitrary points $y(i) \in Y_{\Gamma}$ in such a way that

$$
F\left(y(i)_{1}\right)=i
$$

Introduce a sequence of vectors $v(n) \in \mathbb{C}^{m}$ by the formula

$$
v(n)_{i}=\varphi_{n}(y(i))
$$

From (36) for any $y \in Y$ we have

$$
\left|\varphi_{n}(y)-v(n)_{F\left(y_{1}\right)}\right| \leq \delta,
$$

and from (35), (36) we have

$$
|Q v(n)-v(n+1)| \leq \delta \cdot\|Q\|
$$

To prove Lemma 2, it suffices now to establish the following
Proposition 9 Let $V$ be a finite-dimensional complex linear space, let $S: V \rightarrow$ $V$ be a linear operator and let $V^{+} \subset V$ be the subspace spanned by vectors corresponding to Jordan cells of $S$ with eigenvalues exceeding 1 in absolute value. There exists a constant $C>0$ depending only on $S$ such that the following is true. Assume that the vectors $v(n) \in V, n \in \mathbb{N}$, satisfy

$$
|S v(n)-v(n+1)|<\delta
$$

for all $n \in \mathbb{N}$ and some constant $\delta>0$. Then there exists a unique $v \in V^{+}$such that for all $n \in \mathbb{N}$ we have

$$
\begin{equation*}
\left|S^{n} v-v(n)\right| \leq C \cdot \delta \cdot n^{\operatorname{dim} V-\operatorname{dim} V^{+}+1} \tag{38}
\end{equation*}
$$

Proof of Proposition 9 By definition, the subspace $V^{+}$is $S$-invariant and $S$ is invertible on $V^{+}$; we have furthermore that $\left|Q^{-n} v\right| \rightarrow 0$ exponentially fast as $n \rightarrow \infty$. Let $V^{-}$be the subspace spanned by Jordan cells corresponding to eigenvalues of absolute value at most 1 ; for $v \in V^{-}$, we have $\left|Q^{n} v\right|<C n^{\operatorname{dim} V-\operatorname{dim} V^{+}}$ as $n \rightarrow \infty$. We have the decomposition $V=V^{+} \oplus V^{-}$. Let

$$
u(0)=v(0), u(n+1)=v(n+1)-S v(n)
$$

Decompose $u(n)=u^{+}(n)+u^{-}(n)$, where $u^{+}(n) \in V^{+}, u^{-}(n) \in V^{-}$. Denote

$$
v^{+}(n+1)=u^{+}(n+1)+S u^{+}(n)+\cdots+S^{n} u^{+}(1)
$$

$$
\begin{gathered}
v^{-}(n+1)=u^{-}(n+1)+S u^{-}(n)+\cdots+S^{n} u^{-}(1) \\
v=u^{+}(0)+S^{-1} u^{+}(1)+\cdots+S^{-n} u^{+}(n)+\ldots
\end{gathered}
$$

By definition, $\left|v^{-}(n+1)\right|$ is bounded above by $C \delta n^{\operatorname{dim} V-\operatorname{dim} V^{+}+1}$ and there exists $\tilde{C}$ such that $\left|S^{n} v-v^{+}(n)\right|<\tilde{C} \delta$ for all $n \in \mathbb{N}$, whence (38) follows. Uniqueness of $v$ follows from the fact that for any nonzero $v^{\prime} \in V^{+}$the sequence $\left|S^{n} v^{\prime}\right|$ grows exponentially as $n \rightarrow \infty$. Proposition 9 and Lemmas 1, 2 are proved completely.

Let $f \in \operatorname{Lip} p_{w}^{+}(X)$. We then have a measure $\Theta_{f}$ on the semi-ring $\mathfrak{C}_{0}^{+}$given, for $\gamma \in \mathfrak{C}_{0}^{+}$, by the formula

$$
\Theta_{f}(\gamma)=\int_{\gamma} f d \Phi_{1}^{+}
$$

By (27), the measure $\Theta_{f}$ satisfies the assumptions of Lemma 1. Let $\Xi_{f}^{+} \in \mathcal{Y}^{+}$ be the measure assigned to $\Theta_{f}$ by Lemma 1 .

Lemma 3 Let $f \in \operatorname{Lip}_{w}^{+}(X), \Phi^{-} \in \mathcal{Y}_{\Gamma}^{-}$. Then

$$
\begin{equation*}
<\Xi_{f}^{+}, \Phi^{-}>=\int_{X} f d m_{\Phi^{-}} \tag{39}
\end{equation*}
$$

Proof: Choose the points $x_{i}^{(n)} \in X$ satisfying (29). As above, let $\tilde{v} \in E^{-}$be the vector corresponding to $\Phi^{-}$by (20) and let $\tilde{v}^{(n)}=\left(Q^{t}\right)^{n} \tilde{v}, n \in \mathbb{Z}$. For any $\varepsilon>0$ and $n>0$ sufficiently large, by definition, we have

$$
\begin{equation*}
\left|m_{\Phi^{-}}(f)-\sum_{i=1}^{m}\left(\int_{\gamma_{n}^{+}\left(x_{i}^{(n)}\right)} f d \Phi_{1}^{+}\right) \cdot\left(\tilde{v}^{(-n)}\right)_{i}\right|<\varepsilon \tag{40}
\end{equation*}
$$

By definition of $\Xi_{f}^{+}$and Lemma 1 we have

$$
\mid \sum_{i=1}^{m}\left(\int_{\gamma_{n}^{+}\left(x_{i}^{(n)}\right)} f d \Phi_{1}^{+}\right) \cdot\left(\tilde{v}^{(-n)}\right)_{i}-\sum_{i=1}^{m}\left(\Xi_{f}^{+}\left(\gamma_{n}^{+}\left(x_{i}^{(n)}\right)\right) \cdot\left(\tilde{v}^{(-n)}\right)_{i}\left|<C_{\Gamma} \cdot n^{m+1}\right| \tilde{v}_{i}^{(-n)} \mid\right.
$$

and, by (28), the right-hand side tends to 0 exponentially fast as $n \rightarrow \infty$.
It remains to notice that, by definition,

$$
\sum_{i=1}^{m}\left(\Xi_{f}^{+}\left(\gamma_{n}^{+}\left(x_{i}^{(n)}\right)\right) \cdot\left(\tilde{v}^{(-n)}\right)_{i}=<\Xi_{f}^{+}, \Phi^{-}>\right.
$$

and the Lemma is proved completely.
We have thus established that $\Xi_{f}^{+}=\Phi_{f}^{+}$, where $\Phi_{f}^{+}$is given by (31).

### 2.6 Orderings.

Following S. Ito [7, A.M. Vershik [15, 16, assume that a partial order $\mathfrak{o}$ is given on $\mathcal{E}(\Gamma)$ in such a way that edges starting at a given vertex are ordered linearly, while edges starting at different vertices are not comparable. An edge will be called maximal (with respect to $\mathfrak{o}$ ) if there does not exist a greater edge; minimal, if there does not exist a smaller edge; and an edge $e$ will be called the successor of $e^{\prime}$ if $e>e^{\prime}$ but there does not exist $e^{\prime \prime}$ such that $e>e^{\prime \prime}>e^{\prime}$.

The ordering $\mathfrak{o}$ is extended to a partial ordering of $X_{\Gamma}$ : we write $x<x^{\prime}$ if there exists $l \in \mathbb{Z}$ such that $x_{l}<x_{l}^{\prime}$ and $x_{n}=x_{n}^{\prime}$ for all $n>l$. Under this ordering each leaf $\gamma_{\infty}^{+}$of the foliation $\mathcal{F}^{+}$is linearly ordered, while points lying on different leaves are not comparable.

Let $\operatorname{Max}(\mathfrak{o})$ be the set of points $x \in X, x=\left(x_{n}\right)_{n \in \mathbb{Z}}$, such that each $x_{n}$ is a maximal edge. Similarly, $\operatorname{Min}(\mathfrak{o})$ denotes the set of points $x \in X, x=\left(x_{n}\right)_{n \in \mathbb{Z}}$, such that each $x_{n}$ is a minimal edge. Since edges starting at a given vertex are ordered linearly, the cardinalities of $\operatorname{Max}(\mathfrak{o})$ and $\operatorname{Min}(\mathfrak{o})$ do not exceed $m$.

If a leaf $\gamma_{\infty}^{+}$does not intersect $\operatorname{Max}(\mathfrak{o})$, then it does not have a maximal element; similarly, if $\gamma_{\infty}^{+}$does not intersect $\operatorname{Min}(\mathfrak{o})$, then it does not have a minimal element.

