Simple Lyapunov spectrum for certain linear cocycles over partially hyperbolic maps*

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Abstract

Criteria for the simplicity of the Lyapunov spectra of linear cocycles have been found by Furstenberg, Guivarc’h–Raugi, Gol’dsheid–Margulis and, more recently, Bonatti–Viana and Avila–Viana. In all the cases, the authors consider cocycles over hyperbolic systems, such as shift maps or Axiom A diffeomorphisms.

In this paper we propose to extend such criteria to situations where the base map is just partially hyperbolic. This raises several new issues concerning, among others, the recurrence of the holonomy maps and the (lack of) continuity of the Rokhlin disintegrations of \(u\)-states.

Our main results are stated for certain partially hyperbolic skew-products whose iterates have bounded derivatives along center leaves. They allow us, in particular, to exhibit non-trivial examples of stable simplicity in the partially hyperbolic setting.

Keywords: Lyapunov exponent, partial hyperbolicity, linear cocycles
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1. Introduction

The theory of linear cocycles is now a classical field of dynamical systems and ergodic theory, grounded on the pioneer works of Furstenberg, Kesten \([11, 13]\) and Oseledets \([19]\). The

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derivatives of smooth dynamical systems are the first examples that come to mind, but the notion of linear cocycle is a lot more broad, and arises naturally in many other situations, e.g. in the spectral theory of Schrödinger operators.

Among the outstanding issues is the problem of simplicity: when is it the case that the dimension of all Oseledets subspaces is equal to 1? This was first studied by Furstenberg [11], Guivarc’h–Raugi [15] and Gol’dsheid–Margulis [14], who obtained explicit simplicity criteria for random i.i.d. products of matrices. Recently, Bonatti–Viana [7] and Avila–Viana [2] extended the theory to include a much broader class of (Hölder continuous) cocycles over hyperbolic maps. There is also much progress in the quasi-periodic case, that is, for linear cocycles over rotations: see Duarte–Klein [10] and references therein.

Our purpose in this paper is to initiate the study of the simplicity problem in the context of linear cocycles over partially hyperbolic maps, that combine features from both the hyperbolic and the quasi-periodic cases.

The theory of partially hyperbolic diffeomorphisms and flows was initiated by Brin–Pesin [8] and Hirsch–Pugh–Shub [16] and has been at the heart of much recent progress in dynamical systems. While boasting many of the important properties of uniformly hyperbolic (Axiom A) systems, partially hyperbolic maps are a lot more flexible and encompass several interesting new phenomena.

Linear cocycles over volume-preserving partially hyperbolic maps were studied previously by Avila–Santamaria–Viana [1]. The issue of simplicity is much better understood when the base map is non-uniformly hyperbolic, meaning that all the center Lyapunov exponents are non-zero. Indeed, Viana [22] proved that simplicity is generic, in a very strong sense, among 2-dimensional cocycles. Backes–Poletti–Varandas [5] extended that conclusion to any dimension $d \geq 2$, under additional assumptions such as fiber-bunching.

For this reason, here we focus on the opposite case, namely, we take the partially hyperbolic map to be mostly neutral along the center direction, meaning that its iterates have bounded derivatives along the leaves of the center foliation. The following simple example illustrates some of the systems we have in mind.

Let $\omega_0, \omega_1$ be real numbers and $f_0, f_1 : S^1 \to S^1$ be the corresponding rotations, that is, $f_i(t) = t + \omega_i \mod Z$ for every $t \in S^1$. Take $\omega_0$ to be irrational. Let $A_0 : S^1 \to \text{SL}(3, \mathbb{R})$ and $A_1 : S^1 \to \text{SL}(3, \mathbb{R})$ be given by

\[ A_0(t) = \begin{pmatrix} 2 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 2 \end{pmatrix} \text{ and } A_1(t) = R_1(t)R_2(t)R_3(t) \]

where $R_i(t)$ denotes the rotation of angle $2\pi t$ around the $i$th axis. Each $A_i$ defines a linear cocycle $F_i$ over the transformation $f_i$. We want to consider the random combination $\hat{F}$ of these two cocycles: at each step one applies either $F_1$ or $F_2$, at random.

The results in this paper (see theorem A and example 2.4) ensure that the Lyapunov spectrum of $\hat{F}$ is simple, and the same is true for any small perturbation in the uniform topology. Concerning this last point, it should be noted that simplicity of the Lyapunov spectrum is usually not an open property, see Wang and You [25].

2. Definitions and statements

Here we state our main result. Beforehand, we must give the precise definitions of the notions involved in the statement. In what follows, $\mathbb{K}$ denotes either the real field $\mathbb{R}$ or the complex field $\mathbb{C}$, indifferently.
2.1. Linear cocycles and Lyapunov exponents

The linear cocycle defined by a measurable matrix-valued function $\hat{A} : \hat{M} \to \text{GL}(d, \mathbb{R})$ over an invertible measurable map $\hat{f} : \hat{M} \to \hat{M}$ is the (invertible) map $\hat{F}_A : \hat{M} \times \mathbb{R}^d \to \hat{M} \times \mathbb{R}^d$ given by

$$\hat{F}_A(\hat{p}, v) = (\hat{f}(\hat{p}), \hat{A}(\hat{p})v).$$

Its iterates are given by $\hat{F}_A^n(\hat{p}, v) = (\hat{f}^n(\hat{p}), \hat{A}^n(\hat{p})v)$ where

$$\hat{A}^n(\hat{p}) = \begin{cases} 
\hat{A}(\hat{f}^{n-1}(\hat{p}))(\hat{A}(\hat{f}^{n-2}(\hat{p}))\ldots \hat{A}(\hat{f}(\hat{p}))) & \text{if } n > 0 \\
\text{id} & \text{if } n = 0 \\
\hat{A}(\hat{f}^{-1}(\hat{p}))(\hat{A}(\hat{f}^{-2}(\hat{p}))\ldots \hat{A}(\hat{f}^{-n}(\hat{p})))^{-1} & \text{if } n < 0.
\end{cases}$$

Let $\hat{\mu}$ be an $\hat{f}$-invariant probability measure on $\hat{M}$ such that $\log \|\hat{A}^{\pm 1}\|$ are integrable. By Oseledets [19], at $\hat{\mu}$-almost every point $\hat{p} \in \hat{M}$ there exist real numbers $\lambda_1(\hat{p}) > \cdots > \lambda_k(\hat{p})$ and a decomposition $\mathbb{R}^d = E_{\hat{p}}^1 \oplus \cdots \oplus E_{\hat{p}}^k$ into vector subspaces such that

$$\hat{A}(\hat{p})E_{\hat{p}}^i = E_{\hat{p}}^i$$

and $\lambda_i(\hat{p}) = \lim_{|n| \to \infty} \frac{1}{n} \log \|\hat{A}^n(\hat{p})v\|$ for every non-zero $v \in E_{\hat{p}}^i$ and $1 \leq i \leq k$. The dimension of $E_{\hat{p}}^i$ is called the multiplicity of $\lambda_i(\hat{p})$.

In this work we assume that the invariant measure $\hat{\mu}$ is ergodic. Then the Lyapunov exponents and the dimensions of the subspaces $E_{\hat{p}}^i$ are constant almost everywhere. The Lyapunov spectrum of the cocycle is the set of all Lyapunov exponents. The following notion is central to the whole paper: the Lyapunov spectrum is simple if it contains exactly $d$ distinct Lyapunov exponents or, equivalently, if every Lyapunov exponent has multiplicity equal to 1.

2.2. Partially hyperbolic skew-products

Let $\hat{\sigma} : \hat{\Sigma} \to \hat{\Sigma}$ be any two-sided finite or countable shift. By this we mean that $\hat{\Sigma}$ is the set of two-sided sequences $(x_n)_{n \in \mathbb{Z}}$ in some set $X \subset \mathbb{N}$ with $X > 1$, and the map $\hat{\sigma}$ is given by

$$\hat{\sigma}((x_n)_{n \in \mathbb{Z}}) = (x_{n+1})_{n \in \mathbb{Z}}.$$ 

Let $\text{dist}_\Sigma : \hat{\Sigma} \times \hat{\Sigma} \to \mathbb{R}$ be the distance defined by

$$\text{dist}_\Sigma(\hat{x}, \hat{y}) = \sum_{k = -\infty}^{\infty} 2^{-|k|} \delta(x_k, y_k), \quad \text{with } \hat{x} = (x_k)_{k \in \mathbb{Z}} \text{ and } \hat{y} = (y_k)_{k \in \mathbb{Z}},$$

where $\delta(x, y) = 1$ if $x = y$ and $\delta(x, y) = 0$ otherwise. Then $\hat{\sigma}$ is a hyperbolic homeomorphism (in the sense of [22]), as we are going to explain.

Given any $\hat{x} \in \hat{\Sigma}$, we define the local stable and unstable sets of $\hat{x}$ with respect to $\hat{\sigma}$ by

$$W^s_{\text{loc}}(\hat{x}) = \{ \hat{y} : x_k = y_k \text{ for every } k \geq 0 \}$$

and

$$W^u_{\text{loc}}(\hat{x}) = \{ \hat{y} : x_k = y_k \text{ for every } k \leq 0 \}.$$ 

From now on we fix $\lambda = 1/2$ and $\tau = 1/2$. Then,

(i) $\text{dist}_\Sigma(\hat{\sigma}^n(\hat{y}_1), \hat{\sigma}^n(\hat{y}_2)) \leq \lambda^n \text{dist}_\Sigma(\hat{y}_1, \hat{y}_2)$ for any $\hat{y} \in \hat{\Sigma}$, $\hat{y}_1, \hat{y}_2 \in W^u_{\text{loc}}(\hat{y})$ and $n \geq 0$;

(ii) $\text{dist}_\Sigma(\hat{\sigma}^{-n}(\hat{y}_1), \hat{\sigma}^{-n}(\hat{y}_2)) \leq \lambda^n \text{dist}_\Sigma(\hat{y}_1, \hat{y}_2)$ for any $\hat{y} \in \hat{\Sigma}$, $\hat{y}_1, \hat{y}_2 \in W^s_{\text{loc}}(\hat{y})$ and $n \geq 0$;
gives rise, also through inducing, to a linear cocycle associated to a transition matrix \( \hat{f} \) and \( \hat{f}^\mu \) for every \( \hat{x} \) and \( \hat{y} \), where \( \hat{x} \) and \( \hat{y} \) are \( \hat{x} \)-invariant measures. All the conditions that follow are not affected by this procedure.

Moreover, every linear cocycle \( F \) of the form \( \sigma \) is dominated in the spirit of [16]: it means that any expansion and contraction of \( f_i \) along the fibers \( \{ \hat{x} \} \times K \) are dominated by the hyperbolicity of the base map \( \hat{\sigma} \). For our purposes, its main relevance is that it ensures the existence of strong-stable and strong-unstable ‘foliations’ for \( f \), as we explain next.

Let the product \( \hat{M} = \hat{\Sigma} \times K \) be endowed with the distance defined by

\[
\text{dist}_{\hat{M}}((\hat{x}_1, t_1), (\hat{x}_2, t_2)) = \text{dist}_{\hat{\Sigma}}(\hat{x}_1, \hat{x}_2) + \text{dist}_K(t_1, t_2),
\]

(iii) if \( \text{dist}_{\hat{\Sigma}}((\hat{x}, \hat{y})) \leq \tau \), then \( \hat{W}_{\text{loc}}^u(\hat{x}) \) and \( \hat{W}_{\text{loc}}^s(\hat{y}) \) intersect in a unique point, which is denoted by \([x, y]\) and depends continuously on \( \hat{x} \) and \( \hat{y} \).

By a partially hyperbolic skew-product over the shift map \( \hat{\sigma} \) we mean a homeomorphism \( \hat{f} : \hat{\Sigma} \times K \to \hat{\Sigma} \times K \) of the form

\[
\hat{f}(\hat{x}, t) = (\hat{\sigma}(\hat{x}), \hat{f}_i(t))
\]

where \( K \) is a compact Riemannian manifold and the maps \( \hat{f}_i : K \to K \) are diffeomorphisms satisfying

\[
\lambda \|D\hat{f}_i(t)\| < 1 \quad \text{and} \quad \lambda \|D\hat{f}_i^{-1}(t)\| < 1 \quad \text{for every} \quad (\hat{x}, t) \in \hat{\Sigma} \times K,
\]

where \( \lambda \) is a constant as in (i) and (ii). We also assume the following Hölder condition: there exist \( C > 0 \) and \( \alpha > 0 \) such that the \( C^1 \)-distance between \( \hat{f}_i \) and \( \hat{f}_j \) is bounded by \( C \text{dist}_{\hat{\Sigma}}(\hat{x}, \hat{y})^\alpha \) for every \( \hat{x}, \hat{y} \in \hat{\Sigma} \).

We say that \( \hat{f} \) has mostly neutral center direction if the maps \( \hat{f}_i^n : K \to K \) defined for \( n \in \mathbb{Z} \) and \( \hat{x} \in \hat{\Sigma} \) by

\[
\hat{f}_i^n = \begin{cases} f_{i}^{n-1}(i) \circ \cdots \circ f_i & \text{if } n > 0 \\ \text{id} & \text{if } n = 0 \\ f_{\hat{\sigma}^{-1}(i)}^{-1} \circ \cdots \circ f_{\hat{\sigma}^{-1}(i)} \end{cases}
\]

have bounded derivatives, that is, if there exists \( C > 0 \) such that

\[
\|D\hat{f}_i^n\| \leq C \quad \text{for every} \quad \hat{x} \in \hat{\Sigma} \quad \text{and} \quad n \in \mathbb{Z}.
\]

**Remark 2.1.** Clearly, this implies that the \( \{ \hat{f}_i^n : j \in \mathbb{Z} \text{ and } \hat{x} \in \hat{\Sigma} \} \) is equicontinuous. When the maps \( \hat{f}_i \) are \( C^{1+r} \), equicontinuity alone suffices for all our purposes (see remark 3.2).

In the definition of partially hyperbolic skew-product, one may replace the shift \( \hat{\sigma} : \hat{\Sigma} \to \hat{\Sigma} \) with a sub-shift \( \hat{\sigma}_T : \hat{\Sigma}_T \to \hat{\Sigma}_T \) associated to a transition matrix \( T = (T_{ij})_{i,j \in \mathbb{X}} \). By this we mean that \( T_{ij} \in \{0, 1\} \) for every \( i,j \in \mathbb{X} \) and \( \hat{\sigma}_T \) is the restriction of the shift map \( \hat{\sigma} \) to the subset \( \hat{\Sigma}_T \) of sequences \( (\sigma_n)_{n \in \mathbb{Z}} \) such that \( T_{\sigma_n \sigma_{n+1}} = 1 \) for every \( n \in \mathbb{Z} \).

One way to reduce the sub-shift case to the full shift case is through inducing. Namely, fix any one-cylinder \( [0, i] = \{ (\sigma_n)_{n \in \mathbb{Z}} \in \hat{\Sigma}_T : x_0 = i \} \) with positive measure and consider the first return map \( g : [i] \to [i] \) of \( \hat{\sigma}_T \) to \([i]\). This is conjugate to a full countable shift (with the return times as symbols) and it preserves the normalized restriction to the cylinder of the \( \hat{\sigma}_T \)-invariant measure. All the conditions that follow are not affected by this procedure. Moreover, every linear cocycle \( F \) over \( \hat{\sigma}_T \) gives rise, also through inducing, to a linear cocycle over \( g \) whose Lyapunov spectrum is just a rescaling of the Lyapunov spectrum of \( F \). In particular, simplicity may also be read out from the induced cocycle.

2.3. Stable and unstable linear holonomies

Property (2) is a condition of domination (or normal hyperbolicity, in the spirit of [16]): it means that any expansion and contraction of \( f_i \) along the fibers \( \{ \hat{x} \} \times K \) are dominated by the hyperbolicity of the base map \( \hat{\sigma} \). For our purposes, its main relevance is that it ensures the existence of strong-stable and strong-unstable ‘foliations’ for \( f \), as we explain next.

Let the product \( \hat{M} = \hat{\Sigma} \times K \) be endowed with the distance defined by

\[
\text{dist}_{\hat{M}}(\hat{x}_1, t_1), (\hat{x}_2, t_2)) = \text{dist}_{\hat{\Sigma}}(\hat{x}_1, \hat{x}_2) + \text{dist}_K(t_1, t_2),
\]
where \( \text{dist}_K \) denotes the distance (1) on \( \hat{\Sigma} \) and \( \text{dist}_K \) is the distance induced by the Riemannian metric on \( K \).

We consider the stable holonomies
\[
h^s_{x,y} : K \rightarrow K, \quad h^s_{x,y} = \lim_{n \to \infty} (\hat{f}_x^n)^{-1} \circ \hat{f}_y^n,
\]
defined for every \( x \) and \( y \) with \( x \in W^s_{\text{loc}}(y) \), and unstable holonomies
\[
h^u_{x,y} : K \rightarrow K, \quad h^u_{x,y} = \lim_{n \to \infty} (\hat{f}_y^{-n})^{-1} \circ \hat{f}_x^n
\]
defined for every \( x \) and \( y \) with \( x \in W^u_{\text{loc}}(y) \). That these families of maps exist follows from the assumption (2), using arguments from [6]. See for instance [4], which deals with a similar setting.

We define the local strong-stable set and the local strong-unstable set of each \((\hat{x},t) \in \hat{M}\)
to be
\[
W^s_{\text{loc}}(\hat{x}, t) = \{ (\hat{y}, s) \in \hat{M} : \hat{y} \in W^s_{\text{loc}}(\hat{x}) \text{ and } s = h^s_{x,y}(t) \} \quad \text{and} \quad W^u_{\text{loc}}(\hat{x}, t) = \{ (\hat{y}, s) \in \hat{M} : \hat{y} \in W^u_{\text{loc}}(\hat{x}) \text{ and } s = h^u_{x,y}(t) \},
\]
respectively. It is easy to check that
\[
(\hat{y}, s) \in W^s_{\text{loc}}(\hat{x}, t) \Rightarrow \lim_{n \to +\infty} \text{dist}_{\hat{M}}(\hat{f}^n(\hat{y}, s), \hat{f}^n(\hat{x}, t)) = 0
\]
and analogously on strong-unstable sets for time \( n \to -\infty \).

### 2.4. Measures with partial product structure

Recall that \( \hat{M} = \hat{\Sigma} \times K \). Throughout, we take \( \hat{\mu} \) to be an \( \hat{f} \)-invariant measure with partial product structure, that is, a probability measure of the form \( \hat{\mu} = \hat{\rho} \mu^x \times \mu^u \times \mu^c \) where:

- \( \hat{\rho} : \hat{M} \rightarrow (0, +\infty) \) is a continuous function bounded away from zero and infinity;
- \( \mu^x \) is a probability measure supported on \( \Sigma^- = \mathbb{Z}^{<0} \);
- \( \mu^u \) is a probability measure supported on \( \Sigma^+ = \mathbb{Z}^{>0} \);
- \( \mu^c \) is a probability measure on the manifold \( K \).

For notational convenience, we formulate the boundedness condition as follows: there exists \( \kappa > 0 \) such that
\[
\frac{1}{\kappa} \leq \frac{\hat{\rho}(x^t, x^u)}{\hat{\rho}(\hat{x}, \hat{\epsilon}^u)} \leq \kappa \quad \text{and} \quad \frac{1}{\kappa} \leq \frac{\hat{\rho}(x^t, x^u)}{\hat{\rho}(\hat{\epsilon}^t, x^u)} \leq \kappa
\]
for every \( x^t, \hat{\epsilon}^u \in \Sigma^- \) and \( x^u, \hat{\epsilon}^t \in \Sigma^- \), where \( \hat{\rho} : \hat{\Sigma} \rightarrow \mathbb{R} \) is defined by
\[
\hat{\rho}(\hat{x}) = \int \hat{\rho}(\hat{x}, t) d\mu^c(t).
\]
Observe that when \( \hat{\Sigma} \) is a finite shift space this is an immediate consequence of compactness and the continuity of \( \hat{\rho} \).