For $x(1), x(2) \in \gamma_{\infty}^{+}$, let

$$
[x(1), x(2)]=\left\{x^{\prime} \in \gamma_{\infty}^{+}: x(1) \leq x^{\prime} \leq x(2)\right\} .
$$

The sets $(x(1), x(2)],[x(1), x(2)),(x(1), x(2))$ are defined similarly.
Proposition 10 Let $x \in X$. If $\gamma_{\infty}^{+}(x) \cap \operatorname{Max}(\mathfrak{o})=\emptyset$, then for any $t \geq 0$ there exists a point $x^{\prime} \in \gamma_{\infty}^{+}(x)$ such that

$$
\begin{equation*}
\Phi_{1}^{+}\left(\left[x, x^{\prime}\right]\right)=t . \tag{41}
\end{equation*}
$$

Proof. Let $V(x)=\left\{t: \exists x^{\prime} \geq x: \Phi_{1}^{+}\left(\left[x, x^{\prime}\right]\right)=t\right\}$. Since $\gamma_{\infty}^{+}(x) \cap \operatorname{Max}(\mathfrak{o})=$ $\emptyset$, for any $n$ there exists $x^{\prime \prime} \in \gamma_{\infty}^{+}(x)$ such that all points in $\gamma_{n}^{+}\left(x^{\prime \prime}\right)$ are greater than $x$. Since $\Phi_{1}^{+}\left(\gamma_{n}^{+}\left(x^{\prime \prime}\right)\right)$ grows exponentially, uniformly in $x^{\prime \prime}$, as $n \rightarrow \infty$, the set $V(x)$ is unbounded. Furthermore, since $\Phi_{1}^{+}\left(\gamma_{n}^{+}\left(x^{\prime \prime}\right)\right)$ decays exponentially, uniformly in $x^{\prime \prime}$, as $n \rightarrow-\infty$, the set $V(x)$ is dense in $\mathbb{R}_{+}$. Finally, by compactness of $X$, the set $V(x)$ is closed, which concludes the proof of the Proposition.

A similar proposition, proved in the same way, holds for negative $t$.
Proposition 11 Let $x \in X$. If $\gamma_{\infty}^{+}(x) \cap \operatorname{Min}(\mathfrak{o})=\emptyset$, then for any $t \geq 0$ there exists a point $x^{\prime} \in \gamma_{\infty}^{+}(x)$ such that

$$
\begin{equation*}
\Phi_{1}^{+}\left(\left[x^{\prime}, x\right]\right)=t . \tag{42}
\end{equation*}
$$

Define an equivalence relation $\sim$ on $X$ by writing $x \sim x^{\prime}$ if $x \in \gamma_{\infty}^{+}\left(x^{\prime}\right)$ and $\Phi_{1}^{+}\left(\left[x, x^{\prime}\right]\right)=\Phi_{1}^{+}\left(\left[x^{\prime}, x\right]\right)=0$. The equivalence classes admit the following explicit description, which is clear from the definitions.

Proposition 12 Let $x, x^{\prime} \in X$ be such that $x \in \gamma_{\infty}^{+}\left(x^{\prime}\right), x<x^{\prime}$ and $\Phi_{1}^{+}\left(\left[x, x^{\prime}\right]\right)=$ 0 . Then there exists $n \in \mathbb{Z}$ such that

1. $x_{n}^{\prime}$ is a successor of $x_{n}$;
2. $x$ is the maximal element in $\gamma_{n}(x)$;
3. $x^{\prime}$ is the minimal element in $\gamma_{n}\left(x^{\prime}\right)$.

In other words, $\Phi_{1}^{+}\left(\left[x, x^{\prime}\right]\right)=0$ if and only if $\left(x, x^{\prime}\right)=\emptyset$. In particular, equivalence classes consist at most of two points and, $\nu$-almost surely, of only one point.

Denote $X_{\mathfrak{o}}=X / \sim$, let $\pi_{\mathfrak{o}}: X \rightarrow X_{\mathfrak{o}}$ be the projection map and set $\nu_{\mathfrak{o}}=$ $\left(\pi_{\mathfrak{o}}\right)_{*} \nu$. The probability spaces $\left(X_{\mathfrak{o}}, \nu_{\mathfrak{o}}\right)$ and $(X, \nu)$ are measurably isomorphic; in what follows, we shall often omit the index $\mathfrak{o}$. The foliations $\mathcal{F}^{+}$and $\mathcal{F}^{-}$ descend to the space $X_{\mathfrak{o}}$; we shall denote their images on $X_{\mathfrak{o}}$ by the same letters and, as before, denote by $\gamma_{\infty}^{+}(x), \gamma_{\infty}^{-}(x)$ the leaves containing $x \in X_{\mathfrak{o}}$.

Now let $x \in X_{\mathfrak{o}}$ satisfy $\gamma_{\infty}^{+}(x) \cap \operatorname{Max}(\mathfrak{o})=\emptyset$. By Proposition 10, for any $t \geq 0$ there exists a unique $x^{\prime}$ satisfying (41). Denote $h_{t}^{+}(x)=x^{\prime}$. Similarly, if $x \in X_{\mathfrak{o}}$ satisfy $\gamma_{\infty}^{+}(x) \cap \operatorname{Min}(\mathfrak{o})=\emptyset$. By Proposition 11, for any $t \geq 0$ there exists a unique $x^{\prime}$ satisfying (42). Denote $h_{-t}^{+}(x)=x^{\prime}$.

We thus obtain a flow $h_{t}^{+}$, which is well-defined on the set

$$
X_{\mathfrak{o}} \backslash\left(\bigcup_{x \in \operatorname{Max}(\mathfrak{o}) \cup \operatorname{Min}(\mathfrak{o})} \gamma_{\infty}^{+}(x)\right)
$$

and, in particular, $\nu$-almost surely on $X_{\mathfrak{o}}$. By (16), the flow $h_{t}^{+}$preserves the measure $\nu$.

More generally, it is clear from the definitions that for any $\Phi^{-} \in \mathcal{Y}^{-}$, the measure $m_{\Phi^{-}}$, defined by (26), satisfies

$$
\left(h_{t}^{+}\right)_{*} m_{\Phi^{-}}=m_{\Phi^{-}},
$$

similarly to G. Forni's invariant distributions [5] 6].
Remark. S.Ito in [7] gives a construction of a flow similar to the one above. The flow $h_{t}^{+}$is a continuous-time analogue of a Vershik automorphism [15] (of which a variant also occurs in Ito's work [7]), and, in fact, is a suspension flow over the corresponding Vershik's automorphism, a point of view adopted in [4].

### 2.7 Decomposition of Arcs.

We assume that an ordering $\mathfrak{o}$ is fixed on $\Gamma$. Denote by $\mathfrak{C}(\mathfrak{o})$ the semi-ring of subsets of $X_{\Gamma}$ of the form $\left[x, x^{\prime}\right)$, where $x<x^{\prime}$. Any measure $\Phi^{+} \in \mathcal{Y}^{+}$can be extended to $\mathfrak{C}(\mathfrak{o})$ in the following way.

Let $\mathfrak{R}_{n}^{+}$be the ring generated by the semi-ring $\mathfrak{C}_{n}^{+}$. For $\gamma \in \mathfrak{C}(\mathfrak{o})$, denote by $\gamma(n)$ the smallest (by inclusion) element of the ring $\mathfrak{R}_{-n}^{+}$containing $\gamma$ and
let $\hat{\gamma}(n)$ be the greatest (by inclusion) element of the ring $\mathfrak{R}_{-n}^{+}$contained in $\gamma$ (possibly, $\hat{\gamma}(n)=\emptyset$ ). By definition,

$$
\begin{gather*}
\hat{\gamma}(n) \subset \hat{\gamma}(n+1) \subset \gamma(n+1) \subset \gamma(n) \\
\gamma(n) \backslash \hat{\gamma}(n)=\bigsqcup_{i=1}^{l_{n}} \gamma_{i}^{(n)} \tag{43}
\end{gather*}
$$

where $\gamma_{i}^{(n)} \in \mathfrak{C}_{-n}^{+}, l_{n} \leq\|Q\|$, and

$$
\begin{equation*}
\gamma(n) \backslash \gamma(n+1)=\bigsqcup_{i=1}^{L_{n}} \gamma_{i}^{(n+1)} \tag{44}
\end{equation*}
$$

where $\gamma_{i}^{(n+1)} \in \mathfrak{C}_{-n-1}^{+}, L_{n} \leq 2\|Q\|$.
By definition, if $\Phi^{+} \in \mathcal{Y}^{+}$, then there are only $m$ possible values of $\Phi^{+}(\gamma)$ for $\gamma \in \mathfrak{C}_{-n}^{+}$, and the maximum of these decays exponentially as $n \rightarrow \infty$. We thus have

Proposition 13 There exists positive constants $C_{\Gamma}$, depending only on $\Gamma$, such that the following is true. Let $v_{0}=0, v_{1}, \ldots, v_{l} \in E^{+}, Q v_{i}=\exp (\theta) v_{i}+v_{i-1}$. Assume $v \in \mathbb{C} v_{1} \oplus \cdots \oplus \mathbb{C} v_{l}$ satisfies $|v|=1$ and let $\Phi_{v}^{+}=\mathcal{I}_{\Gamma}(v)$. Then for any $\gamma \in \mathfrak{C}(\mathfrak{o})$ we have

$$
\begin{aligned}
& \left|\Phi_{v}^{+}(\gamma(n))-\Phi_{v}^{+}(\gamma(n+1))\right| \leq C_{\Gamma} n^{l-1} \exp (-(\Re \theta) n) \\
& \left|\Phi_{v}^{+}(\hat{\gamma}(n))-\Phi_{v}^{+}(\hat{\gamma}(n+1))\right| \leq C_{\Gamma} n^{l-1} \exp (-(\Re \theta) n)
\end{aligned}
$$

decay exponentially as $n \rightarrow \infty$. In particular, if $v \in E^{+}, Q v=\exp (\theta) v,|v|=1$, then

$$
\begin{aligned}
&\left|\Phi_{v}^{+}(\gamma(n))-\Phi_{v}^{+}(\gamma(n+1))\right| \leq C_{\Gamma} \exp (-(\Re \theta) n) \\
&\left|\Phi_{v}^{+}(\hat{\gamma}(n))-\Phi_{v}^{+}(\hat{\gamma}(n+1))\right| \leq C_{\Gamma} \exp (-(\Re \theta) n)
\end{aligned}
$$

Consequently, for any $\Phi^{+} \in \mathcal{Y}^{+}, \gamma \in \mathfrak{C}(\mathfrak{o})$, the sequence $\Phi^{+}(\gamma(n))$ converges as $n \rightarrow \infty$, and we set

$$
\Phi^{+}(\gamma)=\lim _{n \rightarrow \infty} \Phi^{+}(\gamma(n))
$$

By (43), we also have

$$
\Phi^{+}(\gamma)=\lim _{n \rightarrow \infty} \Phi^{+}(\hat{\gamma}(n))
$$

Proposition 14 The measure $\Phi^{+}$is finitely-additive on $\mathfrak{C}(\mathfrak{o})$.
Proof: Let $v \in E^{+}$be such that $\Phi^{+}=\Phi_{v}^{+}$and let $\gamma_{0}, \gamma_{1}, \ldots, \gamma_{k} \in \mathfrak{C}(\mathfrak{o})$ satisfy

$$
\gamma_{0}=\bigsqcup_{i=1}^{k} \gamma_{i}
$$

Consider the $\operatorname{arcs} \gamma_{0}(n), \gamma_{1}(n), \ldots, \gamma_{k}(n)$. We have

$$
\begin{equation*}
\gamma_{0}(n) \subset \bigcup_{i=1}^{k} \gamma_{i}(n) \tag{45}
\end{equation*}
$$

and decompose

$$
\gamma_{i}(n)=\bigsqcup \gamma_{i j}(n+1),
$$

where $\gamma_{i j}(n+1) \in \mathfrak{C}_{-n-1}^{+}$.
By (45), each of the arcs $\gamma_{0 j}(n+1)$ is also encountered among the arcs $\gamma_{i j}(n+1)$ (possibly, more than once, but not more than $k$ times). Consider the collection $\gamma_{i j}(n+1)$ and cross out all the arcs $\gamma_{0 j}(n+1)$; by maximality, and since our ordering is linear on each leaf of the foliation $\mathcal{F}^{+}$, there will remain not more than $2 k\|Q\|$ arcs, whence we obtain

$$
\left|\sum_{i=1}^{k} \Phi^{+}\left(\gamma_{i}(n)\right)-\Phi^{+}\left(\gamma_{0}(n)\right)\right| \leq 2 k\|Q\| \cdot\left|Q^{-n-1} v\right|,
$$

and, since the right-hand side decays exponentially as $n \rightarrow \infty$, the Proposition is proved.

Lemma 4 There exists a constant $C_{\Gamma}$ depending only on $\Gamma$ such that the following is true. Let $f \in \operatorname{Lip} p_{w}^{+}\left(X_{\Gamma}\right)$ and let $\Phi_{f}^{+} \in \mathcal{Y}^{+}$be given by (31). For any $\gamma \in \mathfrak{C}(\mathfrak{o})$ we have

$$
\begin{equation*}
\left|\int_{\gamma} f d \Phi_{1}^{+}-\Phi_{f}^{+}(\gamma)\right| \leq C_{\Gamma}| | f \|_{L i p_{w}^{+}}\left(1+\log \left(1+\Phi_{1}^{+}(\gamma)\right)^{m+1} .\right. \tag{46}
\end{equation*}
$$

Indeed, for $\gamma \in \mathfrak{C}^{+}$this follows from Lemma and for all other arcs from Proposition 13

### 2.8 Ergodic averages of the flow $h_{t}^{+}$.

Let $\Phi^{+} \in \mathcal{Y}^{+}$and denote $\Phi^{+}[x, t]=\Phi_{i}^{+}\left(\left[x, h_{t}^{+} x\right]\right)$. The function $\Phi^{+}(x, t)$ is an additive cocycle over the flow $h_{t}^{+}$. Let $f \in \operatorname{Lip} p_{w}^{+}\left(X_{\Gamma}\right)$, and let $\Phi_{f}^{+}$be defined by (31). By definition, $\Phi_{f \circ h_{t}^{+}}=\Phi_{f}^{+}$; recall from (32) that $\Phi_{f \circ \sigma}^{+}=\sigma^{*} \Phi_{f}^{+}$. Lemma 4 implies
Theorem 5 There exists a positive constant $C_{\Gamma}$ depending only on $\Gamma$ such that for any $f \in \operatorname{Lip}_{w}^{+}\left(X_{\Gamma}\right)$, for all $x \in X$ and all $T>0$ we have

$$
\left|\int_{0}^{T} f \circ h_{t}^{+}(x) d t-\Phi_{f}^{+}(x, t)\right| \leq C_{\Gamma}| | f \|_{L i p}(1+\log (1+T))^{m+1} .
$$

Given a bounded measurable function $f: X \rightarrow \mathbb{R}$ and $x \in X$, introduce a continuous function $\mathfrak{S}_{n}[f, x]$ on the unit interval by the formula

$$
\begin{equation*}
\mathfrak{S}_{n}[f, x](\tau)=\int_{0}^{\tau \exp \left(n \theta_{1}\right)} f \circ h_{t}^{+}(x) d t . \tag{47}
\end{equation*}
$$

The functions $\mathfrak{S}_{n}[f, x]$ are $C[0,1]$-valued random variable on the probability space $\left(X_{\Gamma}, \nu_{\Gamma}\right)$.

Theorem 6 If $Q$ has a simple real second eigenvalue $\exp \left(\theta_{2}\right), \theta_{2}>0$, then there exists a continuous functional $\alpha: \operatorname{Lip}_{w}^{+}(X) \rightarrow \mathbb{R}$ and a compactly supported non-degenerate measure $\eta$ on $C[0,1]$ such that for any $f \in L i p_{w, 0}^{+}(X)$ satisfying $\alpha(f) \neq 0$ the sequence of random variables

$$
\frac{\mathfrak{S}_{n}[f, x]}{\alpha(f) \exp \left(n \theta_{2}\right)}
$$

converges in distribution to $\eta$ as $n \rightarrow \infty$.
Remark. Compactness of the support of $\eta$ is understood in the sense of the Tchebycheff topology on $C[0,1]$. Nondegeneracy of the measure $\eta$ means that if $\varphi \in C[0,1]$ is distributed according to $\eta$, then for any $t_{0} \in(0,1]$ the distribution of the real-valued random variable $\varphi\left(t_{0}\right)$ is not concentrated at a single point.

The measure $\eta$ is constructed as follows: let $v_{2}$ be an eigenvector with eigenvalue $\exp \left(\theta_{2}\right)$, set $\Phi_{2}^{+}=\mathcal{I}\left(v_{2}\right)$ (see (19)); then $\eta$ is the distribution of $\Phi_{2}^{+}(x, \tau)$, $0 \leq \tau \leq 1$, considered as a $C[0,1]$-valued random variable on the space $\left.X_{\Gamma}, \nu_{\Gamma}\right)$. The functional $\alpha(f)$ is constructed as follows: under the assumptions of Theorem 6, the matrix $Q^{t}$ also has the simple real second eigenvalue $\exp \left(\theta_{2}\right)$; let $\tilde{v}_{2}$ be the eigenvector with eigenvalue $\exp \left(\theta_{2}\right)$, normalized in such a way that $\sum_{i=1}^{m}\left(v_{2}\right)_{i}\left(\tilde{v}_{2}\right)_{i}=1$; set $\Phi_{2}^{-}=\tilde{\mathcal{I}}\left(\tilde{v}_{2}\right)$ (see (21)), and let $m_{\Phi_{2}^{-}}$be given by (26); then

$$
\alpha(f)=\int f d m_{\Phi_{2}^{-}}
$$

### 2.9 The diagonalizable case.

As an illustration, consider the case when $\left.Q\right|_{E+}$ is diagonalizable with eigenvalues $\exp \left(\theta_{i}\right), i=1, \ldots, r, \Re\left(\theta_{i}\right)>0$. The Perron-Frobenius vector $h$ corresponds to $\exp \left(\theta_{1}\right)$; let $v_{2}, \ldots, v_{r}$ be eigenvectors corresponding to $\exp \left(\theta_{i}\right)$ : thus $Q v_{i}=\exp \left(\theta_{i}\right) v_{i}, i=2, \ldots, r$ and

$$
E^{+}=\mathbb{C} h \oplus \mathbb{C} v_{2} \oplus \cdots \oplus \mathbb{C} v_{r}
$$