Now define, for \( \hat{x} \in \hat{\Sigma} \),
\[
\hat{\rho}(\hat{x}, \cdot) = \frac{\hat{\rho}(\hat{x}, \cdot)}{\hat{\rho}(\hat{x})} \quad \text{and} \quad \hat{\mu}^x = \hat{\rho}(\hat{x}, \cdot) \mu^c.
\]
In other words, \( \check{\mu}^c \) is the normalization of \( \hat{\rho}(\check{x}, \cdot)\mu^c \). Note that \( \{ \check{\mu}^c : \check{x} \in \check{\Sigma} \} \) is a (continuous) disintegration of \( \check{\mu} \) along vertical fibers, that is, with respect to the partition \( \mathcal{P} = \{ \{ \check{x} \} \times K : \check{x} \in \check{\Sigma} \} \).

The assumption that \( \check{\mu} \) is invariant under \( \check{f} \), together with the fact that \( \check{\mu}^c \) depends continuously on \( \check{x} \), implies that
\[
(\hat{f}_{\check{x}})_* \check{\mu}^c = \check{\mu}^c(\check{x}) \quad \text{for every } \check{x} \in \check{\Sigma}.
\]
We will also see in section 5 that this disintegration is holonomy invariant:
\[
\begin{align*}
(\check{h}^s_{\check{x}})_* \check{\mu}^c &= \check{\mu}^c \quad \text{whenever } \check{y} \in W^s(\check{x}) \\
(\check{h}^u_{\check{x}})_* \check{\mu}^c &= \check{\mu}^c \quad \text{whenever } \check{y} \in W^u(\check{x}).
\end{align*}
\]

**Remark 2.2.** In particular, if \( \check{x} \) is a fixed point of the shift map then \( \check{\mu}^c \) is invariant under \( \check{f} \). Clearly, it is equivalent to \( \mu^c \). Moreover, if \( \check{y} \) is a homoclinic point of \( \check{x} \), that is, a point in \( W^s(\check{x}) \cap W^u(\check{x}) \), then \( (\check{h}^s_{\check{x}})_* \check{\mu}^c = (\check{h}^u_{\check{x}})_* \check{\mu}^c = \check{\mu}^c \).

### 2.5. Linear cocycles with holonomies

Let \( \check{A} : \check{M} \to \text{GL}(d, \mathbb{K}) \) be a \( \alpha \)-Hölder continuous map for some \( \alpha > 0 \). By this we mean that there exists \( C > 0 \) such that
\[
\| \check{A}(\check{p}) - \check{A}(\check{q}) \| \leq C \text{dist}_{\mathbb{K}}(\check{p}, \check{q})^\alpha \quad \text{for any } \check{p}, \check{q} \in \check{M}.
\]
The linear cocycle defined by \( \check{A} \) over the transformation \( \check{f} : \check{M} \to \check{M} \) is the map \( \check{F} : \check{M} \times \mathbb{K}^d \to \check{M} \times \mathbb{K}^d \) defined by
\[
\check{F}(\check{p}, v) = (\check{f}(\check{p}), \check{A}(\check{p})v).
\]
In what follows we take the cocycle to admit stable and unstable linear holonomies. Let us explain this.

By stable linear holonomies we mean a family of linear maps \( \check{H}^s_{\check{p}, \check{q}} : \mathbb{K}^d \to \mathbb{K}^d \), defined for each \( \check{p}, \check{q} \in \check{M} \) with \( \check{q} \in W^s_{\text{loc}}(\check{p}) \) and such that, for some constant \( L > 0 \),
\[
\begin{align*}
(a) \quad & \check{H}^s_{\check{p}, \check{q}}(j)(\cdot) = \check{A}^j(\check{q}) \circ \check{H}^s_{\check{p}, \check{q}}(0) \circ \check{A}^{-j}(\check{p})^{-1} \quad \text{for every } j \geq 1; \\
(b) \quad & \check{H}^s_{\check{p}, \check{p}} = \text{id} \quad \text{and} \quad \check{H}^s_{\check{z}, \check{p}} = \check{H}^s_{\check{z}, \check{q}} \circ \check{H}^s_{\check{p}, \check{q}} \quad \text{for any } \check{z} \in W^s(\check{p}); \\
(c) \quad & \| \check{H}^s_{\check{p}, \check{q}} - \text{id} \| \leq L \text{dist}_{\mathbb{K}}(\check{p}, \check{q})^\alpha; \\
(d) \quad & (\check{p}, \check{q}) \mapsto \check{H}^s_{\check{p}, \check{q}} \text{ is uniformly continuous on } \{ (\check{p}, \check{q}) : \check{q} \in W^s_{\text{loc}}(\check{p}) \} \subset \check{M} \times \check{M}.
\end{align*}
\]

Unstable linear holonomies \( \check{H}^u_{\check{p}, \check{q}} : \mathbb{K}^d \to \mathbb{K}^d \) are defined analogously, for the pairs \( (\check{p}, \check{q}) \in \check{M} \times \check{M} \) with \( \check{q} \in W^u_{\text{loc}}(\check{p}) \).

These notions were introduced in [1, 6], where they were called simply stable and unstable holonomies. We add the adjective linear to avoid any confusion with the holonomies \( h^s \) and \( h^u \) in the previous paragraph, that concern only the base dynamics, whereas \( h^s \) and \( h^u \) pertain to the linear cocycle.

It was shown in [1] that stable and unstable linear holonomies do exist, in particular, when the cocycle is fiber-bunched. By the latter we mean that there exist \( C > 0 \) and \( \theta < 1 \) such that
\[
\| \check{A}^n(\check{p}) \| \| \check{A}^n(\check{p})^{-1} \| \lambda^\theta \leq C \theta^n \quad \text{for every } \check{p} \in \check{M} \text{ and } n \geq 0,
\]
where $\lambda$ is a hyperbolicity constant for $\hat{f}$ as in conditions (i) and (ii) above. Then stable and unstable linear holonomies are given by

$$H_{\hat{p}\hat{q}}^s = \lim_{n \to \infty} \hat{A}^n(\hat{q})^{-1} \circ \hat{A}(\hat{p}), \quad \text{and} \quad H_{\hat{p}\hat{q}}^u = \lim_{n \to \infty} \hat{A}^{-n}(\hat{q})^{-1} \circ \hat{A}^{-n}(\hat{p}).$$

### 2.6. Pinching and twisting

Now we state our criterion for simplicity of the Lyapunov spectrum. It is assumed that the cocycle admits stable and unstable linear holonomies.

We call $\hat{F}$ pinching if there exists some fixed (or periodic) vertical leaf $\ell = \{\hat{x}\} \times K$ such that the restriction to $\ell$ of every external power $\Lambda^k \hat{F}$ has simple Lyapunov spectrum, relative to the $f_k$-invariant measure $\hat{\mu}_k^n$ (recall remark 2.2). In other words, the Lyapunov exponents $\lambda_1, \cdots, \lambda_d$ are such that, for each $1 \leq k \leq d - 1$ and $\hat{\mu}_k^n$-almost every $t \in K$, the sums

$$\lambda_{i_1}(\hat{x}, t) + \cdots + \lambda_{i_k}(\hat{x}, t), \quad 1 \leq i_1 < \cdots < i_k \leq d$$

are all distinct.

Next, take $\hat{F}$ to be pinching and let $\|\cdot\| = E^s(\hat{x}, t) \oplus \cdots \oplus E^u(\hat{x}, t)$ be the Oseledets decomposition at each point $(\hat{x}, t) \in K$. This is defined on a full $\hat{\mu}_k^n$-measure set. Choose (measurably) unit vectors $\hat{e}^s(\hat{x}, t) \in E^s(\hat{x}, t)$. Let $\hat{y}$ be a homoclinic point of $\hat{x}$. Given $t \in S^1$, denote $t_1 = h_{\hat{y}\hat{x}}(t)$ and $t_2 = h_{\hat{y}\hat{x}}(t_1)$. Then define $B(t)$ to be the matrix of the linear map

$$H_{\hat{y}(t_1),\hat{x}(t_2)} \circ H_{\hat{y}(t_2),\hat{x}(t_1)} : \mathbb{K}^d \to \mathbb{K}^d$$

relative to the bases $\{\hat{e}^s(t_2), \cdots, \hat{e}^u(t_2)\}$ and $\{\hat{e}^s(t), \cdots, \hat{e}^u(t)\}$, respectively. Observe that $t_2$ also varies on a full $\hat{\mu}_k^n$-measure set, since the composition of the holonomies preserves $\hat{\mu}_k^n$ (remark 2.2).

We call the cocycle $\hat{F}$ twisting if, for some choice of the homoclinic point $\hat{y}$, all the algebraic minors $m_{IJ}(t)$ of $B(t)$ are non-zero for $\hat{\mu}_k^n$-almost every $t \in K$ and they decay sub-exponentially along the orbits of $f_k$, meaning that

$$\lim_{n \to \infty} \frac{1}{n} \log |m_{IJ}(\hat{f}_k^n(t))| = 0 \quad \text{for} \quad \hat{\mu}_k^n \text{-almost every} \ t \in K$$

and any proper subsets $I$ and $J$ of $\{1, \ldots, d\}$.

**Remark 2.3.** It is well-known (see [18, proposition 2.2] or [23, corollary 3.11]) that the property (9) holds whenever the function $\log |m_{IJ}| \circ f_k - \log |m_{IJ}|$ is $\hat{\mu}_k^n$-integrable. In example 2.4 we show how to check the twisting condition in a specific case, using this observation.

Finally, we say that the cocycle $\hat{F}$ is simple if it is both pinching and twisting (in addition to admitting stable and unstable linear holonomies).

### 2.7. Main statement

Let $H^\alpha(M)$ denote the space of all $\alpha$-Hölder continuous maps $\hat{A} : \hat{M} \to \text{GL}(d, \mathbb{K})$. The norm

$$\|\hat{A}\|_\alpha = \sup_{\hat{p} \in \hat{M}} \|\hat{A}(\hat{p})\| + \sup_{\hat{p} \neq \hat{q}} \frac{\|\hat{A}(\hat{p}) - \hat{A}(\hat{q})\|}{\text{dist}_{\hat{M}}(\hat{p}, \hat{q})^\alpha}$$

defines a topology in $H^\alpha(M)$ that we call $\alpha$-Hölder topology.
We say that $A$ is a continuity point for the Lyapunov exponents if, for every $1 \leq i \leq d$, the function $\lambda_i : H^\alpha(M) \to \mathbb{R}$ is continuous in $A$.

**Theorem A.** Let $\hat{f} : M \to M$ be a partially hyperbolic skew-product with mostly neutral center direction and $\hat{\mu}$ be a $\hat{f}$-invariant measure with partial product structure. Suppose that $\hat{A} \in H^\alpha(M)$ is such that the corresponding linear cocycle $\hat{F} : \mathcal{M} \times \mathbb{R}^d \to \mathcal{M} \times \mathbb{R}^d$ over $\hat{f}$ is simple. Then $\hat{A}$ is a continuity point for the Lyapunov exponents, and the Lyapunov spectrum of $\hat{F}$ is simple.

By continuity, the Lyapunov spectrum remains simple for every perturbation of $\hat{F}$, that is, for the linear cocycle over $\hat{f}$ corresponding to every element of $H^\alpha(M)$ sufficiently close to $\hat{A}$.

**Example 2.4.** The example presented in the Introduction satisfies all the conditions in theorem A, and so the conclusion applies to it. In order to explain this, let us formalize the example as follows.

Let $\hat{\sigma} : \hat{\Sigma} \to \hat{\Sigma}$ be the shift map on $\hat{\Sigma} = \{0,1\}^\mathbb{Z}$ and let $\hat{\nu}$ be the Bernoulli measure $(\delta_0/2 + \delta_1/2)^\mathbb{Z}$. Let $\hat{f} : \hat{\Sigma} \times S^1 \to \hat{\Sigma} \times S^1$ be defined by $\hat{f}(\hat{\xi},t) = (\hat{\sigma}(\hat{\xi}), f_{0}\theta_0(t))$ and $\hat{\mu}$ be the product of $\hat{\nu}$ by the Haar measure on $S^1$. Finally, let $\hat{M} = \hat{\Sigma} \times S^1$ and $\hat{F} : \hat{M} \times \mathbb{R}^3 \to \hat{M} \times \mathbb{R}^3$ be given by $\hat{F}(\hat{\xi},t,v) = (\hat{f}(\hat{\xi},t), \hat{A}(\hat{\xi},t)v)$ with $\hat{A}(\hat{\xi},t) = A_{0\theta_0}(t)$.

It is clear that $\hat{f}$ is a skew-product with mostly neutral central direction, and $\hat{\mu}$ has partial product structure. Moreover, $\hat{\mu}$ is ergodic. Indeed, let $\hat{\zeta}$ be any ergodic component. Since $\hat{\nu}$ projects down to $\hat{\nu}$, which is ergodic, $\hat{\zeta}$ must project to $\hat{\nu}$. The Lyapunov exponent of $\hat{f}$ along the vertical $S^1$ fibers is zero and so, by the Invariance Principle of [3], there exists a disintegration $\{\hat{\zeta}_z : z \in \hat{\Sigma}\}$ of $\hat{\zeta}$ along the $S^1$ fibers which is invariant under stable and unstable holonomies and is continuous. Consider the fixed point $\hat{x} = (\ldots,0,0,0,\ldots)$ of $\hat{\sigma}$. Since $f_0$ is uniquely ergodic, because we took $\omega_0$ to be irrational, $\hat{\zeta}_z$ must coincide with the Haar measure on $S^1$. Then, by holonomy invariance, $\hat{\zeta}_z$ is uniquely determined at every point, which proves that the ergodic component is unique.

Now consider the homoclinic point $\hat{y} = (\ldots,0,1,0,\ldots)$ of $\hat{x}$, where the sole non-zero entry is in position 0. The corresponding stable and unstable holonomies are given by

$$h^s_{\hat{x},\hat{y}} = \lim_{n \to \infty} \left(f_0 \circ \cdots \circ f_1\right)^{-1} \circ \left(f_0 \circ \cdots \circ f_0\right) = f_1^{-1} \circ f_0 \quad \text{and}$$

$$h^u_{\hat{x},\hat{y}} = \lim_{n \to \infty} \left(f_0^{-1} \circ \cdots f_0^{-1}\right)^{-1} \circ \left(f_0^{-1} \circ \cdots f_0^{-1}\right) = \text{id}.$$

Similarly, the stable and unstable linear holonomies are given by

$$H^s_{(\hat{x},\hat{y}), (\hat{x},\hat{u})} = A_1(t)^{-1} \circ A_0 \quad \text{and} \quad H^u_{(\hat{x},\hat{y}), (\hat{x},\hat{u})} = \text{id}$$

where $s = h^s_{\hat{x},\hat{y}}(t) = f_0^{-1}(f_1(t)) = t + \omega_1 - \omega_0$.

It is also clear that $\hat{F}$ is pinching: its restriction to $\ell = \{\hat{x}\} \times S^1$ corresponds to the constant cocycle

$$A_0(t) = \begin{pmatrix} 2 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 2^{-t} \end{pmatrix},$$

whose Lyapunov spectrum is obviously simple. We are left to check the twisting condition.
Since $A_0$ is constant, so is its Oseledets decomposition $E^1(t) \oplus E^2(t) \oplus E^3(t)$, with $E^1 = \text{span}\{(1, 0, 0)\}$, $E^2 = \text{span}\{(0, 1, 0)\}$ and $E^3 = \text{span}\{(0, 0, 1)\}$. This shows that $\mathcal{B}(t)$ is just the matrix of

$$H^u_{(x, t), (x, t)} \circ H^s_{(x, t), (x, t)} = A_1(t)^{-1} \circ A_0 = R_3(-t)R_2(-t)R_1(-t)A_0,$$

relative to the canonical basis of $\mathbb{R}^3$. It is straightforward to check that all the minors $m_{I,J}(t)$ of this matrix are analytic functions of $t$ not identically zero. In particular, all their zeros have finite order and, consequently, the functions $\log |m_{I,J}|$ are integrable. It follows that

$$\lim_{n \to \infty} \frac{1}{n} \log |m_{I,J}(f^n(t))| = 0$$

for Lebesgue almost every $t \in S^1$,

which shows that $\hat{F}$ is twisting.

In many contexts of linear cocycles over hyperbolic systems, simplicity turns out to be a generic condition: it contains an open and dense subset of cocycles (precise statements can be found in Viana [23]). This is related to the fact that in the hyperbolic setting pinching and twisting are just transversality conditions, and so they clearly hold on the complement of suitable submanifolds with positive codimension.

It would be interesting to find whether this extends to the present partially hyperbolic setting. In dimension $d = 2$, simplicity is equivalent to positivity of the largest Lyapunov exponent and that has been shown to hold for an open and dense subset of linear cocycles over partially hyperbolic skew-products with mostly neutral center direction, by Poletti [20]. In general, by theorem A, it would suffice to prove density of our pinching and twisting conditions. Density of pinching corresponds, roughly, to density of simplicity for linear cocycles over quasi-periodic transformations, a subject that does not seem to have been much investigated beyond the 2-dimensional case (but see [10]). On the other hand, the arguments in example 2.4 suggest that twisting is probably a rather mild requirement on the cocycle.

2.8. Outline of the proof

For every $1 \leq \ell < d$, we want to find complementary $\hat{F}$–invariant measurable sections

$$\xi : \hat{M} \to \text{Grass}(l, d) \quad \text{and} \quad \eta : \hat{M} \to \text{Grass}(d - l, d)$$

(10)

such that the Lyapunov exponents of $\hat{F}$ along $\xi$ are strictly larger than the Lyapunov exponents along $\eta$.

The starting point is to reduce the problem to the case when the maps $\hat{f}_k$ and the matrices $\hat{A}(x, t)$ depend on $\hat{x}$ only through its positive part $x^+$. This we do in section 4, using the stable linear holonomies to conjugate the original dynamics to others with these properties. Then $\hat{f} : \hat{M} \to \hat{M}$ projects to a transformation $f : M \to M$ on $M = \Sigma^+ \times K$ which is a skew–product over the one-sided shift $\sigma : \Sigma^+ \to \Sigma^+$ and, similarly, the linear cocycle $\hat{F} : \hat{M} \times \mathbb{K}^d \to \hat{M} \times \mathbb{K}^d$ projects to a linear cocycle $F : M \times \mathbb{K}^d \to M \times \mathbb{K}^d$ over the transformation $\hat{f}$.

We also denote by $\tilde{F}$ and $F$ the actions

$$\tilde{F} : \hat{M} \times \text{Grass}(l, d) \to \hat{M} \times \text{Grass}(l, d)$$

and

$$F : M \times \text{Grass}(l, d) \to M \times \text{Grass}(l, d)$$
induced by the two linear cocycles on the Grassmannian bundles. Still in section 4, using very classical arguments, we relate the invariant measures of \( \hat{f} \) and \( \hat{F} \) with those of \( f \) and \( F \), respectively.

In section 5 we study \( u \)-states, that is, \( \hat{F} \)-invariant probability measures \( \hat{m} \) whose Rokhlin disintegrations \( \{ m_x : x \in \hat{\Sigma} \} \) are invariant under unstable holonomies, as well as the corresponding \( \hat{F} \)-invariant probability measures \( m \). Here we meet the first important new difficulty arising from the fact that \( \hat{f} \) is only partially hyperbolic. Indeed, in the hyperbolic setting this should extend to the disintegration Viana [2].

So, the previous arguments yield a statement. More precisely (theorem 7.1), there exists a measurable section \( \{ m_x : x \in M \} \) along the fibers \{x\} \times Grass(\(l,d\)) and this fact plays a key part in the arguments of Bonatti–Viana [7] and Avila–Viana [2].

In the partially hyperbolic setting, the situation is far more subtle: the disintegration \( \{ m_x : x \in \Sigma \} \) along the sets \{x\} \times K \times Grass(\(l,d\)) is still continuous, but there is no reason why this should extend to the disintegration \( \{ m_{x,t} : (x,t) \in M \} \) along the sets \{(x,t)\} \times Grass(\(l,d\)), which is what one really needs. The way we make up for this is by proving a kind of continuity: if \((x_t)_t \to x \) in \( \Sigma \) then \((m_{x,t})_t \to m_{x,t} \) in \( L^1(\mu^c) \). See proposition 5.8 for the precise statement.

This also leads to our formulating the arguments in terms of measurable sections \( K \to Grass(l,d) \) of the Grassmannian bundle, which is perhaps another significant novelty in this paper. The properties of such sections are studied in section 6. The key result (proposition 6.1) is that, under pinching and twisting, the graph of every invariant Grassmannian section has zero \( m_x \)-measure, for every \( x \in M \) and any \( u \)-state \( \hat{m} \).

These results build up to section 7, where we prove that every \( u \)-state \( \hat{m} \) has an atomic disintegration. More precisely (theorem 7.1), there exists a measurable section \( \xi : M \to Grass(l,d) \) such that, given any \( u \)-state \( \hat{m} \) on \( M \times Grass(l,d) \), we have

\[
\hat{m} \circ \xi^{-1} = \delta_{\xi(x,t)} \quad \text{for } \hat{\mu} \text{-almost every } (x,t) \in \hat{\Sigma}.
\]

Thus we construct the invariant section \( \xi : \hat{M} \to Grass(l,d) \) in (10).

To find the complementary invariant section \( \eta : \hat{M} \to Grass(d-l,d) \), in section 8 we apply the same procedure to the adjoint cocycle \( \hat{F}^* \), that is, the linear cocycle defined over \( \hat{f}^{-1} : M \to M \) by the function

\[
(\hat{x},t) \mapsto \hat{A}^*(x,t) = \text{adjoint of } \hat{A}(\hat{f}^{-1}(\hat{x},t)).
\]

We check (proposition 8.4) that this cocycle \( \hat{F}^* \) is pinching and twisting if and only if \( \hat{f} \) is. So, the previous arguments yield a \( \hat{F}^* \)-invariant section \( \xi^* : M \to Grass(l,d) \) related to the \( u \)-states of \( \hat{F}^* \). Then we just take \( \eta = (\xi^*)^{-1} \).

Finally, in section 9 we check that the eccentricity, or lack of conformality, of the iterates \( \hat{A}^n \) goes to infinity \( \hat{\mu} \)-almost everywhere (see proposition 9.1) and we use this fact to deduce that every Lyapunov exponent of \( \hat{F} \) along \( \xi \) is strictly larger than any of the Lyapunov exponents of \( \hat{F} \) along \( \eta \). At this stage the arguments are again very classical. This concludes the proof of theorem A.

In appendix we show that continuous maps are dense in the corresponding \( L^1 \) space, whenever the target space is geodesically convex. This is probably well known, but we could not find explicit references.

3. Disintegration along center leaves

Let us start by fixing some terminology. We use \( \text{id}_Y \) to denote the identity transformation in a set \( Y \). Similarly, \( \text{dist}_Y \) will always denote the distance in a metric space \( Y \).
Let \( \hat{\Sigma} = \Sigma^{-} \times \Sigma^{+} \), where \( \Sigma^{-} = X^{\geq 0} \) and \( \Sigma^{+} = X^{\leq 0} \). Thus we write every \( \hat{x} \in \hat{\Sigma} \) as \( (x', x^u) \) with \( x' \in \Sigma^{-} \) and \( x^u \in \Sigma^{+} \). For simplicity, we also write \( \Sigma = \Sigma^{+} \) and \( x = x^u \). Let \( P : \Sigma \to \Sigma \) be the canonical projection given by \( P(\hat{x}) = x \) and let \( \sigma : \Sigma \to \Sigma \) be the one-sided shift. Given points \( \hat{x} \in \hat{\Sigma} \) and \( \hat{q} \in \hat{M} \), denote
\[
x_n = P(\hat{\sigma}^{-n}(\hat{x})) \quad \text{and} \quad q_n = (P \times \text{id}_K)(\hat{f}^{-n}(\hat{q}))
\]
for each \( n \in \mathbb{N} \).

We also consider \( M = \Sigma \times K \) and the projection \( \mu = (P \times \text{id}_K)_* \hat{\mu} \). In other words,
\[
\mu = \rho(x, t) \mu^u \times \mu^c \quad \text{where} \quad \rho(x, t) = \int \hat{\rho}(x', x, t) \, d\mu^u(x').
\]
(12)

Similarly to (5), for each \( x \in M \) let
\[
\varrho(x, \cdot) = \frac{\rho(x, \cdot)}{\int \rho(x, t) \, d\mu^c(t)} \quad \text{and} \quad \mu^c_x = \varrho(x, \cdot) \mu^c.
\]
(13)

Note that \( \{ \mu^c_x : x \in \Sigma \} \) is a continuous disintegration of \( \mu \) with respect to the partition \( P = \{ \{x\} \times K : x \in \Sigma \} \).

In this section we derive some useful properties of these disintegrations (5) and (13). For this, we assume that the base dynamics is such that each \( \hat{f}_s : K \to K \) along the center direction depends only on \( x = P(\hat{x}) \). This is no restriction in our setting, as we will see in section 4.

Then there exists \( f : M \to M \) of the form \( f(x, t) = (\sigma(x), f_s(t)) \) such that
\[
(P \times \text{id}_K) \circ \hat{f} = f \circ (P \times \text{id}_K).
\]

3.1. Holonomy invariance

We call the extremal center Lyapunov exponents of \( \hat{f} \) the limits
\[
\lambda^c(\hat{x}, t) = \lim_n \frac{1}{n} \log \| D\hat{f}_n^c(t) \| \quad \text{and} \quad \lambda^c(\hat{x}, t) = \lim_n \frac{1}{n} \log \| D\hat{f}_n^c(t)^{-1} \|
\]
for \( (\hat{x}, t) \in \hat{\Sigma} \times K \). The Oseledets theorem [19] ensures that these numbers are well defined at \( \hat{\mu} \)-almost every point. In our situation, since the maps \( \hat{f}_n^c \) have uniformly bounded derivatives:

**Lemma 3.1.** \( \lambda^c = \lambda^c = 0 \).

**Remark 3.2.** When the maps \( \hat{f}_n^c \) are \( C^{1+\epsilon} \), equicontinuity alone suffices to get the conclusion of lemma 3.1. This can be shown using Pesin theory, as follows.

Suppose that \( \lambda^c > 0 \). Then we have a Pesin unstable manifold defined \( \hat{\mu} \)-almost everywhere. This implies that there exist \( \hat{x} \in \hat{\Sigma} \) and \( t \neq s \in K \) such that
\[
\text{dist}_K(\hat{f}_n^c(t), \hat{f}_n^c(s)) \to 0.
\]
Then, given points \( t \) and \( s \) in the unstable manifold and given any \( \delta > 0 \), there exists \( n \) such that \( \text{dist}_K(\hat{f}_n^c(t), \hat{f}_n^c(s)) < \delta \). This implies that the family is not equicontinuous. The proof for \( \lambda^{-c} \) is analogous.

Let \( \pi_1 : \hat{M} \to \hat{\Sigma} \) be the projection \( \pi_1(\hat{x}, t) = \hat{x} \), recall that we also assume that \( \hat{f} : \hat{M} \to \hat{M} \) admits \( s \)-holonomies and \( u \)-holonomies and \( \hat{\mu} \) has partial product structure. That implies that \( (\pi_1)_* \hat{\mu} \) has local product structure in the sense of [3].
Lemma 3.3. The map $\hat{x} \mapsto \hat{\mu}^\epsilon$ is continuous, and the disintegration $\{\hat{\mu}^\epsilon : \hat{x} \in \hat{\Sigma}\}$ is both $u$-invariant and $s$-invariant:

(a) $\left(\hat{h}_{\hat{U},\hat{x}}\right)^\epsilon \hat{\mu}^\epsilon = \hat{\nu}^\epsilon$ for every $\hat{x} \in W^u(\hat{y})$ and

(b) $\left(\hat{h}_{\hat{U},\hat{x}}\right)^\epsilon \hat{\mu}^\epsilon = \hat{\nu}^\epsilon$ for every $\hat{x} \in W^s(\hat{z})$.

Proof. By theorem D in [3], there exists a disintegration $\{\hat{\mu}^\epsilon : \hat{x} \in \hat{\Sigma}\}$ which is continuous, $u$-invariant and $s$-invariant. By essential uniqueness, $\hat{\mu}^\epsilon = \hat{\nu}^\epsilon$ for $\hat{\mu}$-almost every $x$. Since both disintegrations are continuous, it follows that they coincide, and so $\{\hat{\mu}^\epsilon : \hat{x} \in \hat{\Sigma}\}$ is continuous, $u$-invariant and $s$-invariant, as claimed.

Corollary 3.4. $\mu^\epsilon = \mu^\epsilon$ for every $\hat{x} \in \hat{\Sigma}$, where $x = P(\hat{x})$.

Proof. The assumption that $\hat{f}_k$ only depends on $x = P(\hat{x})$ implies that $\hat{h}_{\hat{U},\hat{x}}=\text{id}_k$ for every $\hat{x}$ and $\hat{y}$ in the same stable set. By the previous lemma, this implies that $\hat{\mu}^\epsilon = \hat{\nu}^\epsilon$ whenever $\hat{x}$ and $\hat{y}$ are in the same stable set. Then,

$$\hat{\mu}^\epsilon = \int_{\hat{\nu}^\epsilon} \hat{\mu}^\epsilon(\hat{y}) = \hat{\nu}^\epsilon$$

for any $\hat{x}$ with $P(\hat{x}) = x$.

We also have

Corollary 3.5. The disintegration $\{\mu^\epsilon : x \in \Sigma\}$ is $f$-invariant, in the sense that $(f_\sigma)_\epsilon \mu^\epsilon = \mu_{f\sigma(x)}$ for every $x \in \Sigma$.

Proof. We have that $(f_\sigma)_\epsilon \hat{\mu}^\epsilon = \hat{\mu}^\epsilon$$ for $\hat{\mu}$-almost every $\hat{x}$, because $\hat{\mu}$ is $f$-invariant. Since $\hat{x} \mapsto \hat{\mu}^\epsilon$ is continuous, the identity extends to every $\hat{x} \in \hat{\Sigma}$. By corollary 3.4 this implies that $(f_\sigma)_\epsilon \mu^\epsilon = \mu_{f\sigma(x)}$ for every $x \in \Sigma$.

3.2. Jacobians

Denote $\hat{\nu} = (\pi_1)_\epsilon \hat{\mu}$ and $\nu = (\pi_1)_\epsilon \mu$ where $\pi_1$ denotes both canonical projections $\hat{M} \to \hat{\Sigma}$ and $M \to \Sigma$. Recall the functions $\hat{\phi}$ and $\phi$ defined in (5) and (13). Note that $\{\hat{\phi}(x)\hat{\nu}^\epsilon : \hat{x} \in \hat{\Sigma}\}$ is a disintegration of $\hat{\nu}$ with respect to the partition $\{\pi_1^{-1}(\hat{x}) : \hat{x} \in \hat{\Sigma}\}$ of $\hat{M}$ and $\{\phi(x)\mu^\epsilon : x \in \Sigma\}$ is a disintegration of $\mu$ with respect to the partition $\{\pi_1^{-1}(x) : x \in \Sigma\}$ of $\Sigma$.

Given a measurable map $g : N \to N$ and a measure $\eta$ on $N$ we call Jacobian of $g$ with respect to $\eta$ the essentially unique function $J_{\eta g} : N \to \mathbb{R}$ such that $\eta(g(C)) = \int_C J_{\eta g} \, d\eta$ for every measurable set $C \subset N$ where $g$ is invertible. This is well defined whenever $N$ can be covered with countably many domains of invertibility of $g$. See [24, section 9.7] for a detailed discussion.

Remark 3.6. Let $J_{f_\epsilon'} : K \to \mathbb{R}$ be the Jacobian of $f_\epsilon'$ with respect to $\mu^\epsilon$. Using the observation that $\{\phi(x)\mu^\epsilon : x \in \Sigma\}$ is a disintegration of the $f$-invariant measure $\mu$, one easily gets that

$$J_{f_\epsilon'}(t) = \frac{\phi(f_\epsilon'(x,t))}{\phi(x,t)}.$$
In particular, these Jacobians are uniformly bounded from above and below. Analogously, the fact that \( \{ \hat{\mu}_k^x : x \in \Sigma \} \) is invariant under stable and unstable holonomies ensures that the Jacobians \( J^*_h \) of those holonomies with respect to \( \mu^c \) are uniformly bounded from above and below.

**Lemma 3.7.** \( J^*_h \sigma^k(x, t) = J^*_h \sigma^k(x) \) for every \((x, t) \in M \) and \( k \geq 1 \).

**Proof.** Fix some \( k \)-cylinder \( I = [0, x_0, \ldots, x_{k-1}] \) and let \( J \subset I \) and \( C \subset K \). Noting that \( \sigma^k \mid J \) is injective,

\[
\mu(f^k(J \times C)) = \int_{y \in \sigma^1(J)} \mu_c(f^{k}_{\sigma^1(y)}(C)) \, d\nu(y) = \int_{z \in J} J^*_h \sigma^k(z) \mu_c(f^{k}_{\sigma^1(z)}(C)) \, d\nu(z),
\]

where \( z(y) \) is the unique point in \( \sigma^{-k}(y) \cap J \) and we use the change of variables \( z = z(y) \). Using lemma 3.5, it follows that

\[
\mu(f^k(J \times C)) = \int_{J} J^*_h \sigma^k(z) \mu_c(C) \, d\nu(z) = \int_{J \times C} (J^*_h \sigma^k \circ \pi_1) \, d\mu,
\]

which concludes the proof. \( \square \)

Now we find the Jacobian of \( \sigma^k \):

**Lemma 3.8.** \( J^*_h \sigma^k(x) = 1/\nu^*_c(\sigma^k(I)) \), where \( I = [-k; x_0, \ldots, x_{k-1}] \). Consequently, the Jacobians \( J^*_h \sigma^k = J^*_h \sigma^k \circ \pi_1 \) are continuous and bounded from away zero and infinity on every \( k \)-cylinder.

**Proof.** Given \( x \in \Sigma \) and \( n \geq 1 \), let \( J_n = [x_0, \ldots, x_n] \) be the \( n \)-cylinder that contains \( x \). Then,

\[
J^*_h \sigma^k(x) = \lim_{n \to \infty} \frac{\nu(\sigma^k(J_n))}{\nu(J_n)}.
\]

Since \( \hat{\nu} \) is invariant under \( \hat{\sigma} \),

\[
\nu(J_n) = \hat{\nu}(\Sigma^- \times J_n) = \hat{\nu}(\hat{\sigma}^k(\Sigma^- \times J_n)) = \int_{y \in \sigma^1(J_n)} \nu^*_c(I) \, d\nu(y).
\]

It follows that

\[
\frac{1}{J^*_h \sigma^k(x)} = \lim_{n \to \infty} \frac{\int_{y \in \sigma^1(J_n)} \nu^*_c(I) \, d\nu(y)}{\nu(\sigma^k(J_n))} = \nu^*_c(I).
\]

This proves the first part of the conclusion. The second part is a consequence, since the local product structure implies that \( x \mapsto \nu^*_c(I) \) is continuous for every cylinder \( I \). \( \square \)

4. **Convergence of conditional measures**

For each \( 1 \leq l < d \), the linear cocycle \( \hat{F} : \hat{M} \times \mathbb{R}^d \to \hat{M} \times \mathbb{R}^d \) induces a projective cocycle \( \hat{F} : \hat{M} \times \text{Grass}(l, d) \to \hat{M} \times \text{Grass}(l, d) \) through
\[ \tilde{F}(\tilde{q}, \tilde{v}) = (\tilde{f}(\tilde{q}), \tilde{A}(\tilde{q})\tilde{v}). \]  

Let \( \mu = (P \times \text{id}_K)_* \tilde{\mu} \) and, for any Borel probability measure \( \tilde{m} \) on \( M \times \text{Grass}(l, d) \),

\[ m = (P \times \text{id}_K \times \text{id}_{\text{Grass}(l, d)})_* \tilde{m}. \]

We will be especially interested in the case when \( \tilde{m} \) is a \( \tilde{F} \)-invariant probability measure that projects down to \( \tilde{\mu} \) under the canonical projection \( \pi : M \times \text{Grass}(l, d) \to M \) on the first coordinate.

### 4.1. Reduction to the one-sided case

Our first step is to show that, up to conjugating the cocycle in a suitable way, we may suppose that:

(A) the base dynamics \( \tilde{f}_t \) along the center direction depends only on \( x \);

(B) the matrix \( \tilde{A}(\tilde{x}, t) \) depends only on \( (x, t) \).

Next, let us explain how such a conjugacy may be defined using the stable linear holonomies.

Let \( x^e \in \Sigma^e \) be fixed. For any \( \tilde{y} \in \Sigma \), let \( \phi(\tilde{y}) = (x^e, y) \) and then define

\[ h(\tilde{y}, t) = \left( \tilde{y}, h^e_{\phi(\tilde{y})}(t) \right) \].

Then \( \tilde{f} = h^{-1} \circ f \circ h \) is given by

\[ \tilde{f}(\tilde{y}, t) = (\tilde{\sigma}(\tilde{y}), \tilde{f}_t(t)), \quad \text{with} \quad \tilde{f}_t(t) = h^e_{\phi(\tilde{y})}(t) \circ f(\tilde{y}, t). \]

Notice that \( \tilde{f}_t \) does depend only on \( y \) (because \( \phi \) does).

Assume that (A) is satisfied. Define \( \hat{\phi}(\tilde{y}, t) = (\phi(\tilde{y}), t) \) and then let

\[ H(\tilde{y}, t) = H^e_{\phi(\tilde{y})}(t). \]

Define \( \hat{A}(\tilde{y}, t) = H(\tilde{f}(\tilde{y}, t))^{-1} \circ A(\tilde{y}, t) \circ H(\tilde{y}, t) \). Then

\[ \hat{A}(\tilde{y}, t) = H^e_{\tilde{f}(\tilde{y}, t)} \circ A(\phi(\tilde{y}), t). \]

which only depends on \( (y, t) \). Clearly, this procedure does not affect the Lyapunov exponents.

Observe that this new cocycle is continuous but not necessarily Hélder continuous, since the definition involves the holonomies which are not assumed to be Hélder continuous.

However, \( \hat{A} \) does clearly have well defined holonomies, which is what we actually need for the following.

From now on, we assume that both (A) and (B) are satisfied. Then, there exist

\[ f : M \to M, \quad f(x, t) = (\sigma(x), f_t(t)) \quad \text{and} \quad A : M \to \text{GL}(d, \mathbb{K}) \]

such that

\[ (P \times \text{id}_K) \circ \hat{f} = f \circ (P \times \text{id}_K) \quad \text{and} \quad \hat{A} = A \circ (P \times \text{id}_K). \]

Consequently, the map

\[ F : M \times \text{Grass}(l, d) \to M \times \text{Grass}(l, d), \quad F(p, V) = (f(p), A(p)V) \]

satisfies

\[ (P \times \text{id}_K \times \text{id}_{\text{Grass}(l, d)}) \circ \hat{F} = F \circ (P \times \text{id}_K \times \text{id}_{\text{Grass}(l, d)}). \]
The following well known basic fact will be used to characterize the $F$-invariant probability measures:

**Proposition 4.1.** Let $(N, \mathcal{B}, \eta)$ be a Lebesgue probability space and $g : N \to N$ be a measurable map that preserves $\eta$. Let $\{\eta_y : y \in N\}$ be the disintegration of $\eta$ with respect to the partition into pre-images $\mathcal{P} = \{g^{-1}(y) : y \in N\}$. Let

$$G : N \times L \to N \times L, \quad G(x, v) = (g(x), G_v(x))$$

be a measurable skew-product over $g$ and, given any probability measure $m$ on $N \times L$ that projects down to $\eta$, let $\{m_x : x \in N\}$ be its disintegration with respect to the partition into vertical fibers $\{x \times L, x \in N\$. Then $m$ is invariant under $G$ if and only if

$$m_x = \int \hat{G}_z \cdot m_y \, d\eta(z) \quad \text{for $\eta$-almost every } x \in N.$$

As an immediate consequence, we get:

**Corollary 4.2.** In the conditions of proposition 4.1, if $g$ is invertible then $m$ is invariant under $G$ if and only if $m_x = (G_{g^{-1}(x)})_\ast m_{g^{-1}(x)}$ for $\eta$-almost every $x \in N$.