We have a similar direct-sum representation for $Q^{t}$ :

$$
\tilde{E}^{+}=\mathbb{C} \lambda \oplus \mathbb{C} \tilde{v}_{2} \oplus \cdots \oplus \mathbb{C} \tilde{v}_{r}
$$

where $Q^{t} \tilde{v}_{i}=\exp \left(\theta_{i}\right) \tilde{v}_{i}, i=2, \ldots, r$. For $i \neq j$ we have

$$
\begin{equation*}
\sum_{l=1}^{m}\left(v_{i}\right)_{l}\left(\tilde{v}_{j}\right)_{l}=0 \tag{48}
\end{equation*}
$$

and, for normalization, let us assume that for all $i=1, \ldots, r$ we have

$$
\begin{equation*}
\sum_{l=1}^{m}\left(v_{i}\right)_{l}\left(\tilde{v}_{i}\right)_{l}=1 \tag{49}
\end{equation*}
$$

Let $\Phi_{i}^{+}=\mathcal{I}\left(v_{i}\right), \Phi_{i}^{-}=\tilde{\mathcal{I}}\left(\tilde{v}_{i}\right), i=2, \ldots, r$. Since $\Phi_{1}^{+}=\mathcal{I}(h)$, the measures $\Phi_{i}^{+}, i=1, \ldots, r$, form a basis in $\mathcal{Y}^{+}$, for which the measures $\Phi_{1}^{-}=\tilde{\mathcal{I}}(\lambda)$, $\Phi_{2}^{-}, \ldots, \Phi_{r}^{-}$form a dual basis in $\mathcal{Y}^{-}$.

For $i=1, \ldots, r$, from (26) we have the measures $m_{\Phi_{i}^{-}}=\Phi_{1}^{+} \times \Phi_{i}^{-}$. For instance, $m_{\Phi_{1}^{-}}=\nu$. Theorem 5 now implies
Corollary 1 For any $f \in \operatorname{Lip}_{w}^{+}\left(X_{\Gamma}\right)$ we have
$\left|\int_{0}^{T} f \circ h_{t}^{+}(x) d t-T \int_{X} f d \nu-\sum_{i=2}^{r} \Phi_{i}^{+}(x, T)\left(m_{\Phi_{i}^{-}}(f)\right)\right| \leq C_{\Gamma}\|f\|_{L i p}(1+\log (1+T))^{m+1}$,
where $C_{\Gamma}$ is a constant depending only on $\Gamma$.
For the action of the shift we have:

$$
\begin{gather*}
(\sigma)_{*} \Phi_{i}^{+}=\exp \left(-\theta_{i}\right) \Phi_{i}^{+}, i=1, \ldots, r  \tag{50}\\
(\sigma)_{*} \Phi_{i}^{-}=\exp \left(\theta_{i}\right) \Phi_{i}^{-}, i=1, \ldots, r \tag{51}
\end{gather*}
$$

Corollary 1 now yields

$$
\begin{equation*}
\int_{0}^{\tau \exp \left(\theta_{1} n\right)} f \circ h_{t}^{+}(x) d t=\sum_{i=1}^{r} \exp \left(n \theta_{i}\right) m_{\Phi_{i}^{-}}(f) \Phi_{i}^{+}\left(\sigma^{n} x, \tau\right)+O\left(n^{m+1}\right) \tag{52}
\end{equation*}
$$

### 2.10 The Hölder property.

As above, we write $\Phi^{+}(x, t)=\Phi^{+}\left(\left[x, h_{t}^{+} x\right]\right)$. Our next aim is to show that $\Phi^{+}(x, t)$ is Hölder in $t$ for any $x \in X_{\mathfrak{v}}$.

Proposition 15 There exist positive constants $C_{\Gamma}$ and $t_{0}$, depending only on $\Gamma$ such that the following is true. Let $v \in E^{+}, Q v=\exp (\theta) v,|v|=1$. Then for all $x \in X$ and positive $t<t_{0}$ we have

$$
\left|\Phi_{v}^{+}(x, t)\right| \leq C_{\Gamma} t^{\Re \theta / \theta_{1}} .
$$

Proposition 16 There exist positive constants $C_{\Gamma}$ and $t_{0}$, depending only on $\Gamma$ such that the following is true. Let $v_{0}=0, v_{1}, \ldots, v_{l} \in E^{+}, Q v_{i}=\exp (\theta) v_{i}+$ $v_{i-1}$. Assume $v \in \mathbb{C} v_{1} \oplus \cdots \oplus \mathbb{C} v_{l}$ satisfies $|v|=1$. Then for all $x \in X$ and positive $t<t_{0}$ we have

$$
\left|\Phi_{v}^{+}(x, t)\right| \leq C_{\Gamma}|\log t|^{l-1} t^{\Re \theta / \theta_{1}}
$$

Proof of Propositions 15, 16. Denote $\gamma=\left[x, h_{t}^{+} x\right]$. If $t$ is small enough, then $\hat{\gamma}(0)=\emptyset$. Let $n_{0}$ be the smallest positive integer such that $\hat{\gamma}\left(n_{0}\right) \neq \emptyset$. There exist positive constants $C_{1}, C_{2}$, depending only on $\Gamma$, such that

$$
C_{1} t \leq \exp \left(-\theta_{1} n_{0}\right) \leq C_{2} t
$$

and Propositions 15, 16 follow now from Proposition 13 ,

Corollary 2 There exist positive constants $\theta>0$ and $t_{0}>0$ depending only on $Q$ such that for all $v \in E^{+},|v|=1$, all $x \in X$ and all positive $t<t_{0}$ we have

$$
\left|\Phi_{v}^{+}(x, t)\right| \leq t^{\theta / \theta_{1}} .
$$

For $v \in E^{+},|v|=1$ denote

$$
\theta_{v}=\lim _{n \rightarrow \infty} \frac{\log \left|Q^{n} v\right|}{n} .
$$

Corollary 3 For any $\varepsilon>0$ there exists a constant $T_{\varepsilon}$ depending only on $\varepsilon$ and $\Gamma$ such that for any $v \in E^{+},|v|=1$, any $x \in X$ and any $T>T_{\varepsilon}$, we have

$$
\left|\Phi_{v}^{+}(x, T)\right| \leq T^{\theta_{v} / \theta_{1}+\varepsilon} .
$$

Proof: Indeed, let $t_{0}$ be the constant given by Proposition [16 Let $n_{0}=n_{0}(T)$ be the smallest such integer that $T=\tau \exp \left(n(T) \theta_{1}\right)$, where $\tau<t_{0}$. Since $\Phi_{v}^{+}(x, T)=\Phi_{Q^{n} v}^{+}\left(\sigma^{n} x, \tau\right)$ for all $n$, it follows from Proposition 16 that

$$
\left|\Phi^{+}(x, T)\right| \leq C_{\Gamma} n_{0}^{m+1} \exp \left(n_{0} \Re\left(\theta_{v}\right)\right) \leq C_{\Gamma} T^{\theta_{v} / \theta_{1}+\varepsilon}
$$

if $T$ is sufficiently large (depending only on $\varepsilon$ ).
Corollary 4 For any $v \in E^{+}$we have

$$
\begin{equation*}
\limsup _{T \rightarrow \infty} \frac{\log \left|\Phi_{v}^{+}(x, T)\right|}{\log T}=\frac{\theta_{v}}{\theta_{1}} . \tag{53}
\end{equation*}
$$

Indeed, the upper bound for the limit superior follows from Corollary 3 and the lower bound is immediate from the relation $\Phi_{v}^{+}\left(\gamma_{n}(x)\right)=\left(Q^{n} v\right)_{F\left(x_{n+1}\right)}$.
Corollary 5 For any $\tau \in \mathbb{R}$ and any $v \in E^{+}$satisfying $v \neq 0, \sum_{i=1}^{m} v_{i} \lambda_{i}=0$, the function $\Phi_{v}^{+}(x, \tau)$ is not a constant in $x$.

Proof: Indeed, assume $\Phi_{v}^{+}(x, \tau)=c$ identically. Then $\Phi^{+}(x, k \tau)=k c$, which contradicts (53): is $c=0$, then the limit superior is 0 ; if $c \neq 0$, then the limit superior is 1 .

### 2.11 Tightness.

In this subsection, we assume that $Q$ has a simple real second eigenvalue $\exp \left(\theta_{2}\right)$, $\theta_{2}>0$. Let $v_{2}$ be the corresponding eigenvector and let $\Phi_{2}^{+}=\mathcal{I}\left(v_{2}\right)$. Take $x \in X$ and consider $\Phi^{+}(x, \tau)$ as a continuous function of $\tau$ on the unit interval. Let $\eta$ be the distribution of $\Phi_{2}^{+}(x, \tau)$ in $C[0,1]$. Note that by Corollary 5 for any $\tau_{0}$ the value of $\Phi_{2}^{+}(x, \tau)$ is not constant on $X$, so the measure $\eta$ is nondegenerate.

Let $\mathfrak{S}_{n}[f, x]$ be defined by the equation (47). Introduce a sequence of measures $\mu_{n}$ on $C[0,1]$ by the formula $\mu_{n}=\mathfrak{S}[n, f]_{*} \nu_{\Gamma}$.

By Theorem 8.1 in Billingsley [3], p.54, to prove Theorem 6 it suffices to establish the following two Lemmas.

Lemma 5 Finite-dimensional distributions of the measures $\mu_{n}$ weakly converge to those of $\eta$.

Lemma 6 The family $\mu_{n}$ is tight in $C[0,1]$.
Proof of Lemma 5 By Theorem 5

$$
\int_{0}^{T} f \circ h_{t}^{+}(x) d t=\Phi_{f}^{+}(x, T)+O\left((\log T)^{m+1}\right) .
$$

Let $v_{2}$ be the eigenvector corresponding to the eigenvalue $\exp \left(\theta_{2}\right),|v|=1$, and let $\Phi_{2}^{+} \in \mathcal{Y}^{+}$be the corresponding measure. We have

$$
E^{+}=\mathbb{C} v_{2} \oplus E_{3}
$$

where $E_{3}$ is spanned by Jordan cells corresponding to eigenvalues with absolute value less than $\exp \left(\theta_{2}\right)$. Let $\zeta$ be a number smaller than $\theta_{2}$ but greater than the spectral radius of $\left.Q\right|_{E_{3}}$. Write

$$
\begin{equation*}
\Phi_{f}^{+}=\alpha(f) \Phi_{2}^{+}+\beta(f) \Phi_{v_{3}}^{+} \tag{54}
\end{equation*}
$$

where $v_{3} \in E^{+},\left|v_{3}\right|=1$, and $\alpha(f), \beta(f)$ are continuous functionals on $\operatorname{Lip}_{w}^{+}(X)$, so, in particular, we have

$$
|\alpha(f)|<C_{01}\|f\|_{L i p_{w}^{+}} ;|\beta(f)|<C_{02}\|f\|_{L i p_{w}^{+}},
$$

where the constants $C_{01}, C_{02}$ only depend on $\Gamma$.
By Corollary 2 there exists $t_{0}$ depending only on $\Gamma$ such that for any positive $t$ such that $t<t_{0}$, any $x \in X$ and any $v \in E^{+}$satisfying $|v|=1$ we have

$$
\begin{equation*}
\left|\Phi_{v}^{+}(x, t)\right| \leq 1 \tag{55}
\end{equation*}
$$

Write $T=t \exp \left(n \theta_{1}\right)$, where $t<t_{0}$. Since $\Phi_{v_{3}}^{+}(x, T)=\Phi_{Q^{n} v_{3}}^{+}\left(\sigma^{n} x, t\right)$, for all sufficiently large $n$, we have $\left|Q^{n} v_{3}\right|<\exp (\zeta n)$ and therefore

$$
\begin{equation*}
\left|\Phi_{v_{3}}^{+}\left(x, \tau \exp \left(n \theta_{1}\right)\right)\right|<\exp (n \zeta) \tag{56}
\end{equation*}
$$

for all $x \in X$. By Theorem 5w have

$$
\begin{equation*}
\left|\int_{0}^{\tau \exp \left(n \theta_{1}\right)} f \circ h_{t}^{+}(x) d t-\Phi_{f}^{+}\left(x, \tau \exp \left(\theta_{1} n\right)\right)\right|=O\left(n^{m+1}\right) \tag{57}
\end{equation*}
$$

Since

$$
\Phi_{f}^{+}\left(x, \tau \exp \left(n \theta_{1}\right)\right)=\alpha(f) \Phi_{2}^{+}\left(\left(x, \tau \exp \left(n \theta_{1}\right)\right)+\beta(f) \Phi_{v_{3}}^{+}\left(x, \tau \exp \left(n \theta_{1}\right)\right)\right.
$$

combining the equality

$$
\Phi_{2}^{+}\left(x, \tau \exp \left(n \theta_{1}\right)\right)=\exp \left(n \theta_{2}\right) \Phi_{2}^{+}\left(\sigma^{n} x, \tau\right)
$$

with the bound (56), we obtain, for all large $n$ and all $x \in X$, uniformly in $\tau \in[0,1]$, the estimate

$$
\left|\mathfrak{S}_{n}[f, x](\tau)-\alpha(f) \Phi_{2}^{+}\left(\sigma^{n} x, \tau\right)\right| \leq C_{\Gamma}\|f\|_{L i p_{w}^{+}} \exp \left(\left(\zeta-\theta_{2}\right) n\right)
$$

Since $\sigma$ preserves the measure $\nu$, it follows that the $k$-dimensional distributions of $\left(\mathfrak{S}_{n}[f, x]\left(\tau_{1}\right), \mathfrak{S}_{n}[f, x]\left(\tau_{2}\right), \ldots, \mathfrak{S}_{n}[f, x]\left(\tau_{k}\right)\right)$ converge to the $k$-dimensional distribution of $\left(\Phi_{2}^{+}\left(x, \tau_{1}\right), \Phi_{2}^{+}\left(x, \tau_{2}\right), \ldots, \Phi_{2}^{+}\left(x, \tau_{k}\right)\right)$, and Lemma 5 is proved.

The argument above yields also
Proposition 17 There exist positive constants $C_{0}=C_{0}(\Gamma)$ and $T_{0}=T_{0}(\Gamma)$ such that for any $x \in X$, any $f \in L i p_{w, 0}^{+}(X)$ and any $T>T_{0}$ we have

$$
\left|\int_{0}^{T} f \circ h_{t}^{+}(x) d t\right| \leq C_{0} \cdot\|f\|_{L i p_{w}^{+}} \cdot T^{\theta_{2} / \theta_{1}}
$$

Indeed, for sufficiently large $T, T=t \exp \left(n \theta_{1}\right)$, where $t<t_{0}$, from (54) we have

$$
\Phi_{f}^{+}(x, T)=\alpha(f) \exp \left(n \theta_{2}\right) \Phi_{2}^{+}\left(\sigma^{n} x, t\right)+O(\exp (n \zeta)) .
$$

Since, by (55), we have $\left|\Phi_{2}^{+}\left(\sigma^{n} x, t\right)\right| \leq 1$, Proposition 17 is established.
We proceed to the proof of Lemma 6 .
Proposition 18 There exists a constant $C_{\Gamma}$ depending only on $\Gamma$ such that for any $f \in \operatorname{Lip} p_{w, 0}^{+}(X)$, any $n>0$, any $x \in X$ and any $\tau_{1}, \tau_{2} \in[0,1]$, we have

$$
\left|\mathfrak{S}_{n}[x, f]\left(\tau_{2}\right)-\mathfrak{S}_{n}[x, f]\left(\tau_{1}\right)\right| \leq C_{\Gamma}| | f \|_{L i p_{w}^{+}}\left|\tau_{2}-\tau_{1}\right|^{\theta_{2} / \theta_{1}}
$$

Lemma 6 follows from Proposition 18 by the Arzelà-Ascoli Theorem.
Proof of Proposition 18. Let $\tau_{1}, \tau_{2} \in[0,1], \tau_{1}<\tau_{2}$. For brevity, write $\mathfrak{S}_{n}=\mathfrak{S}_{n}[f, x]$. We have then

$$
\mathfrak{S}_{n}\left(\tau_{2}\right)-\mathfrak{S}_{n}\left(\tau_{1}\right)=\frac{1}{\exp \left(n \theta_{2}\right)} \int_{\tau_{1} \exp \left(n \theta_{1}\right)}^{\tau_{2} \exp \left(n \theta_{1}\right)} f \circ h_{t}^{+}(x) d t
$$

Let $T_{0}$ be the constant given by Proposition 17 and assume first that

$$
\left(\tau_{2}-\tau_{1}\right) \cdot \exp \left(n \theta_{1}\right) \geq T_{0}
$$