**Proof.** Each $\eta_y$ must coincide with the Dirac mass at $g^{-1}(y)$.

4.2. Lifting of measures

The next proposition shows that every $\hat{F}$-invariant measure $\hat{m}$ that projects down to $\hat{m}$ may be recovered from the corresponding $F$-invariant measure $m$, defined by (15). Recall that we write $q_n = (P \times \text{id}_K)(\hat{f}^{-n}(\hat{q}))$ for each $\hat{q} \in \hat{M}$ and $n \geq 0$.

The following proposition is borrowed from [7, section 3]:

**Proposition 4.3.** Take $\hat{m}$ to be $\hat{F}$-invariant. Then, for $\hat{\mu}$-almost every $\hat{q} \in \hat{M}$, the sequence $(\hat{A}^n(q_n)_\ast m_{q_n})_n$ converges to $\hat{m}_\hat{q}$ in the weak* topology.

Moreover, for any $k \geq 1$ and any choice of points $y_{n,k}$ such that $f^k(y_{n,k}) = q_n$ and $\{y_{n,k} : n \geq 0\}$ is contained in some set of the form $I_k \times K$, where $I_k$ is a $k$-cylinder in $\Sigma$,

$$\lim_{n \to \infty} A^n(q_n)_\ast m_{q_n} = \lim_{n \to \infty} A^{n+k}(y_{n,k})_\ast m_{y_{n,k}}. \quad (16)$$

Adapting the proof to the present setting is straightforward, so we only make a couple of observations. Firstly, the proof uses the fact that the Jacobians $J_{f^k(y_{n,k})}$ are bounded away from zero. In [7] that is automatic, by continuity and compactness of the ambient space, whereas in the present setting it is ensured by the assumption that $y_{n,k} \in I_k \times K$ for every $n \geq 0$, through lemma 3.8. Finally, note that the identity (16) corresponds to the conclusion (b) in [7, proposition 3.1].

5. Properties of $u$-states

A probability measure $\hat{m}$ on $\hat{M} \times \text{Grass}(I, d)$ is called a $u$-state of $\hat{F}$ if there exists a disintegration $\{\hat{m}_{\hat{q}} : \hat{q} \in \hat{M}\}$ along the partition $\{\{\hat{q}\} \times \text{Grass}(I, d) : \hat{q} \in \hat{M}\}$ which is invariant under unstable linear holonomy: there exists $M \subset M$ with $\hat{\mu}(M) = 1$ such that
\[ \hat{m}_\hat{q} = H_{\hat{p}_\hat{q}}^\text{loc} m_\rho \] for every \( \hat{p}, \hat{q} \in \hat{M} \) with \( \hat{q} \in W^u_{\text{loc}}(\hat{p}) \). \hspace{1cm} (17)

Let \( \pi : \hat{M} \times \text{Grass}(l, d) \to \hat{M} \) be the canonical projection.

**Proposition 5.1.** There is some \( F \)-invariant \( u \)-state \( \hat{m} \) that projects down to \( \hat{\mu} \) under \( \pi \).

This is analogous to [2, proposition 4.2]. In very brief terms, the idea is to fix some \( \hat{x} \in \hat{\Sigma} \) and to construct a homeomorphism between the space of measures in \( \{ \hat{x} \} \times W^u_{\text{loc}}(\hat{x}) \times K \) that project down to \( \mu' \) and the space of all \( u \)-states. Using that the former is weak* compact, we get that the space of \( u \)-states measures is also compact. Moreover, it is \( F \)-invariant. That ensures that any accumulation point of \( n^{-1} \sum_{j=0}^{n-1} F_j \hat{m} \) is also a \( u \)-state.

In the remainder of this section, \( \hat{m} \) denotes any \( F \)-invariant \( u \)-state that projects down to \( \hat{\mu} \) under \( \pi \), and \( \{ \hat{m}_\hat{q} : \hat{q} \in \hat{M} \} \) is taken to be a disintegration as in (17).

### 5.1. Bounded distortion

Let \( \pi_1 : \hat{M} \to \hat{\Sigma} \) be the canonical projection \( \pi_1(\hat{x}, t) = \hat{x} \) and denote \( \hat{\nu} = \pi_1 \hat{\mu} \). Equivalently,

\[
\hat{\nu}(E) = \int_{E \times K} \hat{\rho}(x', x, t) \, d\mu'(x') \, d\mu^u(x) \, d\mu^c(t)
\]

for any measurable set \( E \subset \hat{\Sigma} \). For each \( x \in \Sigma \), define \( \hat{\nu}_x \) to be the normalization of \( \mu^x \)

\[
\hat{\nu}_x = \int \hat{\rho}(\cdot, x, t) \, d\mu^c(t).
\]

Then \( \{ \hat{\nu}_x : x \in \Sigma \} \) is a continuous disintegration of \( \hat{\nu} \) with respect to the partition into local stable sets \( W^s_{\text{loc}}(x) \).

The measure \( \hat{\nu} \) satisfies the properties of local product structure, boundedness and continuity in [2, section 1.2]. In what follows, we recall a few results about this type of measures that we will use later. For each \( \hat{x}^u \in \Sigma^+ \) and \( k \geq 1 \) let the backward average measure \( \mu^u_{k, \hat{x}^u} \) of the map \( \sigma \) be defined by

\[
\mu^u_{k, \hat{x}^u} = \sum_{\sigma^k(z) = \hat{x}^u} \frac{1}{J\sigma^k(z)} \delta_z
\]

where \( J\sigma^k : \Sigma^+ \to \mathbb{R} \) is the Jacobian of \( \mu^u \) with respect to \( \sigma^k \).

**Lemma 5.2 (Lemma 2.6 in [2]).** For any cylinder \( I^u = [0, t_0, \ldots, t_{k-1}] \subset \Sigma^+ \) and \( \hat{z}^u \in I^u \),

\[
\hat{\delta}_k^{\hat{\nu}^u} \hat{\nu}^u = J\sigma^k(\hat{z}) \left( \hat{\nu}_{\sigma^k(\hat{z})} \right)(I^u)
\]

where \( \{ \hat{\nu}_{\sigma^k(\hat{z})} : z^u \in \Sigma^+ \} \) is the disintegration of \( \hat{\nu} \) with respect to the partition \( \{ \Sigma^- \times \{ z^u \} : z^u \in \Sigma^+ \} \).

**Lemma 5.3 (Lemma 2.7 in [2]).** For every cylinder \( [J] \subset \Sigma^+ \) and \( \hat{x}^u \in \Sigma^+ \),

\[
\kappa^u([J]) \geq \limsup_n \frac{1}{n} \sum_{k=0}^{n-1} \mu^u_{k, \hat{x}^u}([J])
\]

\[
\geq \liminf_n \frac{1}{n} \sum_{k=0}^{n-1} \mu^u_{k, \hat{x}^u}([J]) \geq \frac{1}{\kappa^u([J])},
\]

\[ \frac{n-1}{n} \sum_{k=0}^{n-1} \mu^u_{k, \hat{x}^u}([J]) \geq \frac{1}{\kappa^u([J])}, \]
where \( \kappa \) is the bound given in (3).

As a direct consequence, for every cylinder \([J] \subset \Sigma^+\) and \( x^\mu \in \Sigma^+ \),
\[
\limsup_k \mu^\mu_{x^\mu}(|J|) \geq \frac{1}{\kappa} \mu^\mu(|J|).
\] (18)

5.2. \( L^1 \)-continuity of conditional probabilities

For every \( x \in \Sigma \) let
\[
F_x : K \times \text{Grass}(l,d) \to K \times \text{Grass}(l,d), \quad F_x(t,V) = (f_x(t), A(x,t)V)
\]
and for every \( x, y \in \tilde{\Sigma} \) in the same unstable set let
\[
H_{x,y} : K \times \text{Grass}(l,d) \to K \times \text{Grass}(l,d)
\]
be defined by
\[
H_{x,y}(t,V) = (h_{x,y}^t(t), H_{x,y}(l,t)(y)).
\]
Observe that \( F_y \circ H_{x,y} = H_{x,y}(x,y) \circ F_x \). Now define \( \tilde{m}_x : x \in \tilde{\Sigma} \) as
\[
\tilde{m}_x = \int \tilde{m}_x \, d\tilde{\mu}_x(t).
\]
Observe that for any \( \varphi : K \times \text{Grass}(l,d) \to \mathbb{R} \),
\[
\int \varphi \, d(H_{x,y} \, \tilde{m}_x) = \int \varphi(h_{x,y}^t(t), H_{x,y}(l,t)(y)) \, d\tilde{m}_x \, d\tilde{\mu}_x(t)
\]
\[
= \int \varphi(h_{x,y}^t(t), V) \, d\left(H_{x,y}(l,t)(y) \, \tilde{m}_x \, \tilde{\mu}_x(t)\right)
\]
\[
= \int \varphi(h_{x,y}^t(t), V) \, d\tilde{m}_x \, d\tilde{\mu}_x(t)
\]
\[
= \int \varphi(t, V) \, d\tilde{m}_x \, d\tilde{\mu}_x(t)
\]
\[
= \int \varphi(t, V) \, d\tilde{m}_x \, d\tilde{\mu}_x(t),
\]
because \( \tilde{m} \) is a u-state and \( \{\tilde{\mu}_x : x \in \tilde{\Sigma}\} \) is \( h^x \)-invariant. It is also easy to see that \( \{\tilde{m}_x : x \in \tilde{\Sigma}\} \) is a disintegration of \( \tilde{m} \) with respect to the partition \( \{x \times K \times \text{Grass}(l,d) : x \in \tilde{\Sigma}\} \).

The main point with the next corollary is that the conclusion is for every \( x \in \Sigma \).

**Corollary 5.4.** For every u-state \( \tilde{m} \), there exists a disintegration \( \{\tilde{m}_x : x \in \tilde{\Sigma}\} \) of \( \tilde{m} \) such that
\[
\tilde{m}_{x^\otimes(i)} = F_{x^\otimes(i)}^n \tilde{m}_x
\]
for every \( n \geq 1 \), every \( x \in \Sigma \), and \( \tilde{\nu}_\otimes \)-almost every \( x \in W^u_{loc}(x) \).

**Proof.** Up to modifying a 0 measure subset we can assume that de disintegration \( \{\tilde{m}_x : x \in \tilde{\Sigma}\} \) has the property that \( H_{x,y} \, \tilde{m}_x = \tilde{m}_y \) for every \( x \in W^u_{loc}(y) \). Since \( \tilde{m} \) is \( F \)-invariant, the equality is true for all \( n \geq 1 \) and \( \tilde{\nu}_\otimes \)-almost all \( \tilde{z} \in \tilde{\Sigma} \) or, equivalently, for \( \tilde{\nu}_\otimes \)-almost every \( \tilde{z} \in W^u_{loc}(z) \) and \( \nu \)-almost every \( z \in \Sigma \). Consider an arbitrary point \( x \in \Sigma \). Since \( \nu \) is positive
on open sets, \( x \) may be approximated by points \( z \) such that

\[
m_{\delta^n(z)} = F_{\delta^n(z)}^* \hat{m}.
\]

for every \( n \geq 1 \) and \( \hat{\mu}_z \)-almost every \( z \in W^s_{\text{loc}}(\hat{x}) \). Since the conditional probabilities of \( \hat{m} \) are invariant under unstable linear holonomies, it follows that

\[
m_{\delta^n(z)} = (H_{\delta^n(z)}, \delta^n(z))_\# F_{\delta^n(z)}^* \hat{m} = F_{\delta^n(z)}^* (H_{\delta^n(z)} \circ \hat{m} = F_{\delta^n(z)}^* \hat{m}.
\]

for \( \hat{\mu}_z \)-almost every \( z \in W^s_{\text{loc}}(\hat{x}) \), where \( \hat{x} \) is the unique point in \( W^s_{\text{loc}}(x) \cap W^a_{\text{loc}}(\hat{x}) \). Since the measures \( \hat{\mu}_z \) and \( \hat{\mu}_z \) are equivalent, this is the same as saying that the last equality holds for \( \hat{\mu}_z \)-almost every \( \hat{x} \in W^s_{\text{loc}}(x) \), as claimed. \( \square \)

Recall that

\[
\mu = (P \times \text{id}_K)_\# \hat{\mu} \quad \text{and} \quad m = (P \times \text{id}_K \times \text{id}_{\text{Grass}(l,d)})_\# \hat{m}.
\]

Let \( \{m_x : x \in \Sigma\} \) and \( \{m_{x,t} : (x, t) \in M\} \) be disintegrations of \( m \) with respect to the partitions \( \{\{x\} \times K \times \text{Grass}(l,d), x \in \Sigma\} \) and \( \{\{x, t\} \times \text{Grass}(l,d), (x, t) \in \Sigma \times K\} \), respectively. Thus each \( m_x \) is a probability measure on \( K \times \text{Grass}(l,d) \) and each \( m_{x,t} \) is a probability measure on \( \text{Grass}(l,d) \).

It is easy to check that \( x \mapsto m_x \) may be chosen to be continuous with respect to the weak* topology (see corollary 5.9). The corresponding statement for \( x \mapsto m_{x,t} \) is false, in general. However, the main goal in this section is to show that the family \( \{m_{x,t} : (x, t) \in M\} \) does have some continuity property:

**Proposition 5.5.** Let \( \{x_n\}_n \) be a sequence in \( \Sigma \) converging to some \( x \in \Sigma \). Then there exists a sub-sequence \( \{x_{n,k}\}_k \) such that

\[
m_{x_{n,k}} \rightarrow m_{x,t} \text{ as } k \rightarrow \infty
\]

in the weak* topology, for \( \mu^t \)-almost every \( t \in K \).

We will deduce this from a somewhat stronger \( L^1 \)-continuity result, whose precise statement will be given in proposition 5.8. The key ingredient in the proofs is a result about maps on geodesically convex metric spaces that we are going to state in lemma 5.6 and which will also be useful at latter stages of our arguments.

A metric space \( N \) is **geodesically convex** if there exists \( \tau \geq 1 \) such that for every \( u, v \in N \) there exist a continuous path \( \lambda : [0, 1] \rightarrow N \) with \( \lambda(0) = u \), \( \lambda(1) = v \) and

\[
dist_N(\lambda(t), \lambda(s)) \leq \tau \text{dist}(u, v) \text{ for every } s, t \in [0, 1].
\]

Geodesically convex metric spaces include convex subsets of a Banach space, path connected compact metric spaces and complete connected Riemannian manifolds, among other examples. The spaces of maps with values in a geodesically convex metric space are analyzed in appendix.

**Lemma 5.6.** Let \( L \) be a geodesically convex metric space and take \( (K, \mathcal{B}_K, \mu_K) \) to be a probability space such that \( K \) is a normal topological space, \( \mathcal{B}_K \) is the Borel \( \sigma \)-algebra of \( K \) and \( \mu_K \) is a regular measure.

Let \( H_{j,t} : L \rightarrow L \) and \( h_j : K \rightarrow K \), with \( j \in \mathbb{N} \) and \( t \in K \), be such that

\[
(H_{j,t}(x))_j \rightarrow x \quad \text{and} \quad (h_j(t))_j \rightarrow t,
\]

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uniformly in \( t \in K \) and \( x \in L \) and, moreover, the Jacobian \( J h_j(t) \) of each \( h_j \) with respect to \( \mu_K \) is uniformly bounded. Then

\[
\lim_j \int \text{dist}_L(\psi(t), H_j \circ \psi \circ h_j(t)) \, d\mu_K(t) = 0
\]

for every bounded measurable map \( \psi : K \to L \).

**Proof.** Take \( j \in \mathbb{N} \) to be sufficiently large that \( d_L(H_j(x), x) < \epsilon/4 \) for every \( t \) and \( x \). Then,

\[
\int \text{dist}_L(\psi, H_j \circ \psi \circ h_j) \, d\mu_K \leq \int \left( \text{dist}_L(\psi, \psi \circ h_j) + \text{dist}_L(\psi \circ h_j, H_j \circ \psi \circ h_j) \right) \, d\mu_K \leq \int \text{dist}_L(\psi, \psi \circ h_j) \, d\mu_K + \frac{\epsilon}{4}.
\]

Let \( C > 1 \) be a uniform bound for \( J h_j \). By proposition A.1, given \( \epsilon > 0 \) there exists a continuous map \( \tilde{\psi} : K \to L \) such that

\[
\int \text{dist}_L(\tilde{\psi}, \psi) \, d\mu_K < \frac{\epsilon}{4C}.
\]

Then, by change of variables,

\[
\int \text{dist}_L(\tilde{\psi} \circ h_j, \psi \circ h_j) \, d\mu_K \leq C \int \text{dist}_L(\tilde{\psi}, \psi) \, d\mu_K < \frac{\epsilon}{4}.
\]

Then

\[
\int \text{dist}_L(\tilde{\psi}, \psi \circ h_j) \, d\mu_K \leq \int \left( \text{dist}_L(\tilde{\psi}, \tilde{\psi} \circ h_j) + \text{dist}_L(\tilde{\psi} \circ h_j, \psi \circ h_j) \right) \, d\mu_K \leq \int \text{dist}_L(\tilde{\psi}, \tilde{\psi} \circ h_j) \, d\mu_K + \frac{\epsilon}{2}.
\]

By the continuity of \( \tilde{\psi} \), increasing \( j \) if necessary,

\[
\text{dist}_L(\tilde{\psi}(t), \psi \circ h_j(t)) < \frac{\epsilon}{4} \quad \text{for every} \ t \in K.
\]

The conclusion follows from these inequalities. \( \square \)

**Lemma 5.7.** Let \((x_n)\) be a sequence in \( \Sigma \) converging to some \( x \in \Sigma \) and \((j_n)\) be a sequence of integer numbers such that \( z_n = \sigma^{-j_n}(x_n) \) converges to some \( z \in \Sigma \) and \((f_{j_n}^z)\) converges uniformly to some \( g : K \to K \). Then \( g \) is absolutely continuous with respect to \( \mu^c \) with bounded Jacobian. Moreover, the Jacobians of \( f_{j_n}^z \) with respect to \( \mu^c \) are uniformly bounded.

**Proof.** By lemma 3.5 we have that \((f_{j_n}^z)_* \mu^c_{z_n} = \mu^c_{x_n} \). Taking \( n \to \infty \) we get that \( g_* \mu^c_z = \mu^c_x \), which implies that \( J_{\mu^c} g = (x,t)/g(z,t) \) is uniformly bounded. Since, \( J_{\mu^c} f_{j_n}^z \) converges uniformly to \( J_{\mu^c} g \), it follows that the sequence is uniformly bounded. \( \square \)
**Proposition 5.8.** Let \( \varphi : \text{Grass}(l, d) \to \mathbb{R} \) be a continuous function, \((x_n)_n\) be a sequence in \( \Sigma \) converging to some \( x \in \Sigma \) and \((j_n)_n\) be a sequence of integer numbers such that \( \sigma^{-j_n}(x_n) \) converges and \((f^j_{x_n})_n\) converges uniformly to some \( g : K \to K \). Then \( \int \varphi \, dm_{x,j^j_{x_n}(t)} \) converges to \( \int \varphi \, dm_{x,g(t)} \) in \( L^1(\mu^c) \).