By Proposition 17 we have

$$
\int_{\tau_{1} \exp \left(n \theta_{1}\right)}^{\tau_{2} \exp \left(n \theta_{1}\right)} f \circ h_{t}(x) d t \leq C\|f\|_{L i p_{w}^{+}} \cdot\left(\tau_{2}-\tau_{1}\right)^{\theta_{2} / \theta_{1}} \exp \left(n \theta_{2}\right)
$$

and, consequently,

$$
\left|\mathfrak{S}_{n}\left(\tau_{2}\right)-\mathfrak{S}_{n}\left(\tau_{1}\right)\right| \leq C_{33}\left(\tau_{2}-\tau_{1}\right)^{\theta_{2} / \theta_{1}}
$$

where the constant $C_{33}$ only depends on $\Gamma$.
Now let $\tau_{2}-\tau_{1}=\tau_{0} \exp \left(-n \theta_{1}\right), \tau_{0}<T_{0}$. Since

$$
\exp \left(-n \theta_{2}\right)=\left(\left(\tau_{2}-\tau_{1}\right) / \tau_{0}\right)^{\theta_{2} / \theta_{1}}
$$

using boundedness of $f$, write

$$
\begin{gathered}
\frac{1}{\exp \left(n \theta_{2}\right)} \int_{\tau_{1}}^{\tau_{2} \exp \left(n \theta_{1}\right)} f \circ h_{t}^{+}(x) d t \leq \exp \left(-n \theta_{2}\right) \cdot\|f\|_{\infty} \cdot \tau_{0} \leq \\
\leq \tau_{0}^{1-\theta_{2} / \theta_{1}}\|f\|_{\infty}\left(\tau_{2}-\tau_{1}\right)^{\theta_{2} / \theta_{1}} \leq T_{0}^{1-\theta_{2} / \theta_{1}}\|f\|_{\infty}\left(\tau_{2}-\tau_{1}\right)^{\theta_{2} / \theta_{1}},
\end{gathered}
$$

and the Proposition is proved. Theorem 6 is proved completely.

### 2.12 A symbolic coding for translation flows on surfaces.

To derive Theorems 1 , 2 from Theorems 5, 6, it remains to observe that the vertical flow on the stable foliation of a pseudo-Anosov diffeomorphism is isomorphic to a symbolic flow on the asymptotic foliation of a Markov compactum obtained from the decomposition of the underlying surface into Veech's zippered rectangles, see [4], Sec. 4. The identification of $E^{+}$(and, consequently, of $\mathcal{Y}^{+}$) with the corresponding subspace in cohomology is given by Proposition 4.16 in Veech[14]. The fact that the pairing between cocycles corresponds to the cup-product is immediate from Proposition 4.19 in [14].

## 3 Spaces of Markov Compacta.

Let $\mathfrak{G}$ be the set of all oriented graphs on $m$ vertices such that there is an edge starting at every vertex and an edge ending at every vertex. As before, for a graph $\Gamma \in \mathfrak{G}$, we denote by $\mathcal{E}(\Gamma)$ the set of its edges and by $A(\Gamma)$ its incidence matrix: $A_{i j}(\Gamma)=\#\{e \in \mathcal{E}(\Gamma): I(e)=i, F(e)=j\}$. Denote $\Omega=\mathfrak{G}^{\mathbb{Z}}$ :

$$
\Omega=\left\{\omega=\ldots \omega_{-n} \ldots \omega_{n} \ldots, \omega_{i} \in \mathfrak{G}, i \in \mathbb{Z}\right\}
$$

For $\omega \in \Omega$, denote by $X(\omega)$ the corresponding Markov compactum:

$$
X(\omega)=\left\{x=\ldots x_{-n} \ldots x_{n} \ldots, x_{n} \in \mathcal{E}\left(\omega_{n}\right), F\left(x_{n+1}\right)=I\left(x_{n}\right)\right\}
$$

For $x \in X, n \in \mathbb{Z}$, introduce the sets

$$
\begin{gathered}
\gamma_{n}^{+}(x)=\left\{x^{\prime} \in X(\omega): x_{t}^{\prime}=x_{t}, t \geq n\right\} ; \gamma_{n}^{-}(x)=\left\{x^{\prime} \in X(\omega): x_{t}^{\prime}=x_{t}, t \leq n\right\} \\
\gamma_{\infty}^{+}(x)=\bigcup_{n \in \mathbb{Z}} \gamma_{n}^{+}(x) ; \gamma_{\infty}^{-}(x)=\bigcup_{n \in \mathbb{Z}} \gamma_{n}^{-}(x)
\end{gathered}
$$

The sets $\gamma_{\infty}^{+}(x)$ are leaves of the asymptotic foliation $\mathcal{F}_{\omega}^{+}$on $X(\omega)$; the sets $\gamma_{\infty}^{-}(x)$ are leaves of the asymptotic foliation $\mathcal{F}_{\omega}^{-}$on $X(\omega)$.

For $n \in \mathbb{Z}$ let $\mathfrak{C}_{n, \omega}^{+}$be the collection of all subsets of $X(\omega)$ of the form $\gamma_{n}^{+}(x)$, $n \in \mathbb{Z}, x \in X$; similarly, $\mathfrak{C}_{n, \omega}^{-}$is the collection of all subsets of the form $\gamma_{n}^{-}(x)$. Set

$$
\begin{equation*}
\mathfrak{C}_{\omega}^{+}=\bigcup_{n \in \mathbb{Z}} \mathfrak{C}_{n, \omega}^{+} ; \mathfrak{C}_{\omega}^{-}=\bigcup_{n \in \mathbb{Z}} \mathfrak{C}_{n, \omega}^{-} \tag{58}
\end{equation*}
$$

Just as in the periodic case, the collections $\mathfrak{C}_{n, \omega}^{+}, \mathfrak{C}_{n, \omega}^{-}, \mathfrak{C}_{\omega}^{+}, \mathfrak{C}_{\omega}^{-}$are semi-rings.
Remark. To make notation lighter, we shall often omit the subscript $\omega$ and only include it when dependence on $\omega$ is underlined.

### 3.1 Measures and Cocycles.

Let $\sigma$ be the shift on $\Omega$ given by the formula $(\sigma \omega)_{n}=\omega_{n+1}$. Let $\mathbb{P}$ be an ergodic $\sigma$-invariant probability measure on $\Omega$. We then have a natural cocycle $\mathbb{A}$ on the system $(\Omega, \sigma, \mathbb{P})$ defined, for $n>0$, by the formula

$$
\mathbb{A}(n, \omega)=A\left(\omega_{n}\right) \ldots A\left(\omega_{1}\right)
$$

The cocycle $\mathbb{A}$ will be called the renormalization cocycle.
We need the following assumptions on the measure $\mathbb{P}$ and on the cocyle $\mathbb{A}$.
Assumption 1 The matrices $A\left(\omega_{n}\right)$ are almost surely invertible with respect to $\mathbb{P}$. There exists $\Gamma \in \mathfrak{G}$ such that $\mathbb{P}(\Gamma)>0$.

Assumption 2 The logarithm of the renormalization cocycle (and of its inverse) is integrable.

For $n<0$ set

$$
\mathbb{A}(n, \omega)=A^{-1}\left(\omega_{-n}\right) \ldots A^{-1}\left(\omega_{0}\right)
$$

and set $\mathbb{A}(0, \omega)$ to be the identity matrix.
The transpose cocycle $\mathbb{A}^{t}$ over the dynamical system $\left(\Omega, \sigma^{-1}, \mathbb{P}\right)$ defined, for $n>0$, by the formula

$$
\mathbb{A}^{t}(n, \omega)=A^{t}\left(\omega_{1-n}\right) \ldots A^{t}\left(\omega_{0}\right)
$$

Similarly, for $n<0$ write

$$
\mathbb{A}^{t}(n, \omega)=\left(A^{t}\right)^{-1}\left(\omega_{-n}\right) \ldots\left(A^{t}\right)^{-1}\left(\omega_{1}\right)
$$

and set $\mathbb{A}^{t}(0, \omega)$ to be the identity matrix.
By Assumptions 12 for $\mathbb{P}$-almost any $\omega \in \Omega$ we have the decompositions

$$
\mathbb{R}^{m}=E_{\omega}^{+} \oplus E_{\omega}^{-} ; \mathbb{R}^{m}=\tilde{E}_{\omega}^{+} \oplus \tilde{E}_{\omega}^{-}
$$

where $E^{+}$is the Lyapunov subspace corresponding to positive Lyapunov exponents of $\mathbb{A} ; \tilde{E}^{+}$is the Lyapunov subspace corresponding to positive Lyapunov exponents of $\mathbb{A}^{t} ; E^{-}$is the Lyapunov subspace corresponding to zero and negative Lyapunov exponents of $\mathbb{A} ; E^{-}$is the Lyapunov subspace corresponding to
zero and negative Lyapunov exponents of $\mathbb{A}^{t}$. The standard inner product on $\mathbb{R}^{m}$ yields a nondegenerate pairing between the spaces $E_{\omega}^{+}$and $\tilde{E}_{\omega}^{+}$.