**Proof.** Denote \( t_n = f^j_{x_n}(t) \). Fix \( x' \in \Sigma^- \) and let

\[
h^n_u = h^n_u(\xi', x_n, t_n) \circ f^j_{x_n} \quad \text{and} \quad H^n_u = H^n_u(\xi', x_n, t_n) \circ f^j_{x_n},
\]

Let \( \mathcal{M} \) be the space of probability measures on \( \text{Grass}(l, d) \) with the distance

\[
d(\xi, \eta) = \sup \left\{ \left| \int \phi \, d\xi - \int \phi \, d\eta \right| : \sup \left| \phi \right| \leq 1 \right\}.
\]

This generates the weak* topology, and so \( \mathcal{M} \) is compact. By remark 3.6 and lemma 5.7, the Jacobians of \( \tilde{g}^{-1} \circ h^n_u \) with respect to \( \mu^c \) are uniformly bounded. Applying lemma 5.6 with \( L = \mathcal{M} \), \( H_{j,\xi} = (H^n_u)_{*, \xi} \), \( h_j = \tilde{g}^{-1} \circ h^n_u \) and \( \psi(t) = \hat{m}_{\xi', x_n, t_n} \), we get that

\[
\lim_{n \to \infty} \int d \left( \hat{m}_{\xi', x_n, t_n}, (H^n_u)_{*, \xi} \right) \hat{m}_{\xi', x_n, t_n} \, d\mu_j(t) = 0.
\]

Observe that \( (H^n_u)_{*, \xi} \hat{m}_{\xi', x_n, t_n} = \hat{m}_{\xi', x_n, t_n} \) and so the previous relation implies that the sequence \( t \mapsto \int \varphi \, d\hat{m}_{\xi', x_n, t} \) converges to \( t \mapsto \int \varphi \, d\hat{m}_{\xi', x_n, t} \) in \( L^1(\mu^c) \).

Next, by the definition of the disintegration,

\[
m_{\xi, t} = \int \tilde{\rho}(x', x, t) \hat{m}_{\xi', x_n, t} \, d\mu^c(x')
\]

and so

\[
\int \left| \int \varphi(\xi') \, d\hat{m}_{\xi', x_n, t} - \int \varphi(\xi') \, d\hat{m}_{\xi', x_n, t} \right| \, d\mu^c
\leq \int \left| \int \varphi(\xi') \, d\hat{m}_{\xi', x_n, t} \right| \, d\mu^c
\]

\[
- \int \varphi(\xi', x, g(t)) \, d\hat{m}_{\xi', x_n, t} \, d\mu^c
\]

So, noting that the integrand goes to zero as \( n \to \infty \), for every \( x' \in \Sigma^- \), the dominated convergence theorem ensures that

\[
\lim_{n \to \infty} \int \left| \int \varphi(\xi') \, d\hat{m}_{\xi', x_n, t} - \int \varphi(\xi') \, d\hat{m}_{\xi', x_n, t} \right| \, d\mu^c = 0,
\]

as we wanted to prove. \( \square \)

The case when \( j_n = 0 \) for every \( n \in \mathbb{N} \) suffices for proving proposition 5.5 (the full statement will be needed in section 7). Indeed, it gives that if \( (x_n)_n \to x \) and \( \varphi : \text{Grass}(l, d) \to \mathbb{R} \) is continuous then \( t \mapsto \int \varphi \, dm_{x_n, t} \) converges to \( t \mapsto \int \varphi \, dm_{x, t} \) in \( L^1(\mu^c) \). So, there exists a subsequence \( (n_k)_k \) such that

\[
\int \varphi \, dm_{x_n, t} \to \int \varphi \, dm_{x, t} \quad \text{for} \quad \mu^c = \text{almost every} \ t.
\]
Moreover, since the space of continuous functions is separable, one can use a diagonal argument (see e.g. the proof of [24, proposition 2.1.6]) to construct such a sub-sequence independent of $\varphi$. In other words,

$$m_{\varphi,k} \to m_{\varphi}$$

in the weak* topology, for $\mu^c$-almost every $t$.

This proves proposition 5.5.

**Corollary 5.9.** The disintegration $\{m_x : x \in \Sigma\}$ is continuous.

**Proof.** Let $\varphi : K \times \text{Grass}(l,d) \to \mathbb{R}$ be a continuous function. Given any $(x_n)_n \to x$, we have that

$$\left| \int \varphi \, dm_{x_n} - \int \varphi \, dm_{x} \right|$$

$$= \left| \int \int \varphi(t,v) \, dm_{x_n}(v) \rho(x_n,t) \, d\mu^c(t) \right|$$

$$- \left| \int \int \varphi(t,v) \, dm_{x}(v) \rho(x,t) \, d\mu^c(t) \right|$$

$$\leq \int \left| \int \varphi(t,v) \rho(x_n,t) \, dm_{x_n}(v) \right| \, d\mu^c(t) - \left| \int \int \varphi(t,v) \rho(x,t) \, dm_{x}(v) \right| \, d\mu^c(t).$$

By proposition 5.5, up to restricting to a subsequence, we may suppose that $(m_{x_n})_n$ converges to $m_x$ in the weak* sense, for $\mu^c$-almost every $t$. Then

$$\int \varphi(t,v) \rho(x_n,t) \, dm_{x_n}(v) \to \int \varphi(t,v) \rho(x,t) \, dm_{x}(v)$$

for $\mu^c$-almost every $t$. To get the conclusion it suffices to use this observation in the previous inequality, together with dominated convergence.

**Corollary 5.10.** We have $m_x = \int (F_{\varphi}^k)_{*} m_x \, d\nu^k(x)$ for every $x \in \Sigma$ and $k \geq 1$, where $\nu^k_x$ is defined as

$$\nu^k_x = \sum_{y \in \rho^{-1}(x)} \frac{1}{J_{\rho,\sigma^k}(y)} \delta_y.$$

**Proof.** The $F$-invariance of $m$ gives that $m_x = \int (F_{\varphi}^k)_{*} m_x \, d\nu^k(x)$ for $\nu$-almost every $x$, the continuity of the disintegration implies that this extends to every $x \in \Sigma$.

**6. Dual graphs of Grassmannian sections**

Fix $1 \leq l < d$. Let $w_1, \ldots, w_l$ be a basis of a given subspace $W \in \text{Grass}(l,d)$. The exterior product $w_1 \wedge \cdots \wedge w_l$ depends on the choice of the basis, but its projective class does not. Thus we have a well defined map

$$\text{Grass}(l,d) \leftrightarrow \mathbb{P} \Lambda^l(K^d), \quad W \mapsto [w_1 \wedge \cdots \wedge w_l],$$

(20)

which can be checked to be an embedding: it is called the Plücker embedding of Grass$(l,d)$. The image is the projectivization of the space of $l$-vectors.
\[
\Lambda^l_\xi(\mathbb{R}^d) = \{ w_1 \wedge w_2 \wedge \cdots \wedge w_l \in \Lambda^l(\mathbb{R}^d) : w_i \in \mathbb{R}^d \text{ for } 1 \leq i \leq l \},
\]
which we denote by \( \mathbb{P}\Lambda^l_\xi(\mathbb{R}^d) \). This is a closed subset of \( \mathbb{P}\Lambda^l(\mathbb{R}^d) \), and it is invariant under the action induced on \( \mathbb{P}\Lambda^l(\mathbb{R}^d) \) by any linear map \( B : \mathbb{R}^d \rightarrow \mathbb{R}^d \). See [2, section 2] for more information about \( l \)-vectors.

The geometric hyperplane \( \mathcal{H}V \subset \text{Grass}(l, d) \) associated to each \( V \in \text{Grass}(d - l, d) \) is the set \( \mathcal{H}V \) of all subspaces \( W \) in \( \text{Grass}(l, d) \) which are not in general position relative to \( V \). In other words,
\[
\mathcal{H}V = \{ W \in \text{Grass}(l, d) : W \cap V \neq \{0\} \}.
\]
This may also be formulated using the Plücker embedding (20): if \( v \) is any \( (d - l) \)-vector representing \( V \), then \( \mathcal{H}V \) consists of the subspaces \( W \in \text{Grass}(l, d) \) represented by \( l \)-vectors \( w \) such \( v \wedge w = 0 \).

Let \( \text{sec}(K, \text{Grass}(l, d)) \) denote the space of measurable maps \( V \) from some full \( \mu^c \)-measure subset of \( K \) to the Grassmannian manifold of all \( l \)-dimensional subspaces of \( \mathbb{R}^d \). Define the dual graph of each \( V \in \text{sec}(K, \text{Grass}(d - l, d)) \) to be
\[
\text{graph}(\mathcal{H}V) = \{ (t, v) \in K \times \text{Grass}(l, d) : v \in \mathcal{H}V(t) \}.
\]

Let \( \hat{m} \) be any \( u \)-state on \( \hat{M} \times \text{Grass}(l, d) \), \( m \) be its projection to \( M \times \text{Grass}(l, d) \) and \( \{ m_x : x \in \Sigma \} \) be the Rokhlin disintegration of \( m \) along the fibers \( K \times \text{Grass}(l, d) \) (recall section 5.2). The purpose of this section is to prove the following fact:

**Proposition 6.1.** We have \( m_x(\text{graph}(\mathcal{H}V)) = 0 \) every \( V \in \text{sec}(K, \text{Grass}(d - l, d)), u \)-state \( \hat{m} \) and every \( x \in \Sigma \).

The following terminology will be useful. For each \( \hat{x} \in \hat{\Sigma} \), consider the following push-forward maps \( \text{sec}(K, \text{Grass}(l, d)) \rightarrow \text{sec}(K, \text{Grass}(l, d)) \):

(a) \( V \mapsto F_{\hat{x}}V \) given by
\[
F_{\hat{x}}V(t) = \hat{A}(\hat{x}, s)V(s) \text{ with } s = (f_{\hat{x}})^{-1}(t);
\]
(b) \( V \mapsto H_{\hat{x}, \hat{y}}^V \) given, for \( \hat{y} \in W_{\text{loc}}^V(\hat{x}) \), by
\[
H_{\hat{x}, \hat{y}}^V(t) = H_{(\hat{x}, s), (\hat{y}, t)}^V V(s) \text{ with } s = h_{\hat{x}, \hat{y}}^V(t);
\]
(c) \( V \mapsto H_{\hat{x}, \hat{y}}^{\text{loc}} \) given, for \( \hat{y} \in W_{\text{loc}}^{\text{loc}}(\hat{x}) \), by
\[
H_{\hat{x}, \hat{y}}^{\text{loc}}(t) = H_{(\hat{x}, s), (\hat{y}, t)}^{\text{loc}} V(s) \text{ with } s = h_{\hat{x}, \hat{y}}^{\text{loc}}(t).
\]
These are well defined because the \( h^t \) and \( h^a \) holonomy maps are absolutely continuous with respect to \( \mu^c \), as a consequence of (7).

### 6.1. Graphs have measure zero

Starting the proof of proposition 6.1, recall that each \( m_x \) is a probability measure on \( \hat{K} \times \text{Grass}(l, d) \), and
\[
m_x = \int m_{x,t} \varrho(x, t) \, d\mu^c(t),
\]
where each \( m_{x,t} \) is a probability measure on \( \{(x, t) \} \times \text{Grass}(l, d) \). Recall also that \( x \mapsto m_x \) is continuous, by corollary 5.9.
Let \( x \in \Sigma \) be fixed for the time being, and consider the functions
\[
G : K \times \text{Grass}(d - l, d) \to \mathbb{R}, \quad G(t, V) = m_{x,t}(\mathcal{S}V) \text{ and } g : K \to \mathbb{R}, \quad g(t) = \sup\{m_{x,t}(\mathcal{S}Z) : Z \in \text{Grass}(d - l, d)\}.
\]

**Lemma 6.2.** \( G : K \times \text{Grass}(d - l, d) \to \mathbb{R} \) and \( g(t) : K \to \mathbb{R} \) are measurable functions.

**Proof.** Let \( \mathcal{P}_1 \prec \mathcal{P}_2 \prec \cdots \) be an increasing sequence of finite partitions of \( \text{Grass}(d - l, d) \) such that \( \mathcal{P} = \bigvee_{i \in \mathbb{N}} \mathcal{P}_i \) is the partition into points (that such a sequence exists is clear, e.g. because the Grassmannian is compact). Write \( \mathcal{P}_i = \{P_1, \ldots, P_n_i\} \) and then define
\[
G_n : K \times \text{Grass}(d - l, d) \to \mathbb{R}, \quad G_n(t, V) = \sum_{j=1}^{n_i} m_{x,t}(\mathcal{S}P_j) \chi_{P_j}(V)
\]
where \( \chi_{P_j} : \text{Grass}(d - l, d) \to \mathbb{R} \) denotes the characteristic function of a measurable set \( B \subset \text{Grass}(d - l, d) \). By the Rokhlin disintegration theorem, each \( t \mapsto m_{x,t}(\mathcal{S}P_j) \) is a measurable function. It follows that \( G_n \) is measurable for every \( n \). Moreover, \( (G_n)_n \) converges to \( G \) at every point. Thus, \( G \) is measurable. Analogously,
\[
g_n : K \to \mathbb{R}, \quad g_n(t, V) = \max\{m_{x,t}(\mathcal{S}P_j) : j = 1, \ldots, n\}
\]
is measurable for every \( n \), and \( (g_n)_n \) converges pointwise to \( g \). Thus the map \( g \) is measurable. \(\Box\)

For each fixed \( t \in K \), the function \( V \mapsto G(t, V) \) is upper semicontinuous: if \( (V_n)_n \) converges to \( V \) then \( \mathcal{S}V_n \) is contained in a small neighborhood of \( \mathcal{S}V \), for every large \( n \), and then \( m_{x,t}(\mathcal{S}V_n) \) can not be much larger than \( m_{x,t}(\mathcal{S}V) \). Since \( \text{Grass}(d - l, d) \) is compact, it follows that the set
\[
\Gamma(t) = \{V \in \text{Grass}(d - l, d) : G(t, V) = g(t)\}
\]
is compact and non-empty (the supremum in the definition of \( g \) is attained) for every \( t \in K \).

**Theorem 6.3** (Theorem III.30 in [9]). \( \) Let \( (X, \mathcal{B}, \mu) \) be a complete probability space and \( Y \) be a separable complete metric space. Denote by \( \mathcal{B}(Y) \) the Borel \( \sigma \)-algebra of \( Y \). Let \( \kappa(Y) \) be the space of compact subsets of \( Y \), with the Hausdorff topology. The following are equivalent:

1. a map \( x \to K_x \) from \( X \) to \( \kappa(Y) \) is measurable;
2. its graph \( \{(x,y) \in X \times Y : y \in K_x\} \) is in \( \mathcal{B} \otimes \mathcal{B}(Y) \);
3. \( \{x \in X : K_x \cap U \neq \emptyset\} \in \mathcal{B} \) for any open set \( U \subset Y \).

Moreover, any of these conditions implies that there exists a measurable map \( \sigma : X \to Y \) such that \( \sigma(x) \in K_x \) for every \( x \in X \).

**Lemma 6.4.** A given \( V \in \text{sec}(K, \text{Grass}(d - l, d)) \) realizes the supremum of
\[
\{m_{x,\text{graph}\mathcal{S}V} : V \in \text{sec}(K, \text{Grass}(d - l, d))\}
\]
if and only if \( V(t) \in \Gamma(t) \) for \( \mu^\ast\)-almost every \( t \in K \). Moreover, there exists some \( \mathcal{V}_x \in \text{sec}(K, \text{Grass}(d - l, d)) \) that does realize this supremum.
**Proof.** By lemma 6.2, the set
\[ \{(t, V) : V \in \Gamma(t)\} = \{(t, V) : G(t, V) = g(t)\} \]
is a measurable subset of \( K \times \text{Grass}(d - l, d) \). Compare the second condition in theorem 6.3. Thus, from the last claim in the theorem, there exists some measurable map \( \mathcal{V}_x : K \rightarrow \text{Grass}(d - l, d) \) such that \( \mathcal{V}_x(t) \in \Gamma(t) \) for every \( t \in K \). In other words,
\[ m_{\mu}(\delta \mathcal{V}_x(t)) = G(t, \mathcal{V}_x(t)) = g(t) = \sup_Z m_{\mu}(\delta Z) \]
for every \( t \in K \). Given any \( \mathcal{V} \in \text{sec}(K, \text{Grass}(d - l, d)) \) we have
\[
m_{\mu}(\text{graph}\mathcal{V}) = \int m_{\mu}(\delta \mathcal{V}(t)) \varrho(x, t) \, d\mu^c(t) \\
\leq \int \sup_Z m_{\mu}(\delta Z) \varrho(x, t) \, d\mu^c(t) \\
= \int m_{\mu}(\delta \mathcal{V}_x(t)) \varrho(x, t) \, d\mu^c(t) = m_{\mu}(\text{graph}\mathcal{V}_x).
\] (21)
Thus, \( \mathcal{V}_x \) does realize the supremum. Moreover, (21) is an equality if and only if \( G(x, \mathcal{V}(t)) = g(t) \) for \( \mu^c \)-almost every \( t \in K \).

So far, we kept \( x \in \Sigma \) fixed. The next proposition shows that the supremum in lemma 6.4 is actually independent of \( x \). Denote
\[ \gamma = \sup \{ m_{\mu}(\text{graph}\mathcal{V}) : \mathcal{V} \in \text{sec}(K, \text{Grass}(d - l, d)), x \in \Sigma \}. \]

**Proposition 6.5.** \( \sup \{ m_{\mu}(\text{graph}\mathcal{V}) : \mathcal{V} \in \text{sec}(K, \text{Grass}(d - l, d)) \} = \gamma \) for every \( x \in \Sigma \).

**Proof.** Given any cylinder \([J] \subset \Sigma \), choose a positive constant \( c < \mu([J]) / \kappa \), where \( \kappa > 0 \) is the constant in (18). Consider any \( \tilde{x} \in \Sigma \) and \( \tilde{\mathcal{V}} \in \text{sec}(K, \text{Grass}(d - l, d)) \). For each \( k \geq 1 \) and \( y \in \sigma^{-k}(\tilde{x}) \), define \( \mathcal{V}_y^k \in \text{sec}(K, \text{Grass}(d - l, d)) \) by
\[ \mathcal{F}_\gamma^k \mathcal{V}_y^k = \tilde{\mathcal{V}}, \quad \text{that is,} \quad A^k(t, y) \mathcal{V}_y^k(t) = \tilde{\mathcal{V}}(f^k_y(t)) \quad \text{for each} \ t \in K. \]
By corollary 5.10, \( m_{\mu}(\text{graph}\tilde{\mathcal{V}}) = \int m_{\mu}(\text{graph}\mathcal{V}_y^k) \, d\mu^c_{\kappa,\gamma}(y) \) and so
\[
m_{\mu}(\text{graph}\tilde{\mathcal{V}}) \leq \mu^c_{\kappa,\gamma}([J]) \sup \{ m_{\mu}(\text{graph}\mathcal{V}) : \mathcal{V} \in \text{sec}(K, \text{Grass}(d - l, d)), x \in [J] \} + (1 - \mu^c_{\kappa,\gamma}([J])) \gamma. \]
By (18), there exist arbitrary large values of \( k \) such that \( \mu^c_{\kappa,\gamma}([J]) \geq c \). Thus
\[
m_{\mu}(\text{graph}\tilde{\mathcal{V}}) \leq c \sup \{ m_{\mu}(\text{graph}\mathcal{V}) : \mathcal{V} \in \text{sec}(K, \text{Grass}(d - l, d)), x \in [J] \} + (1 - c) \gamma. \]
Varying \( \tilde{x} \in \Sigma \) and \( \tilde{\mathcal{V}} \in \text{sec}(K, \text{Grass}(d - l, d)) \), we can make the left-hand side arbitrarily close to \( \gamma \). It follows that
\[
\sup \{ m_{\mu}(\text{graph}\mathcal{V}) : \mathcal{V} \in \text{sec}(K, \text{Grass}(d - l, d)), x \in [J] \} \geq \gamma.
\]
The converse inequality is obvious. Thus, we have shown that the supremum over any cylinder \([J] \) coincides with \( \gamma \).
So, given any \( x \in \Sigma \) we may find a sequence \( (x_n)_n \rightarrow x \) such that the sequence

\[
(m_{x_n}(\text{graph} \, \gamma_{x_n}))_n
\]

converges to \( \gamma \), where (see lemma 6.4) each \( \gamma_{x_n} \) realizes the supremum at \( x_n \). Moreover, by proposition 5.5, up to restricting to a subsequence we may assume that \( (m_{x_n})_n \rightarrow m_x \) for every \( t \) in some full \( \mu' \)-measure set

\[
X \subseteq K.
\]

Then

\[
\gamma = \lim_n m_{x_n}(\text{graph} \, \gamma_{x_n})
= \lim_n \int m_{x_n}((\gamma)_{x_n}(t)) \varrho(x_n,t) \, d\mu'(t)
\leq \lim \sup_k m_{x_k}((\gamma)_{x_k}(t)) \varrho(x_n,t) \, d\mu'(t).
\]

For each fixed \( t \in X \), consider a sub-sequence \( (x_m)_m \) along which the lim sup is realized. It is no restriction to suppose that \( (\gamma_{x_m}(t))_m \) converges to some \( V \in \text{Grass}(d - l, d) \) as \( k \rightarrow \infty \). For any \( \epsilon > 0 \), let \( V_\epsilon \) be the closed \( \epsilon \)-neighborhood of \( V \). The fact that \( \gamma_{x_k}(t) \subseteq V_\epsilon \) for every large \( k \) implies that

\[
\lim \sup_k m_{x_k}((\gamma)_{x_k}(t)) \varrho(x_n,t) \, d\mu'(t).
\]

(because \( V_\epsilon \) is closed). Thus, making \( \epsilon \rightarrow 0 \) on the right-hand side,

\[
\lim \sup_k m_{x_k}((\gamma)_{x_k}(t)) \leq m_{x_k}(\gamma) \varrho \leq m_{x_k}(\gamma) V_\epsilon.
\]

Replacing this in (22), we find that \( \gamma \leq \int m_{x_k}((\gamma)_{x_k}(t)) \varrho(x,t) \, d\mu'(t) \) as claimed. \( \square \)

Having proved proposition 6.5, the proofs of the following two lemmas are analogous to those of lemmas 5.2 and 5.3 in [7], and so we omit them.

**Lemma 6.6.** Given any \( x \in \Sigma \) and \( \gamma \in \text{sec}(K, \text{Grass}(d - l, d)) \), we have that

\[
m_x(\text{graph} \, \gamma) = \gamma \text{ if and only if } m_x(\text{graph} \, \gamma) = \gamma \text{ for every } y \in \sigma^{-1}(x).
\]

As introduced in section 3.2, let \( \hat{\nu}_x : x \in \Sigma \) be the disintegration of \( \hat{\nu} = (\pi_1)_* \tilde{\mu} \) with respect to the partition into stable sets \( \{\Sigma^+ \times \{x\} : x \in \Sigma\} \). Observe that every \( \hat{\nu}_x \) is equivalent to \( \mu' \).

**Lemma 6.7.** For any \( x \in \Sigma \) and any \( \gamma \in \text{sec}(K, \text{Grass}(d - l, d)) \) we have that

\[
m_x(\text{graph} \, \gamma) \leq \gamma \text{ for } \hat{\nu}_{x_k} \text{ almost every } x_k \in W^\text{loc}_{x_k}(x).
\]

Hence, \( m_x(\text{graph} \, \gamma) = \gamma \) if and only if \( m_x(\text{graph} \, \gamma) = \gamma \) for \( \hat{\nu}_x \)-almost every \( x_k \in W^\text{loc}_{x_k}(x) \).

### 6.2. Sections over a periodic point

Let \( \hat{p} \) be a fixed (or periodic) point of \( \hat{\sigma} \) and \( \hat{z} \) be a homoclinic point associated to \( \hat{p} \). More precisely, we fix \( \hat{z} \) and \( i \geq 1 \) such that \( \hat{z} \in W^u_{\text{loc}}(\hat{p}) \) and \( \hat{\sigma}^i(\hat{z}) \in W^s_{\text{loc}}(\hat{p}) \). Denote \( \hat{p} = P(\hat{p}) \) and \( z = P(\hat{z}) \).

By the pinching hypothesis in section 2.6, the Oseledets decomposition of \( F \) restricted to \( K \)

has the form \( E^i(t) \) for \( i = 1 \) and \( E^i(t) \) for \( i > 1 \) and \( \hat{e}_i(t) \) for \( i \geq 1 \). Fix a measurable family \( e_i(t) \) of bases of \( \mathbb{R}^d \) with \( e_i(t) \in E^i(t) \) for every \( i \). The matrices of the iterates \( \hat{A}(p,t) \) relative to these bases are diagonal:
We are going to use the associated linear bases of $\Lambda^{(d-1)}(\mathbb{K}^d)$ and $\Lambda'(\mathbb{K}^d)$, defined at $\mu^c$-almost every $t \in K$ by

$$
\{e^I(t) = e^{i_1}(t) \wedge \cdots \wedge e^{i_d}(t), \text{ for } I = \{i_1 < \cdots < i_d-1\}\}
$$

and

$$
\{e^J(t) = e^{j_1}(t) \wedge \cdots \wedge e^{j_l}(t), \text{ for } J = \{j_1 < \cdots < j_l\}\}
$$

respectively.

By lemma 6.4 and proposition 6.5, we may choose $V^0 \in \sec(K, \text{Grass}(d-l,d))$ such that $m_p(\text{graph}V^0) = \gamma$. Define $V^j = F^{-j}_p V^0$ for $j \geq 1$. By proposition 6.6, we also have $m_p(\text{graph}V^j) = \gamma$ for every $j \geq 1$. Let $V^0 : K \to \Lambda^{(d-1)}(\mathbb{K}^d)$ be a representative of $V^0 \in \sec(K, \text{Grass}(d-l,d))$, in the sense that $V^0(t)$ is the projective class of $V^0(t)$ for each $t$. Then denote $V^j = F^{-j}_p V^0$ for each $j \geq 1$. Expressing $V^0$ in terms of the linear bases (23) of $\Lambda^{(d-1)}(\mathbb{K}^d)$,

$$
V^0(t) = \sum_{I \in t} v_I(t) e^I(t),
$$

we find that

$$
V^j(t) = \sum_{I \in t} \frac{v_I(f^j_p(t))}{a^{i_1}(t)} e^I(t),
$$

with $a^{i_1} = a^{i_1}(t) \cdots a^{i_{d-1}}(t)$. Note that $\lim_j (1/j) \log |a^{i_j}| = \lambda_i$, and so

$$
\lim_{j} \frac{1}{j} \log |a^{i_j}| = \lambda_i + \cdots + \lambda_{i_{d-1}}.
$$

(25)

Order the multi-indices

$$
I = \{i_1 < \cdots < i_{d-1}\}
$$

in such a way that the sums $\lambda_i + \cdots + \lambda_{i_{d-1}}$ are in increasing order (by the pinching condition these sums are all distinct).

Let $\tilde{I} = \{i_1, \cdots, i_{d-1}\}$ be the first multi-index, in this ordering, for which $v_I$ is not essentially zero. In what follows we assume that $f_{\tilde{I}}$ is ergodic for $\mu^c$. Then $\tilde{I}$ is the same for every $t \in K$ in a full $\mu^c$-measure set. The non-ergodic case can be reduced to this one by ergodic decomposition.

**Lemma 6.8.** The section $t \mapsto E^j(t)$ satisfies $m_p(\text{graph}E^j) = \gamma$.

**Proof.** By the Birkhoff ergodic theorem,

$$
\lim_n \frac{1}{n} \sum_{j=0}^{n-1} |v_I(f^j_p(t))| = \int |v_I| d\mu^c_p,
$$

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for \( \mu_{p}^{*} \)-almost every \( t \in K \). So, there exist some \( \delta > 0 \) such that
\[
\lim_{n \to \infty} \frac{1}{n} \sum_{j=0}^{n-1} |\nu_{j}(f_{j}^{h}(t))| > \delta > 0
\]
for every \( t \) in some full \( \mu_{p}^{*} \)-measure set. For any \( t \) in that set we may consider a sub-sequence \( (\nu_{j})_{k} \) such that \( |\nu_{j}(f_{j}^{h})| > \delta > 0 \). Then
\[
\lim_{k} \frac{1}{|V_{j}(t)|} V_{j}(h(t)) = \bar{e}(t),
\]
and so
\[
\lim_{k} V_{j}(h(t)) = E_{j}^{0}(t) + E_{j}^{1}(t) + \cdots + E_{j}^{\hat{r}-1}(t) = E_{j}^{\hat{r}}(t)
\]
for \( \mu' \)-almost every \( t \). We also have that \( m_{p,j}(V_{j}(t)) = \sup_{p} m_{p,j}(V) \), and then lemma 6.4 implies that \( m_{p,j}(E_{j}(t)) = \sup_{p} m_{p,j}(V) \) for \( \mu' \)-almost every \( t \), as we claimed. 

This means that, from the start, we may take \( V_{0}(t) \) to coincide with one of the invariant sections \( E_{j}^{0} \) given by the Oseledets decomposition, for \( \mu_{p}^{*} \)-almost every \( t \in K \). Define \( V' = F_{\hat{z}}^{-1} V^{0} \). We have that \( m_{c}(\text{graph}\nu_{j} V') = \gamma \) and, by lemma 6.7, \( \hat{m}_{(\hat{z},c)}(\text{graph}\nu_{j} V') = \gamma \) for \( \mu' \)-almost all \((\hat{z},c) \in W_{\text{loc}}^{c}(\hat{z})\). For each \((x',p) \in W_{\text{loc}}^{c}(\hat{p})\), define \( \nu_{j}(x',p) = H_{(\hat{z},c),(x',p)}^{c}(V') \), where \((x',c)\) is the unique point in \( W_{\text{loc}}^{c}((x',p)) \cap W_{\text{loc}}^{c}(\hat{z}) \). Since \( \hat{m} \) is a \( u \)-state, and \( h_{(\hat{z},p),\hat{c}}^{c} \delta_{c}^{\hat{z}} = \mu_{p}^{*} \), this implies that
\[
\hat{m}_{(\hat{z},c)}(\nu_{j}(x',p)) = \gamma \quad \text{for} \quad \mu' \text{-almost every} \quad (x',p) \in W_{\text{loc}}^{c}(p).
\]

Denote \( V_{j}^{c}(x',p) = F_{(\hat{z},c)}^{-j} V_{j}(x',p) \) for each \((x',p)\) and \( j \geq 1 \). In particular, \( V_{j}^{c}(p) = F_{(\hat{z},c)}^{-j} V_{j}(p) \). We are going to prove that for a large set of \( j \) the \( V_{j}^{c}(p) \) have no intersection.

**Proposition 6.9.** There exists \( N \geq 1 \) such that for every \( M \in \mathbb{N} \) and \( \delta > 0 \) there exist \( m_{1} < m_{2} < \cdots < m_{M} \) and \( K \subset K \), with \( \mu(K) > 1 - \delta \) and \( V_{j}^{m_{j}}(t) \cap \cdots \cap V_{j}^{M_{j}}(t) = \emptyset \) for any choice of \( m_{1} < m_{2} < \cdots < m_{M} \) and \( t \in K \).

**Proof.** Let \( V_{j} : K \to \Lambda^{d-j}([\mathbb{K}]^{d}) \) be such that \( V_{j}(t) \) is a unitary \( d-j \) vector that represents \( V_{j}(t) \) in \( \Lambda^{d-j}([\mathbb{K}]^{d}) \). We can write it as
\[
V_{j}(t) = \sum_{l} v_{l}(t) e_{l}^{j}(t).
\]

Then
\[
F_{(\hat{z},c)}^{-j} V_{j}(t) = \Lambda^{1}(\hat{p},t)^{-1} V_{j}(f_{j}^{h}(t)) = \sum_{l} \frac{v_{l}(f_{j}^{h}(t))}{d_{j}^{h}(t)} e_{j}^{h}(t).
\]

Let \( N = \dim \Lambda^{1}([\mathbb{K}]^{d}) \). Given any \( m_{1} < m_{2} < \cdots < m_{N} \) and \( t \in K \) such that \( V_{j}^{m_{j}}(t) \cap \cdots \cap V_{j}^{m_{j}}(t) \neq \emptyset \), there is some non-zero \( W(t) \in \Lambda^{1}([\mathbb{K}]^{d}) \) such that
\[
W(t) \cap F_{(\hat{z},c)}^{-m_{j}} V_{j}(t) = 0.
\]

(27)
Write
\[
W(t) = \sum_I \omega_I(t)e_I(t) \quad \text{where } J = \{1, 2, \ldots , l\} \setminus I.
\]

Then (27) can be written as
\[
\sum_I \frac{\nu_I(f^m(t))}{d_I^m(t)} \omega_I(t)z_I = 0,
\]
for every \(1 \leq k \leq N\), where \(z_I = \pm 1\) is the sign of \(e_1^l(t) \land \cdots \land e_{l-k}^l(t) \land e_1^l(t) \land \cdots \land e_{l-k}^l(t)\).

This may be written as
\[
B(t)x = 0
\]
(28)
where
\[
B(t) = \begin{pmatrix}
\frac{\nu_1(f^m(t))}{d_1^m(t)} & \cdots & \frac{\nu_n(f^m(t))}{d_n^m(t)} \\
\vdots & \ddots & \vdots \\
\frac{\nu_1(f^m(t))}{d_1^m(t)} & \cdots & \frac{\nu_n(f^m(t))}{d_n^m(t)}
\end{pmatrix}
\]
(29)
and \(x = (z_{I_1}\omega_{I_1}, \ldots , z_{I_N}\omega_{I_N})^T\).

So, in order to prove that the intersection is necessarily empty, it suffices to show that (28) has no non-zero solutions, in other words, that \(\det B(t) \neq 0\). We are going to use the following fact:

**Lemma 6.10.** Let \(b_i^n : K \to \mathbb{K}\), for \(1 \leq i \leq d\) and \(n \in \mathbb{N}\), be measurable functions and suppose there exist \(\chi_1 < \chi_2 < \cdots < \chi_d\) such that
\[
\lim_{n \to \infty} \frac{1}{n} \log |b_i^n(t)| = \chi_i \text{ for } \mu^c\text{-almost every } t.
\]
(30)

Then for every \(M \in \mathbb{N}\) and \(\delta > 0\) there exist \(n_1 < n_2 < \cdots < n_M\) and \(K \subset K\) with \(\mu(K) > 1 - \delta\), such that for any choice of a set \(\{k_1, \cdots , k_d\} \subset \{1, \ldots , M\}\) with \(k_1 < \cdots < k_d\), the matrix
\[
B(t) \in \mathbb{K}^{d \times d}, \quad B_{ij}(t) = b_{ij}^n(t),
\]
has non-zero determinant for every \(t \in K\).

For the proof we need the following simple algebraic fact:

**Lemma 6.11.** Let \(C = (c_{ij})_{1 \leq i,j \leq d}\) be a square matrix with \(c_{i}^1 \neq 0\) for every \(i = 1, \ldots , d\).
Then
\[
\det C = \prod_{i=1}^d c_{i}^1 \cdot \det E
\]
where \(E = (e_{ij})_{2 \leq i,j \leq d}\) is defined by
\( e_j^i = \frac{c_j^i}{c_i^i} - \frac{c_j^i}{c_i^i} \). \hspace{1cm} (31)

**Proof.** The assumption ensures that we may write

\[
\det C = c_1^1 \ldots c_d^1 \begin{bmatrix}
1 & 1 & \cdots & 1 \\
\vdots & \ddots & \ddots & \vdots \\
\vdots & \ddots & \ddots & \vdots \\
c_1^d & c_2^d & \cdots & c_d^d
\end{bmatrix}.
\]

Subtracting the first column from each one of the others, we end up with

\[
\det C = c_1^1 \ldots c_d^1 \begin{bmatrix}
e_2^1 & \cdots & e_d^1 \\
\vdots & \ddots & \vdots \\
e_2^1 & \cdots & e_d^1
\end{bmatrix}
\]

as claimed. \(\Box\)

**Proof of lemma 6.10.** Let us write \( b_i^{1,n} = b_i^n \) for \( i = 1, \ldots, d \) and \( n \geq 1 \). The hypothesis \((30)\) implies that there exist \( n_1 \geq 1 \) and \( K_1 \subset K \) with \( \mu(K_1) > 1 - \delta/M \) such that

\[
b_i^{1,n}(t) \neq 0 \text{ for } n \geq n_1, t \in K_1 \text{ and } i = 1, \ldots, d.
\] \hspace{1cm} (32)

Let \( n_1 \) be fixed and define (compare \((31)\))

\[
b_i^{2,n_1}(t) = \frac{b_i^{1,n}(t)}{b_i^{1,n}(t)} - \frac{b_i^{1,n}(t)}{b_i^{1,n}(t)} \text{ for } i = 2, \ldots, d \text{ and } n > n_1.
\] \hspace{1cm} (33)

From \((30)\), and the observation that \( \chi_1 > \chi \), we get that

\[
\lim_{n \to \infty} \frac{1}{n} \log |b_i^{2,n_1}(t)| = \lim_{n \to \infty} \frac{1}{n} \log |b_i^{1,n}(t)| = \chi.
\] \hspace{1cm} (34)

In particular, there exists \( n_2 > n_1 \) and \( K_2 \subset K_1 \) with \( \mu(K_2) > 1 - 2\delta/M \) such that

\[
b_i^{2,n_1}(t) \neq 0 \text{ and } b_i^{1,n}(t) \neq 0 \text{ for } n \geq n_2, t \in K_2 \text{ and } i = 2, \ldots, d
\] \hspace{1cm} (35)

(the second condition follows immediately from \((32)\) and the fact that \( n_2 > n_1 \), but we mention it explicitly, for consistency with what follows).

Next, proceed by induction on \( l \leq M \): Suppose that we have defined an increasing sequence of numbers \( n_1 < \cdots < n_l \), a decreasing sequence of sets \( K_l \supset \cdots \supset K_1 \) with \( \mu(K_l) > 1 - l\delta/M \), and for every \( 1 \leq j \leq \max\{l, d\} \) a family of measurable functions \( b_i^{j,n_1-\cdots-n_1} : K_l \to \mathbb{R} \setminus \{0\} \) with \( 1 \leq k_1 < \cdots < k_j \leq l \) satisfying the following relation:

\[
b_i^{j,n_1-\cdots-n_1}(t) = \frac{b_i^{j-1,n_1-\cdots-n_1}(t)}{b_i^{j-1,n_1-\cdots-n_1}(t)} - \frac{b_i^{j-1,n_1-\cdots-n_1}(t)}{b_i^{j-1,n_1-\cdots-n_1}(t)},
\] \hspace{1cm} (36)

for \( i = j, \ldots, d \) and \( t \in K_j \).

Suppose \( l < M \). Fix, \( 1 < j \leq \max\{l + 1, d\} \) and \( 1 \leq k_1 < \cdots < k_j \leq l \), define the functions \( b_i^{j+1,n_1-\cdots-n_1} : K_l \to \mathbb{R} \) inductively by
for $i = j + 1, \ldots, d$, $t \in K_l$ and $n \geq n_l$. Then, arguing as in (34) and using induction on $j$,

$$\lim_{n \to \infty} \frac{1}{n} \log |b_i^{j+1,n_1,\ldots,n_l}(t)| = \chi_i.$$  

(38)

Hence we can find $n_{l+1} > n_l$ and $K_{l+1} \subset K_l$ with

$$\mu(K_{l+1}) > 1 - (l + 1)\delta/M$$

such that, for every $1 < j \leq \max\{l + 1, d\}$ and $1 \leq k_1 < \cdots < k_l \leq l$,

$$b_i^{j+1,n_1,\ldots,n_l}(t) \neq 0 \text{ for } i = j + 1, \ldots, d \text{ and } t \in K_{l+1}.$$  

Now fix $\{k_1, \ldots, k_d\} \subset \{1, \ldots, M\}$ with $k_1 < \cdots < k_d$, and define for $t \in K_M$ the matrix

$$B(t) \in \mathbb{R}^{d \times d}, \quad B_{ij}(t) = b_i^{k_j}(t).$$

Then, in view of the recursive relations (36) and (37), we may apply lemma 6.11 $d$-times to $C = B(t)$ to conclude that

$$\det B(t) = \prod_{i=1}^d b_i^{k_1}(t) \prod_{i=2}^d b_i^{k_2,k_1}(t) \cdots \prod_{i=d}^d b_i^{k_d,\ldots,k_1}(t).$$

This completes our argument. \qed

Let us go back to proving proposition 6.9. The twisting condition (section 2.6), implies that

$$\lim_{n \to \infty} \frac{1}{n} \log |\nu_{l}(f^n(t))| = 0 \text{ for } \mu^\circ\text{-almost every } t \in K$$

and $I = \{i_1 < \cdots < i_{d-1}\}$. Then, by (25),

$$\lim_{n \to \infty} \frac{1}{n} \log \frac{|\nu_{l}(f^n(t))|}{|d^{\circ,n}(t)|} = -(\lambda_{i_1} + \cdots + \lambda_{i_{d-1}}) \text{ for } \mu^\circ\text{-almost every } t \in K.$$

The pinching condition ensures that these sums are all distinct. Then we may apply lemma 6.10 to the functions

$$b_{i}^{k_l}(t) = \frac{\nu_{l}(f^n(t))}{d^{\circ,n}(t)}.$$  

We get that there exist $m_1 < \cdots < m_d$ and $\tilde{K} \subset K$ with $\mu(\tilde{K}) > 1 - \delta$ such that for every $\{k_1, \ldots, k_d\} \subset \{1, \ldots, M\}$ with $k_1 < \cdots < k_d$, the matrix $B(t)$ defined in (29) is invertible for every $t \in \tilde{K}_M$.