In particular, by Assumption 1, the spaces $E_{\omega}^{+}$and $\tilde{E}_{\omega}^{+}$each contain a unique vector all whose coordinates are positive; we denote these vectors by $h^{(\omega)}$ and $\lambda^{(\omega)}$, respectively, and assume that they are normalized by (12).

Let $v \in E_{\omega}^{+}$and for all $n \in \mathbb{Z}$ set $v^{(n)}=\mathbb{A}(n, \omega) v$. Introduce a finitelyadditive complex-valued measure $\Phi_{v}^{+}$on the semi-ring $\mathfrak{C}_{\omega}^{+}$(defined in (58)) by the formula

$$
\begin{equation*}
\Phi_{v}^{+}\left(\gamma_{n+1}^{+}(x)\right)=\left(v^{(n)}\right)_{F\left(x_{n+1}\right)} \tag{59}
\end{equation*}
$$

As before, the measure $\Phi_{v}^{+}$is invariant under holonomy along $\mathcal{F}^{-}$: by definition, we have the following

Proposition 19 If $F\left(x_{n}\right)=F\left(x_{n}^{\prime}\right)$, then $\Phi_{v}^{+}\left(\gamma_{n}^{+}(x)\right)=\Phi_{v}^{+}\left(\gamma_{n}^{+}\left(x^{\prime}\right)\right)$.
The measures $\Phi_{v}^{+}$span a complex linear space, which is denoted $\mathcal{Y}_{\omega}^{+}$. The $\operatorname{map} \mathcal{I}_{\omega}: v \rightarrow \Phi_{v}^{+}$is an isomorphism between $E_{\omega}^{+}$and $\mathcal{Y}_{\omega}^{+}$. Set $\Phi_{1, \omega}^{+}=\mathcal{I}_{\omega}\left(h^{(\omega)}\right)$.

Now for $\tilde{v} \in \tilde{E}^{+}$and for all $n \in \mathbb{Z}$ set $\tilde{v}^{(n)}=\mathbb{A}^{t}(n, \omega) \tilde{v}$ and introduce a finitely-additive complex-valued measure $\Phi_{\tilde{v}}^{-}$on the semi-ring $\mathfrak{C}_{\omega}^{-}$(defined in (581) ) by the formula

$$
\begin{equation*}
\Phi_{\tilde{v}}^{-}\left(\gamma_{n}^{-}(x)\right)=\left(\tilde{v}^{(-n)}\right)_{I\left(x_{n}\right)} . \tag{60}
\end{equation*}
$$

By definition, the measure $\Phi_{\tilde{v}}^{-}$is invariant under holonomy along $\mathcal{F}^{+}$: more precisely, we have the following

Proposition 20 If $I\left(x_{n}\right)=I\left(x_{n}^{\prime}\right)$, then $\Phi_{\tilde{v}}^{-}\left(\gamma_{n}^{-}(x)\right)=\Phi_{\tilde{v}}^{-}\left(\gamma_{n}^{-}\left(x^{\prime}\right)\right)$.
Let $\mathcal{Y}_{\omega}^{-}$be the space spanned by the measures $\Phi_{\tilde{v}}^{-}, \tilde{v} \in \tilde{E}^{+}$. The map $\tilde{\mathcal{I}}_{\omega}: \tilde{v} \rightarrow \Phi_{\tilde{v}}^{-}$is an isomorphism between $\tilde{E}_{\omega}^{+}$and $\mathcal{Y}_{\omega}^{-}$. Set $\Phi_{1, \omega}^{-}=\tilde{\mathcal{I}}_{\omega}\left(\lambda^{(\omega)}\right)$.

Define a map $t_{\sigma}: X_{\omega} \rightarrow X_{\sigma \omega}$ by $\left(t_{\sigma} x\right)_{i}=x_{i+1}$. The map $t_{\sigma}$ induces a map $t_{\sigma}^{*}: \mathcal{Y}_{\sigma \omega}^{+} \rightarrow \mathcal{Y}_{\omega}^{+}$given, for $\Phi_{\sigma \omega}^{+} \in \mathcal{Y}_{\sigma \omega}^{+}$and $\gamma \in \mathfrak{C}_{\omega}^{+}$, by the formula

$$
t_{\sigma}^{*} \Phi^{+}(\gamma)=\Phi_{\sigma \omega}^{+}\left(t_{\sigma} \gamma\right)
$$

We have the following commutative diagrams:

$$
\begin{gathered}
E_{\omega}^{+} \xrightarrow{\mathcal{I}_{\omega}} \mathcal{Y}_{\omega}^{+} \\
\downarrow^{\mathbb{A}(1, \omega)} \\
E_{\sigma \omega}^{+} \xrightarrow{\mathcal{I}_{\sigma \omega}} \mathcal{Y}_{\sigma \omega}^{+} \\
\tilde{E}_{\sigma}^{*} \\
\tilde{E}_{\omega}^{+} \xrightarrow{\tilde{\mathcal{I}}_{\omega}} \mathcal{Y}_{\omega}^{-} \\
\uparrow_{\mathbb{A}^{t}(1, \sigma \omega)} \\
\tilde{E}_{\sigma \omega}^{+} \xrightarrow{\tilde{\mathcal{I}}_{\sigma \omega}}{ }^{t_{\sigma}^{*}} \\
\mathcal{Y}_{\sigma \omega}^{-}
\end{gathered}
$$

### 3.2 Pairings and weakly Lipschitz functions.

Given $\Phi^{+} \in \mathcal{Y}_{\omega}^{+}, \Phi^{-} \in \mathcal{Y}_{\omega}^{-}$, introduce a finitely additive measure $\Phi^{+} \times \Phi^{-}$on the semi-ring $\mathfrak{C}$ of cylinders in $X(\omega)$ : for any $C \in \mathfrak{C}$ and $x \in C$, set

$$
\begin{equation*}
\Phi^{+} \times \Phi^{-}(C)=\Phi^{+}\left(\gamma_{\infty}^{+}(x) \cap C\right) \cdot \Phi^{-}\left(\gamma_{\infty}^{-}(x) \cap C\right) \tag{61}
\end{equation*}
$$

Note that by Propositions 19, 20, the right-hand side in (61) does not depend on $x \in C$.

As above, for $\Phi^{-} \in \mathcal{Y}_{\omega}^{-}$, denote

$$
\begin{equation*}
m_{\Phi^{-}}=\Phi_{1}^{+} \times \Phi^{-} \tag{62}
\end{equation*}
$$

In particular, we have a positive countably additive measure

$$
\nu_{\omega}=\Phi_{h^{(\omega)}}^{+} \times \Phi_{\lambda(\omega)}^{-}
$$

There is a natural $\mathbb{C}$-linear pairing $<,>$ between the spaces $\mathcal{Y}_{\omega}^{+}$and $\mathcal{Y}_{\omega}^{-}$: for $\Phi^{+} \in \mathcal{Y}_{\omega}^{+}, \Phi^{-} \in \mathcal{Y}_{\omega}^{-}$, set

$$
\begin{equation*}
<\Phi^{+}, \Phi^{-}>=\Phi^{+} \times \Phi^{-}(X(\omega)) \tag{63}
\end{equation*}
$$

As in Sec. 2.3, we have
Proposition 21 Let $v \in E_{\omega}^{+}, \tilde{v} \in \tilde{E}_{\omega}^{+}, \Phi_{v}^{+}=\mathcal{I}_{\omega}(v), \Phi_{\tilde{v}}^{-}=\tilde{\mathcal{I}}_{\omega}(\tilde{v})$. Then

$$
\begin{equation*}
<\Phi_{v}^{+}, \Phi_{\tilde{v}}^{-}>=\sum_{i=1}^{m} v_{i} \tilde{v}_{i} \tag{64}
\end{equation*}
$$

The pairing $<,>$ is non-degenerate and $t_{\sigma}^{*}$-invariant.
The function space $\operatorname{Lip}_{w}^{+}(X(\omega))$ is introduced in the same way as before: a bounded Borel-measurable function $f: X(\omega) \rightarrow \mathbb{C}$ belongs to the space $L i p_{w}^{+}(X)$ if there exists a constant $C>0$ such that for all $n \geq 0$ and any $x, x^{\prime} \in X$ satisfying $F\left(x_{n+1}\right)=F\left(x_{n+1}^{\prime}\right)$, we have

$$
\begin{equation*}
\left|\int_{\gamma_{n}^{+}(x)} f d \Phi_{1}^{+}-\int_{\gamma_{n}^{+}\left(x^{\prime}\right)} f d \Phi_{1}^{+}\right| \leq C \tag{65}
\end{equation*}
$$

and, if $C_{f}$ is the infimum of all $C$ satisfying (65), then we norm $\operatorname{Lip}_{w}^{+}(X)$ by setting

$$
\|f\|_{L i p_{w}^{+}}=\sup _{X} f+C_{f} .
$$

As before, we denote by $\operatorname{Lip}_{w, 0}^{+}(X(\omega))$ the subspace of functions of $\nu_{\omega}$-integral zero.