**Proof of proposition 6.1.** Assume for the sake of contradiction that $\gamma > 0$. Then let $2\delta < \gamma$ and take $C > 0$ large enough that $C(\gamma - 2\delta) > 1$. Consider the sequence of integers $I = \{n_1, n_2, \ldots, n_{CN}\}$ given by proposition 6.9. Then there exists $\tilde{K} \subset K$ with $\mu^\circ(\tilde{K}) > 1 - \delta$ such that

$$\nu_{m_1}(t) \cap \cdots \cap \nu_{m_d}(t) = \emptyset$$

(39)
for every $t \in \tilde{K}$ and every $\{k_1 < \cdots < k_\gamma\} \subset \{1, \ldots, CN\}$.

First, suppose that $\tilde{m}_p(\text{graph}\gamma)V^j_p = \gamma$. The property (39) means that the sets $\mathcal{V}_m(t)$ are $N$-wise disjoint for every $t \in \tilde{K}$. Then,

$$\tilde{m}_p \left( \bigcup_{j \in I} \text{graph}\gamma V^j_p \right) \geq \tilde{m}_p \left( \bigcup_{j \in I} \text{graph}\gamma V^j_p \mid \tilde{K} \right) \geq \frac{1}{N} \sum_{j} \tilde{m}_p \left( \text{graph}\gamma V^j_p \mid \tilde{K} \right) \geq C(\gamma - \delta) > 1.$$ 

This is a contradiction because the measure $\tilde{m}_p$ is a probability.

Now we treat the general case. By (26), $\tilde{m}_{(x',p)}(\text{graph}\gamma)V^j_{m'}(x',p) = \gamma$ for every $j$ and $\mu^l$-almost every $(x',p) \in W^*_{\text{loc}}(p)$. In particular, we may a sequence $(x'_k, p)_k \rightarrow p$ with that property. Moreover, let $B_k(t)$ be the matrix defined by a system of equations as in (29), with coefficients depending on $V^m_{(x',p)}$ instead of $V^m_p$. Keep in mind that, by definition,

$$V^j_{(x',p)} = \mathcal{H}^{j}_{(x',p)}(V').$$

The sequence $\mathcal{H}^{j}_m(\gamma)_{(x',p)}$ converges uniformly to $\mathcal{H}^{j}_{m'}(V')$ when $k \rightarrow \infty$. Let $\mathcal{V}_m = \mathcal{H}^{j}_{m'}(V')$. By lemma 5.6 (together with the observation that $L^1$ convergence implies convergence almost everywhere over some subsequence), up to restricting to some subsequence of values of $k$ we have

$$\lim_k \mathcal{V}^m_{(x',p)}(t) = \mathcal{V}^m_p(t) \text{ for } \mu^l\text{-almost every } t \in K.$$ 

This proves that $B_k$ converges almost everywhere to $B$.

Recall that $\det B(t) \neq 0$ for every $t \in K$, by lemma 6.10. Then, there exist $\tilde{L} \subset \tilde{K}$ with $\mu^L(\tilde{L}) > 1 - 2\delta$ and $k_0 \geq 1$ such that $\det B_k(t) \neq 0$ for every $t \in \tilde{K}$ and $k \geq k_0$. Then, applying the previous argument with $(x',p)$ and $B_k$ instead of $\tilde{p}$ and $B$, we get that

$$\tilde{m}_{(x',p)} \left( \bigcup_{j \in I} \text{graph}\gamma V^j_{(x',p)} \right) \geq \frac{1}{N} \sum_{j} \tilde{m}_{(x',p)} \left( \text{graph}\gamma V^j_{(x',p)} \mid \tilde{L} \right) \geq C(\gamma - 2\delta) > 1$$

for every $k \geq k_0$. Thus, again we get a contradiction (because $\tilde{m}_{(x',p)}$ is a probability). \qed

7. Convergence to Dirac measures

The goal of this section is to prove the following theorem:

**Theorem 7.1.** There exists a measurable map $\hat{\xi} : \hat{M} \rightarrow \text{Grass}(l, d)$ such that, given any $u$-state $\hat{m}$ on $M \times \text{Grass}(l, d)$, we have

$$\hat{m}_{x,t} = \delta_{\hat{\xi}(x,t)} \text{ for } \hat{\mu}\text{-almost every } (\hat{x}, t) \in \hat{M}.$$ 

In particular, there exists a unique $u$-state.

7.1. Quasi-projective maps

We begin by recalling the notion of quasi-projective map, which was introduced by Furstenberg [12] and extended by Gol’dsheid, Margulis [14]. See also [2, section 2.3] for a related discussion.
Let $v \mapsto [v]$ be the canonical projection from $\mathbb{K}^d$ minus the origin to the projective space $\mathbb{P}\mathbb{K}^d$. We call $P.# : \mathbb{P}\mathbb{K}^d \to \mathbb{P}\mathbb{K}^d$ a projective map if there is some $P \in GL(d, \mathbb{K})$ that induces $P.#$ through $P.#([v]) = [P(v)]$. The space of projective maps has a natural compactification, the space of quasi-projective maps, defined as follows. The quasi-projective map $Q.#$ induced in $\mathbb{P}\mathbb{K}^d$ by a non-zero, possibly non-invertible, linear map $Q : \mathbb{K}^d \to \mathbb{K}^d$ is given by

$$Q.#([v]) = [Q(v)].$$

Observe that $Q.#$ is well defined and continuous on the complement of the kernel $\ker Q.# = \{ [v] : v \in \ker Q \}$.

More generally, one calls $P.# : \text{Grass}(l, d) \to \text{Grass}(l, d)$ a projective map if there is $P \in GL(d, \mathbb{K})$ that induces $P.#$ through $P.#(\xi) = P(\xi)$. Furthermore, the quasi-projective map $Q.#$ induced in $\text{Grass}(l, d)$ by a non-zero, possibly non-invertible, linear map $Q : \mathbb{K}^d \to \mathbb{K}^d$ is given by

$$Q.#(\xi) = Q(\xi).$$

Observe that $Q.#$ is well defined and continuous on the complement of the kernel $\ker Q.# = \{ \xi \in \text{Grass}(l, d) : \xi \cap \ker Q \neq \{0\} \}$.

The space of quasi-projective maps inherits a topology from the space of non-zero linear maps, through the natural projection $Q \mapsto Q.#$. Clearly, every quasi-projective map $Q.#$ is induced by some linear map $Q$ such that $\|Q\| = 1$. It follows that the space of quasi-projective maps on any $\text{Grass}(l, d)$ is compact for this topology.

The following two lemmas are borrowed from section 2.3 of [2]:

**Lemma 7.2.** The kernel $\ker Q.#$ of any quasi-projective map is contained in some hyperplane of $\text{Grass}(l, d)$.

**Lemma 7.3.** If $(P_n)_n$ is a sequence of projective maps converging to some quasi-projective map $Q$ of $\text{Grass}(l, d)$, and $(\nu_n)_n$ is a sequence of probability measures in $\text{Grass}(l, d)$ converging weakly to some probability $\nu$ with $\nu(\ker Q) = 0$, then $(P_n)_n, \nu_n$ converges weakly to $Q, \nu$.

72. Convergence

Recall that, given $1 \leq l \leq d - 1$ and $1 \leq i_1 < \cdots < i_l \leq d$, we write

$$E^{i_1\cdots i_l}(t) = E^{i_1}(t) \wedge \cdots \wedge E^{i_l}(t) \in \Lambda^l(\mathbb{K}^d)$$

for every $t \in K$ such that the Oseledets subspaces $E^t_i$ are defined. By a slight abuse of language, we also denote by $E^{i_1\cdots i_l}(t)$ the associated vector subspace, that is,

$$E^{i_1}(t) \oplus \cdots \oplus E^{i_l}(t) \in \text{Grass}(l, d).$$

In this way, each $E^{i_1\cdots i_l}$ becomes an element of $\sec(K, \text{Grass}(l, d))$.

Let $p \in \bar{\Sigma}$ be the fixed point of $\sigma$ and $\hat{z} \in \Sigma$ be a homoclinic point of $p$ with $\hat{z} \in W^u_{\text{loc}}(p)$. Fix $t \in \mathbb{N}$ such that $\hat{\sigma}^t(\hat{z}) \in W^u_{\text{loc}}(p)$. For each $k \geq 0$, denote $\hat{z}_k = \hat{\sigma}^{-k}(\hat{z})$ and $z_k = P(\hat{z}_k)$. Observe that $f_{\hat{z}_k} = f_{z_k}$ and, similarly, $A(p, t) = A(p, t)$. We take advantage of this fact to simplify the notations a bit (removing ‘hats’) in the arguments that follow. Recall that, given $y \in W^u_{\text{loc}}(\hat{z})$, $t \in K$ and $V \in \sec(K, \text{Grass}(l, d))$, we take $H^{u}_{\hat{z}_k} V \in \sec(K, \text{Grass}(l, d))$ to be defined by

$$H^{u}_{\hat{z}_k} V(t) = H^{u}_{(i_1, \cdots, i_l)} V(s)$$

with $s = h^{u}_{\hat{z}_k}(t)$. 

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Proposition 7.4. Let $\eta = H^{\mu_0}_{\hat{p}_0, E^1, \ldots, E^d} \in \sec(K, \text{Grass}(l, d))$. For every sequence $(k_j)_j \to \infty$ there exists a sub-sequence $(k'_j)$, such that
\[
\lim_{i \to \infty} A^k(z_{k_i}, t_{k_i}) m_{z_{k_i}, t_{k_i}} = \delta_{\eta(t)}, \text{ where } t_k = (f^k_\eta)^{-1}(t),
\]
for $\mu$-almost every $t \in K$.

Proof. To simplify our notations, let $h = h^\eta_{\hat{p}_0}$ and $h_k = h^\eta_{\hat{p}_k}$. Observe that $f^k = h^\eta_{\hat{p}_k} \circ f^k_\eta \circ h_k$ and
\[
A^k(z_k, t_k) = H^\mu_{(p, h_k(t_k))}(A^k(p, h_k(t_k))) H^\mu_{(z_k, h_k(t_k))}.
\]
So $(A^k)_{z_k, t_k}$ is equal to
\[
(H^\mu_{(p, h_k(t_k))} A^k(p, h_k(t_k))) (H^\mu_{(z_k, h_k(t_k))}) m_{z_k, t_k}.
\]
Note that $H^\mu_{(z_k, h_k(t_k))}$ converges uniformly to the identity map $id$, because $\hat{z}_k$ converges to $\hat{p}$.

Let $K_0 \subset K$ be a full $\mu$-measure such that the conclusion of the Oseledets theorem holds at $(\hat{p}, t)$ for every $t \in K_0$. We claim that for any $t \in K_0$ and every sub-sequence of
\[
A^k(p, h_k(t_k))
\]
that converges, the limit is a quasi-projective transformation $Q_\#$ that maps every point outside $K$ to $E^1, \ldots, E^d(h_k(t_k)) \subset \text{Grass}(l, d)$. This can be seen as follows.

Given $w \in N(K^d)$ and $k \geqslant 1$, we may write
\[
w = \bigoplus_{1 \leqslant i_1 < \cdots < i_d \leqslant d} w^1_{i_1} \cdots w^d_{i_d} E^{i_1, \ldots, i_d}(f^k)^{-1}(h^\eta_{\hat{p}_k}(t))
\]
with coefficients $w^1, \ldots, w^d \in K$. It follows from the sub-exponential decay of angles of the Oseledets splitting that $|w^k|$ grows sub-exponentially in $k$ for every $i = 1, \ldots, N$. Recall that $(f^k)^{-1}(h(t)) = h_k(t_k)$. Then, the action of $A^k(p, h_k(t_k))$ in the projectivization of the exterior power is given by
\[
A^k(p, h_k(t_k)) w = \bigoplus_{j=1}^N w_j^k \left|A^k(p, (f^k_\eta)^{-1}(h(t))) \cdot E^j_{(f^k_\eta)^{-1}(h(t))}\right| A^k(p, (f^k_\eta)^{-1}(h(t)))| E^j_{(f^k_\eta)^{-1}(h(t))}.
\]
The quotient of the norms converges to zero for any $j > 1$. Thus, we have that either $A^k(p, h_k(t_k)) w \to E^j_{h(t)}$ or $A^k(p, h_k(t_k)) w \to 0$. The latter case means that $w$ is in the kernel of the limit. Thus, any limit quasi-projective transformation does map the complement of the kernel to $E^j_{h(t)}$, as claimed.

As an immediate consequence we get that for any $t \in K_0$ and every sub-sequence of
\[
H^\mu_{(p, h_k(t_k))} A^k(p, h_k(t_k))
\]
that converges, the limit is a quasi-projective transformation that maps every point outside the kernel to $H^\mu_{(p, h_k(t_k))} E^j_{h(t)}$. 

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By remark 2.1, the family \( \{ f_n^i : n, k \geq 1 \} \) is equicontinuous. Using Arzelà–Ascoli, it follows that we can find a sub-sequence of \( (k_j)_j \) along which the family \( (f_k^i)^{-1} \) converges to some \( g : K \to K \). Then, by proposition 5.5, there exists a further subsequence \( (k'_j)_j \) and a full \( \mu^c \)-measure set \( K_1 \subset K \) such that
\[
m_{z_{t'}, t'} \to m_{p, g(t)}
\]
for every \( t \in K_1 \).

By proposition 6.1 and lemma 5.7, there exists a full \( \mu^c \)-measure set \( K_2 \subset K \) such that \( m_{p, g(t)} \) gives zero weight to every hyperplane of \( \text{Grass}(l, d) \) for every \( t \in K_2 \). Then, by lemma 7.3 and the previous observations,
\[
\lim_{k \to \infty} A^k(z_k, t_k)_* m_{z_k, t_k} = \delta_{\eta(t)}
\]
along any sub-sequence such that \( A^k(z_k, t_k) \) converges. This yields the claim of the proposition. \( \square \)

**Remark 7.5.** The argument remains valid when one replaces the homoclinic point \( \hat{z} \) by any other point in \( W^u(\hat{p}) \).

It follows from proposition 4.3 that there is a full \( \mu^c \times \mu^d \)-measure subset of points \( \hat{x} \in \hat{\Sigma} \) such that
\[
\lim_{n \to \infty} A^n(x_n, t_n)_* m_{x_n, t_n} = m_{\hat{x}, \eta}
\]
for \( \mu^c \)-almost every \( t \in K \), \( x_n = P(\hat{\sigma}^{-n}(\hat{x})) \) and \( t_n = (f_n^i)^{-1}(t) \). Since the shift is ergodic with respect to the projection of \( \mu \) on \( \hat{\Sigma} \), one may also require that
\[
\lim_{j \to \infty} \hat{\sigma}^{-\eta_j}(\hat{x}) = \hat{z}
\]
for some sub-sequence \( (\eta_j)_j \to \infty \).

Fix any \( \hat{x} \in \hat{\Sigma} \) such that both conditions hold. Let \( k \geq 1 \) be fixed, for the time being. Then (40) implies that
\[
\lim_{j \to \infty} A^n(x_n, t_n)_* m_{x_n, t_n} = \lim_{j \to \infty} A^{n+k}(x_{n+k}, t_{n+k})_* m_{x_{n+k}, t_{n+k}}
\]
\[
= \lim_{j \to \infty} A^n(x_n, t_n)_* A^k(x_{n+k}, t_{n+k})_* m_{x_{n+k}, t_{n+k}}
\]
\[
= \lim_{j \to \infty} A^n(x_n, t_n)_* A^k(x_{n+k}, t_{n+k})_* m_{x_{n+k}, t_{n+k}}.
\]

Note also that, by definition,
\[
t_{n+k} = f_{n+k}^k(t_{n+k}).
\]

We use once more the fact that \( \{ f_n^i : n \in \mathbb{Z} \text{ and } \hat{x} \in \hat{\Sigma} \} \) is equicontinuous (remark 3.2). Using Ascoli–Arzelà, it follows that there exists a sequence \( (\eta_j)_j \to \infty \) such that \( (f_n^i)^{-1} \) converges to some \( g : K \to K \). Up to further restricting to a sub-sequence if necessary, proposition 5.8 ensures that
\[
m_{x_{n+k}, t_{n+k}} \text{ converges to } m_{z_k, g(t)} \text{ for } \mu^c \text{-almost every } t,
\]
where \( z_k = P(\hat{\sigma}^{-k}(\hat{z})) \) and \( g(t) = (f_n^i)^{-1}(g(t)) \).
Fix any \( t \in K \) such that the previous claims are fulfilled. Let \( (n'_i) \) be any sub-sequence of \( (n_j) \) such that \( A^n (x_{n_j}, t_{n_j}) \) converges to some quasi-projective map \( Q : \text{Grass}(l, d) \to \text{Grass}(l, d) \).

Then (41) may be written as

\[
Q \ast A^k (z_k, g(t)) \ast m_{\gamma, 0}(t) \ast.
\]

If \( \eta(g(t)) / \in \ker Q \) then, making \( k \to \infty \), we may use lemma 7.3 and proposition 7.4 to conclude that \( \hat{\sigma}^m (\hat{x}, t) = \delta_{Q \eta(g(t))} \). This gives the conclusion of theorem 7.1 under this assumption.

Let us show that we can always reduce the proof to this case. Recall that \( \hat{t} \in \mathbb{Z} \) was chosen so that \( \hat{\sigma}^n (\hat{x}) \in W_{\text{loc}} (\hat{p}) \). Define \( \hat{y}, \hat{w} \in \hat{\Sigma} \) by

\[
\hat{\sigma}^n (\hat{x}) \in W_{\text{loc}} (\hat{w}) \text{ and } \hat{\sigma}^m (\hat{x}) \in W_{\text{loc}} (\hat{w}) \text{ \cap W}_{\text{loc}} (z).
\]

See figure 1. Note that \( \hat{y} \) depends on \( k \) and \( j \) and \( \hat{w} \) that depends on \( k \). We denote \( y = P(\hat{y}) \) and \( w = P(\hat{w}) \). Moreover, \( y_n = P(\hat{\sigma}^n (\hat{y})) \) and \( w_n = P(\hat{\sigma}^n (\hat{w})) \) for each \( n \geq 0 \) Let \( m \in \mathbb{N} \) be fixed, for the time being. We have that \( x_i = y_i \) with \( 0 \leq i \leq n_j + k \). So,

\[
\hat{\sigma}^i (P(\hat{\sigma}^{n_j+k-i-m}(\hat{y}))) = y_{n_j+k} = x_{n_j+k}.
\]

Also \( \hat{\sigma}^{n_j+k}(y) \to \hat{\sigma}^i (\hat{w}) \), and so \( \hat{\sigma}^{n_j+k-i-m}(y) \to \hat{\sigma}^m (\hat{w}) \) when \( j \to \infty \). Therefore, by proposition 4.3

\[
\hat{m}_{j,i} = \lim_{j \to \infty} A^{n_j+k}(x_{n_j+k}, t_{n_j+k}) \ast m_{x_{n_j+k}, t_{n_j+k}} \ast
\]

\[
= \lim_{j \to \infty} A^m (y_n, t_m) \ast m_{y_n, t_m} \ast.
\]

where \( m_j = n_j + k + i + m \). The last expression may be rewritten as

\[
A^n (x_{n_j}, t_{n_j}) \ast A^{k+i} (y_{n_j+k+i}, t_{n_j+k+i}) \ast A^m (y_n, t_m) \ast m_{y_n, t_m} \ast.
\]

Making \( j \to \infty \),
\[(f_w^{-k+i+1})^{-1} \to (f_w^{-i+1})^{-1} \circ g\]

\[A^{k+i}(y_{n_i+k+i}, t_{n_i+k+i}) \to A^{k+i}(w, (f_w^{k+i})^{-1}g(t))\]

\[A^m(y_{n_i}, t_{n_i}) \to A^m(w, (f_w^{k+i+m})^{-1}g(t))\]

and, restricting to a sub-sequence if necessary,

\[m_{n_i} \to m \quad w_{n_i}(f_w^{k+i+m})^{-1}g(t)\]

for \(\mu^c\)-almost every \(t\).