Take $\Phi^{-} \in \mathcal{Y}^{-}$. Any function $f \in \operatorname{Lip}_{w}^{+}(X)$ is integrable with respect to the measure $m_{\Phi^{-}}$in the same sense as in Sec. [2.4, and a measure $\Phi_{f}^{+} \in \mathcal{Y}^{+}$is defined by the requirement that for any $\Phi^{-} \in \mathcal{Y}^{-}$we have

$$
\begin{equation*}
<\Phi_{f}^{+}, \Phi^{-}>=\int_{X(\omega)} f d m_{\Phi^{-}} \tag{66}
\end{equation*}
$$

Note that the mapping $\Xi_{\omega}^{+}: \operatorname{Lip}_{w}^{+}(X(\omega)) \rightarrow \mathcal{Y}_{\omega}^{+}$given by $\Xi_{\omega}^{+}(f)=\Phi_{f}^{+}$is continuous by definition and satisfies

$$
\begin{equation*}
\Xi_{\sigma \omega}^{+}\left(f \circ t_{\sigma}\right)=\left(t_{\sigma}\right)^{*} \Xi_{\omega}^{+}(f) . \tag{67}
\end{equation*}
$$

From the definitions we also have
Proposition $22 \operatorname{Let} \Phi^{+}(1), \ldots, \Phi^{+}(r)$ be a basis in $\mathcal{Y}_{\omega}^{+}$and let $\Phi^{-}(1), \ldots, \Phi^{-}(r)$ be the dual basis in $\mathcal{Y}_{\omega}^{-}$with respect to the pairing $<,>$. Then for any $f \in$ $\operatorname{Lip}_{w}^{+}(X(\omega))$ we have

$$
\Phi_{f}^{+}=\sum_{i=1}^{r}\left(m_{\Phi^{-}(i)}(f)\right) \Phi^{+}(i)
$$

### 3.3 Orderings and flows.

Assume that for $\mathbb{P}$-almost every $\omega$ a partial ordering $\mathfrak{o}(\omega)$ is given on $\mathcal{E}\left(\omega_{n}\right)$ for all $n \in \mathbb{Z}$ in such a way that edges starting at a given vertex are ordered linearly, while edges starting at different vertices are incomparable. Assume, moreover, that the orders $\mathfrak{o}(\omega)$ are $\sigma$-invariant, in the sense that the ordering $\mathfrak{o}(\omega)$ on $\mathcal{E}\left(\omega_{n}\right)$ is the same as the ordering $\mathfrak{o}(\sigma \omega)$ on $\mathcal{E}\left((\sigma \omega)_{n-1}\right)$.

Similarly to the above, construct spaces $X_{\mathfrak{o}}(\omega)$ and introduce a flow $h_{t}^{(+, \omega)}$ on each $X_{\mathfrak{o}}(\omega)$. The shift $\sigma$ renormalizes the flows $h_{t}^{(+, \omega)}$ : if we set

$$
\begin{equation*}
H^{(1)}(n, \omega)=\|\mathbb{A}(n, \omega)\| \tag{68}
\end{equation*}
$$

then for any $t \in \mathbb{R}$ we have a commutative diagram


As before, each measure $\Phi^{+} \in \mathcal{Y}_{\omega}^{+}$yields a Hölder cocycle over the flow $h_{t}^{(+, \omega)}$; we shall denote the cocycle by the same letter as the measure.

Note that for any $\Phi^{-} \in \mathcal{Y}_{\omega}^{-}$the measure $m_{\Phi^{-}}$defined by (62) satisfies

$$
\left(h_{t}^{(+, \omega)}\right)_{*} m_{\Phi^{-}}=m_{\Phi^{-}},
$$

similarly to G. Forni's invariant distributions [5] 6].
Note that the mapping $\Xi_{\omega}^{+}: \operatorname{Lip}_{\omega}^{+}(X(\omega)) \rightarrow \mathcal{Y}_{\omega}^{+}$given by $\Xi_{\omega}^{+}(f)=\Phi_{f}^{+}$by definition satisfies

$$
\begin{equation*}
\Xi_{\omega}^{+}\left(f \circ h_{t}^{(+, \omega)}\right)=\Xi_{\omega}^{+}(f) \tag{69}
\end{equation*}
$$

We thus have the following

Theorem 7 Let $\mathbb{P}$ be an ergodic $\sigma$-invariant probability measure on $\Omega$ satisfying the assumptions 1, 园, For any $\varepsilon>0$ there exists a positive constant $C_{\varepsilon}$ depending only on $\mathbb{P}$ such that the following holds. For $\mathbb{P}$-almost any $\omega$ there exists a continuous mapping $\Xi_{\omega}^{+}: \operatorname{Lip}_{w}^{+}(X(\omega)) \rightarrow \mathcal{Y}_{\omega}^{+}$such that for any $f \in \operatorname{Lip}_{w}^{+}(X(\omega))$, any $x \in X(\omega)$ and all $T>0$ we have

$$
\left|\int_{0}^{T} f \circ h_{t}^{(+, \omega)}(x) d t-\Xi_{\omega}^{+}(f)(x, t)\right| \leq C_{\varepsilon}\|f\|_{{L i p_{w}^{+}}^{+}}\left(1+T^{\varepsilon}\right) .
$$

The mapping $\Xi_{\omega}^{+}$satisfies the equality $\Xi_{\omega}^{+}\left(f \circ h_{t}^{(+, \omega)}\right)=\Xi_{\omega}^{+}(f)$. The diagram

is commutative.
The mapping $\Xi_{\omega}^{+}$is given by $\Xi_{\omega}^{+}(f)=\Phi_{f}^{+}$, where $\Phi_{f}^{+}$is defined by (66).
Now assume that the second Lyapunov exponent $\theta_{2}$ of the renormalization cocycle $\mathbb{A}$ is positive and simple. Let $v_{2} \in E_{\omega}^{+}$be a Lyapunov vector corresponding to the $\operatorname{exponent} \exp \left(\theta_{2}\right)$ (such a vector is defined up to multiplication by a scalar). Introduce a multiplicative cocycle $H^{(2)}(n, \omega)$ over $\sigma$ by the formula

$$
\begin{equation*}
H^{(2)}(n, \omega)=\frac{\left|\mathbb{A}(n, \omega) v_{2}^{(\omega)}\right|}{\left|v_{2}^{(\omega)}\right|} \tag{70}
\end{equation*}
$$

Recall that the cocycle $H^{(1)}(n, \omega)$ is given by (68). Similarly to the above, given a bounded measurable function $f: X(\omega) \rightarrow \mathbb{R}$ and $x \in X(\omega)$, introduce a continuous function $\mathfrak{S}_{n}[f, x]$ on the unit interval by the formula

$$
\begin{equation*}
\mathfrak{S}_{n}[f, x](\tau)=\int_{0}^{\tau H^{(1)}(n, \omega)} f \circ h_{t}^{(+, \omega)}(x) d t \tag{71}
\end{equation*}
$$

The functions $\mathfrak{S}_{n}[f, x]$ are $C[0,1]$-valued random variables on the probability space $\left(X(\omega), \nu_{\omega}\right)$.

Theorem 8 Let $\mathbb{P}$ be an ergodic $\sigma$-invariant probability measure on $\Omega$ satisfying the assumptions 1, 园 and such the second Lyapunov exponent of the renormalization cocycle $\mathbb{A}$ with respect to $\mathbb{P}$ is positive and simple.

For $\mathbb{P}$-almost any $\omega^{\prime} \in \Omega$ there exists a non-degenerate compactly supported measure $\eta_{\omega^{\prime}}$ on $C[0,1]$ and, for $\mathbb{P}$-almost any pair $\left(\omega, \omega^{\prime}\right)$ there exists a sequence of moments $l_{n}=l_{n}\left(\omega, \omega^{\prime}\right)$ such that the following holds.

For $\mathbb{P}$-almost any $\omega$ there exists a continuous functional

$$
\mathfrak{a}^{(\omega)}: \operatorname{Lip}_{w}^{+}(X(\omega)) \rightarrow \mathbb{R}
$$

such that for $\mathbb{P}$-almost any $\omega^{\prime}$ and any $f \in \operatorname{Lip}_{w, 0}^{+}(X(\omega))$ satisfying $\mathfrak{a}^{(\omega)}(f) \neq 0$ the sequence of random variables

$$
\frac{\mathfrak{S}_{l_{n}\left(\omega, \omega^{\prime}\right)}[f, x]}{\mathfrak{a}^{(\omega)}(f) H^{(2)}\left(l_{n}\left(\omega, \omega^{\prime}\right), \omega\right)}
$$

converges in distribution to $\eta_{\omega^{\prime}}$ as $n \rightarrow \infty$.
Theorems 7, 8 imply Theorems 3, 4, The proofs of Theorems 7, 8 follow the same pattern as those of Theorems [5, 6) detailed proofs will appear in the sequel to this paper.

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