**Lemma 7.6.** Denote \(\tilde{\eta}(s) = H_{\tilde{h}}^n(\tilde{h}(s), (\tilde{h}(s))_s^{-1}, (\tilde{h}(s))_s^{-1})\) with \(\tilde{h}(s) = h_{n,i}^p(s)\). Then there exists a full \(\mu^c\)-measure set \(\tilde{K} \subset K\) and a sub-sequence \((k_j)\) such that for every \(t \in \tilde{K}\) there exists a sub-sequence \((n_j') = n_j'(t)\), of \((n_j)(t)\) such that

\[A^{k_j}(x_{n_j}', t_{n_j}') \circ A^{k+i}(y_{n_i+k+i}, t_{n_i+k+i})\]

converges to some quasi-projective transformation \(\tilde{Q}\). Moreover, \(\tilde{\eta}((f_w^{k+i+1})^{-1}g(t))\) is not in \(\ker \tilde{Q}\) if \(j\) is sufficiently large, depending on \(t\).

**Proof.** As before denote \(h = h_{n,i}^p\) and \(h_k = h_{n,i}^{k,p}\). We begin by constructing the sub-sequence \((k_j)\). Note that \((f_w^{k+i+1})^{-1} = (f_w^{k+i})^{-1}(f_w^{k})^{-1}\) and \(w \to \tilde{z}\) when \(k \to \infty\). First, take a sub-sequence of values of \(k\) such that \((f_w^{k})^{-1}\) converges uniformly to some \(\phi\). Since \(h_k\) converges to the identity map, \((f_w^{k})^{-1}h = h_{n,i}^{k,p}\) also converges uniformly to \(\phi\). Note that \(\phi\) is absolutely continuous with respect to \(\mu^c\), by lemma 5.7. Recall that

\[\tilde{\eta}((f_w^{k+i+1})^{-1}g(t)) = H_{\tilde{h}}^n(\tilde{h}(s), (\tilde{h}(s))_s^{-1}, (\tilde{h}(s))_s^{-1}) E_{\tilde{h}}^{k+i+1} (f_w^{k+i+1})^{-1}\]

with \(\tilde{h}(s) = h_{n,i}^{p}(s)\). Up to restricting the sub-sequence of values of \(k\), we may use lemma 5.6 to get that

\[\tilde{\eta}((f_w^{k+i+1})^{-1}g(t)) \to \tilde{\eta}((f_w^{k+i})^{-1}g(t)) \text{ and} \]

\[E((f_w^{k+i})^{-1}h g(t)) \to E(g(t))\]

for every \(t\) in some full \(\mu^c\)-measure set \(K_t\). This defines the sub-sequence \((k_j)\) in the statement. In what follows, all the statements on \(k\) are meant restricted to this sub-sequence.

The twisting condition implies that

\[H_{\tilde{h}}^n F_{\tilde{h}}^{k+i+1} \cap (E^{k+i+1} + \cdots + E^{k}) = \{0\}\]

for any \(j_{i+1}, \ldots, j_d \in \{1, \ldots, d\}\) and a full \(\mu^c\)-measure set of values of \(t \in K\). In other words, \(A^t(\mu(h(t)))\) does not belong to any of the hyperplanes of Grass(l, d) determined by the Oseledets decomposition at the point \((p, t)\). Since \(\phi\) and \(g\) are absolutely continuous, there exists a full \(\mu^c\)-measure set \(K_2\) of values of \(t\) such that (43) holds with \(t\) replaced by \(\phi(g(t))\).

Take \(\tilde{K} = K_1 \cap K_2\). Fix any \(t \in \tilde{K}\) such that in addition \((\tilde{h}, h g(t))\) satisfies the conclusion of the Oseledets theorem. Consider any sub-sequence \((n_j')\), of \((n_j)\), such that \(A^{n_j'}(x_{n_j}', t_{n_j}')\) converges to some quasi-projective transformation \(\tilde{Q}\) when \(i \to \infty\). Then

\[A^{n_j'}(x_{n_j}', t_{n_j}') \circ A^{k+i}(y_{n_i+k+i}, t_{n_i+k+i})\]
converges to $\tilde{Q} = Q \circ A^{k+1}(w, (f_n^{k+1})^{-1}g(t))$ when $i \to \infty$. Moreover,
\[
\ker \tilde{Q} = A^{k+1}(w, (f_n^{k+1})^{-1}g(t))^{-1} \ker Q
= A^i(w, (f_n^{k+1})^{-1}g(t))^{-1} A^k(z_k, (f_n^k)^{-1}g(t))^{-1} \ker Q.
\]
Next, observe that
\[
A^k(z_k, (f_n^k)^{-1}g(t))^{-1} = \Theta_k A^{-k}(p, hg(t)) \Theta
\]
where $h = h^{n}_{p, z}$ and
\[
\Theta = H^{n}_{G(z(t)), (p, hg(t))} \text{ and } \Theta_k = H^{n}_{G(f_n^k p, (f_n^k)^{-1}hg(t))} (G(z(t)), (f_n^k)^{-1}g(t))^{-1}.
\]
By lemma 7.2, the kernel of $Q$ is contained in some hyperplane $\tilde{\mathcal{V}}$ of Grass($l$, $d$). Hence, $\Theta(\ker Q)$ is contained in the hyperplane $\Theta(\tilde{\mathcal{V}})$, of course. Since we take $t \in K$ to be such that the Oseledets theorem holds at $(\hat{p}, t)$, the backward iterates $A^{-k}(p, hg(t)) \Theta(\tilde{\mathcal{V}})$ are exponentially asymptotic to some hyperplane section $\tilde{\mathcal{V}}$ that is defined by a $(d - l)$-dimensional sum $E$ of Oseledets subspaces. This remains true for $\Theta_k A^{-k}(p, hg(t)) \Theta(\tilde{\mathcal{V}})$ because $\Theta_k$ converges exponentially fast to the identity map, since $\hat{z}_k$ converges to $\hat{p}$ exponentially fast. In other words, using (44),
\[
\text{dist}_{\text{Grass}((l, d))}(A^i(z_k, (f_n^k)^{-1}g(t))^{-1}\tilde{\mathcal{V}}, E((f_n^k)^{-1}hg(t))) \to 0
\]
exponentially fast as $k \to \infty$. Then, by (42), we have that $A^k(z_k, (f_n^k)^{-1}g(t))^{-1}\tilde{\mathcal{V}}$ converges to $E(\phi g(t))$. So,
\[
\ker \tilde{Q} \subset A^i(z, (f_n^i)^{-1}\phi g(t))^{-1}\tilde{\mathcal{V}} E(\phi g(t)).
\]
Keep in mind that $\hat{z} \in W^s_{\text{loc}}(\hat{p})$ and $\ell \in \mathbb{N}$ is such that $\sigma^\ell(\hat{z}) \in W^s_{\text{loc}}(\hat{p})$. Recall also (from section 4.1) that in the present setting all the local stable holonomies $h'$ and $H'$ are trivial. Define
\[
V^i(t) = H^{n}_{G(z(t)), (p, z)} H^{l}_{(p, z), \sigma(z(t))} E(t_2)
= H^{n}_{G(z(t)), (p, z)} \hat{A}^{-\ell}(\sigma(z), s) H^{l}_{(p, z), \sigma(z(t), s)} \hat{A}^i(p, t_2) E(t_2)
= H^{n}_{G(z(t)), (p, z)} \hat{A}^{-\ell}(\sigma(z), s) E(s)
\]
with $t_1 = h^{n}_{p, z}(t)$, $t_2 = h^{l}_{p, z}(t_1)$ and $s = f_n^l(t_2) = f_n^l(t_1)$. Then, by the twisting condition, $\cup_{i \in J} V^i$ cannot intersect any sum of the form $\cup_{i \in I} E^i$ with $\#I + \#J = d$. In particular, the distance between
\[
H^{n}_{G(z(t)), (p, z)} \hat{A}^{-\ell}(\sigma(z), s) E(s) \quad \text{and} \quad E^i(s)
\]
is positive. Equivalently, the distance between
\[
\hat{A}^{-\ell}(\sigma(z), s) E(s) \quad \text{and} \quad E^i(s)
\]
is positive. Then $\eta((f_n^e)^{-1}\phi g(t))$ does not intersect
\[
\hat{A}^{-\ell}(\sigma(z), \phi g(t)) E(\phi g(t)) = A^i(z, (f_n^i)^{-1}\phi g(t))^{-1} E(\phi g(t)),
\]
which implies that $\eta((f_n^e)^{-1}\phi g(t)) \notin \ker \tilde{Q}$. \qed
Having established lemma 7.6, we can now use the same argument as previously, to conclude that \( \hat{m}_{iL} = \delta_{\hat{\xi} \tilde{t}} \) at \( \mu' \)-almost every point also in this case. To do this, observe that for every \( m \) and \( k \) fixed there exist a sub-sequence \( (m'_i)_k \) of \( (m_j)_l \) such that
\[
m_{E_n} \rightarrow m_{E_n, (\tilde{\rho}'_n + a)^{-1} g(t)} \quad \text{for } \mu' \text{-almost every } t.
\]
Using a diagonal argument, we may choose \( (m'_i)_k \) to be independent of \( k \) and \( m \). Fix, once and for all, a full \( \mu' \)-measure subset \( K' \) such that (45) and the conclusions of lemma 7.6 and proposition 7.4 (more precisely, remark 7.5) hold for every \( t \in K' \).

For each fixed \( t \in K' \), fixing \( k \) sufficiently large and making \( m'_i \) go to infinity (along the sub-sequence given by lemma 7.6), we find that
\[
\hat{m}_{iL} = \hat{\tilde{Q}}_n (A^m(w_m,f_k^{-1})^{-1} g(t)) \cdot m_{\tilde{\rho}'_n, (f_k^{-1})^{-1} g(t)}.
\]
Then, making \( m \rightarrow \infty \) and using lemma 7.3 and proposition 7.4,
\[
\hat{m}_{iL} = \delta_{\xi(x,t)},
\]
where \( \hat{\tilde{Q}}_n = \hat{\delta}_n (f_k^{1-k} g(t)) \).

Thus we proved that \( \hat{m}_{iL} \) is a Dirac measure for \( \sim \)-almost every \( \tilde{x} \in \hat{\Sigma} \) and \( \hat{\rho}'_n \)-almost every \( t \in K \). Note also that the set \( \hat{M} \subset M \) of points \( \hat{x}, t \in \hat{M} \) such that \( \hat{m}_{iL} \) is a Dirac measure is measurable, since the map \( (\hat{x}, t) \rightarrow \hat{m}_{iL} \) is measurable and the set of Dirac measures is closed in the weak* topology is closed, then \( \hat{M} \) is measurable. Thus we have shown that \( \hat{M} \) has total \( \hat{\mu} \)-measure, which completes the proof of theorem 7.1.

8. Orthogonal complement

8.1. Eccentricity

Let \( L : \mathbb{K}^d \rightarrow \mathbb{K}^d \) be a linear isomorphism and \( 1 \leq l \leq d \). The \( l \)-dimensional eccentricity of \( L \) is defined by
\[
E(l, L) = \sup \left\{ \frac{m(L \mid \xi)}{\| L \| \| \xi \|} : \xi \in \text{Grass}(l, d) \right\}, \quad m(L \mid \xi) = \| (L \mid \xi)^{-1} \|^{-1}.
\]
We call any \( l \)-subspace \( \xi \in \text{Grass}(l, d) \) that realizes the supremum as most expanded \( l \)-subspace. These always exist, since the Grassmannian is compact and the expression depends continuously on \( \xi \).

These notions may be expressed in terms of the polar decomposition of \( L = K'DK \) with respect to any orthonormal basis: denoting by \( a_1, \ldots, a_d \) the eigenvalues of the diagonal operator \( D \), in non-increasing order, then \( E(l, L) = a_l/a_{l+1} \). The supremum is realized by any subspace \( \xi \) whose image under \( K \) is a sum of \( l \)-eigenspaces of \( D \) such that the product of the eigenvalues is \( a_1 \cdot \cdots \cdot a_l \). It follows that \( E(l, L) \geq 1 \), and the most expanded \( l \)-subspace is unique if and only if the eccentricity is strictly larger than 1.

**Proposition 8.1.** For every \( 0 < c < 1 \), there exists a set \( M_c \subset M \) with \( \mu(\hat{M}_c) > c \) such that \( E(l, A^n(f^{-n}(\hat{x}, t))) \rightarrow \infty \), and the image of the most expanded subspace by \( A^n(f^{-n}(\hat{x}, t)) \) converges to \( (\hat{x}, t) \), restricted to the iterates such that \( f^{-n}(\hat{x}, t) \in \hat{M}_c \).

For the proof, let us recall the following fact, whose proof can be found in [2]:

**Proposition 8.2.** Let \( \mathcal{N} \) be a weak* compact family of probabilities on \( \text{Grass}(l, d) \) such that all \( \nu \in \mathcal{N} \) give zero weight to every hyperplane. Let \( L_\alpha : \mathbb{K}^d \rightarrow \mathbb{K}^d \) be linear isomor-
phisms such that \((L_n)\nu_n\) converges to a Dirac measure \(\delta_\xi\) as \(n \to \infty\), for some sequence \(\nu_n\) in \(\mathcal{N}\). Then the eccentricity \(E_l(L_n)\) goes to infinity and the image \(L_n(\zeta_n)\) of the most expanding \(l\)-subspace of \(L_n\) converges to \(\xi\).

**Proof of proposition 8.1.** Given \(0 < c < 1\) take \(M_c \subset M\) to be a compact set, with \(\mu(M_c) > c\) and such that the restriction of the map \((x, t) \mapsto m(x, t)\) to \(M_c\) is continuous. This implies that

\[
\mathcal{N} = \{m(x, t) : (x, t) \in M_c\}
\]

is a weak* compact subset of the space of probability measures of \(\text{Grass}(l, d)\), and every measure in \(\mathcal{N}\) gives zero weight to every hyperplane. Moreover,

\[
A^*\left(\hat{f}^{-n}(\hat{x}, t)\right), m_{P \times \text{id}((\hat{x}, t))} = \delta_{\xi(\hat{x}, t)}.
\]

Take \(\hat{M}_c = (P \times \text{id})^{-1}(M_c)\). Then the claim follows from proposition 8.2, with \(L_n = A^*\left(\hat{f}^{-n}(\hat{x}, t)\right)\).

### 8.2. Adjoint cocycle

Fix any continuous Hermitian form \(\langle \cdot, \cdot \rangle_{(\hat{x}, t)}\) in \(\hat{M} \times \mathbb{K}^d\). Let \(\hat{F}^* : \hat{M} \times \mathbb{K}^d \to \hat{M} \times \mathbb{K}^d\) be the adjoint cocycle, defined over \(\hat{f}^{-1} : M \to M\) by

\[
\hat{F}^*(\hat{x}, t, v) = (\hat{f}^{-1}(\hat{x}, t), \hat{A}_*(\hat{x}, t)v)
\]

where \(\hat{A}_*(\hat{x}, t)\) is the adjoint \(\hat{A}(\hat{f}^{-1}(\hat{x}, t))^*\) of the matrix \(\hat{A}(\hat{f}^{-1}(\hat{x}, t))\) with respect to the Hermitian form. In other words, \(\hat{A}_*(\hat{x}, t)\) is characterized by

\[
\langle u, \hat{A}_*(\hat{x}, t)v \rangle_{\hat{f}^{-1}(\hat{x}, t)} = \langle \hat{A}(\hat{f}^{-1}(\hat{x}, t)u), v \rangle_{(\hat{x}, t)}
\]

for any \(u, v \in \mathbb{K}^d\) and \((\hat{x}, t) \in \hat{M}\). We have that

\[
W_{\hat{f}^{-1}}^u(\hat{x}, t) = W_{\hat{f}^{-1}}^u(\hat{x}, t) = W_{\hat{f}^{-1}}^u(\hat{x}, t).
\]

It is also easy to see that

\[
H_{\hat{f}}^{\hat{A}_*}(\hat{x}, t) = (H_{\hat{f}}^\hat{A}(\hat{x}, t))^* \quad \text{and} \quad H_{\hat{f}}^{\hat{A}_*}(\hat{x}, t) = (H_{\hat{f}}^\hat{A}(\hat{x}, t))^*,
\]

respectively, for any \((\hat{x}, t), (\hat{y}, s)\) in the same \(\hat{f}^{-1}\)-unstable set and any \((\hat{x}, t), (\hat{z}, r)\) in the same \(\hat{f}^{-1}\)-stable set.

The following fact is well known (see [21, proposition 2.7] for a similar result):

**Proposition 8.3.** The cocycles \(\hat{F}\) and \(\hat{F}^*\) have the same Lyapunov exponents. Moreover, if \(E_j^l, j = 1, \ldots, k\) are the Oseledets spaces of \(\hat{F}\) then the Oseledets spaces of \(\hat{F}^*\) are, respectively,

\[
E_j^l = \left[ E_1^1 \oplus \cdots \oplus E_{l-1}^1 \oplus E_{l+1}^1 \oplus \cdots \oplus E_k^1 \right]^{\perp}, j = 1, \ldots, k.
\]

**Proposition 8.4.** \(\hat{A}_*\) is simple, if and only if, \(\hat{A}_*\) is pinching.

**Proof.** Applying proposition 8.3 to the restriction of \(\hat{F}\) to the periodic leaf \(\{\hat{p}\} \times K\) we get that the cocycle \(A_*\) is pinching, if and only if, \(A\) is pinching. Moreover, the Oseledets
decomposition $E^*_1 \oplus \cdots \oplus E^*_k$ of $\hat{A}$ is given by the orthogonal complements of the Oseledets subspaces of $A$:

$$E^*_j = \left[ E^1 \oplus \cdots \oplus E^{j-1} \oplus E^{j+1} \oplus \cdots \oplus E^k \right] \perp.$$

We are going to use this for proving the twisting property, as follows.

Let $\phi_{\hat{p},\hat{z}} = H^{\hat{u}_{\hat{p},\hat{z}}} \circ H^{\hat{u}_{\hat{p},\hat{z}}} \circ \phi_{\hat{p},\hat{z}}$ and $\phi^*_p = H_{\hat{u},\hat{z}} \circ H_{\hat{u},\hat{z}}^*.$

Denote $h : K \to K, \ h(t) = h_{\hat{p},\hat{z}} \circ h_{\hat{p},\hat{z}}$ and $H_t = H^p_{(\hat{u},\hat{z})_{(i)}(\hat{p},\hat{z})} \circ H^p_{(\hat{u},\hat{z})_{(i)}(\hat{p},\hat{z})}.$

Then, for any $V \in \text{sec}(K, \text{Grass}(l, d)),$

$$\phi_{\hat{p},\hat{z}} V(t) = H_t \left(V(h^{-1}(t))\right) \quad \text{and} \quad \phi^*_p V(t) = H_{h(\hat{t})} V(h(t)).$$

First, we treat the case $l = 1.$ Define measurably for (almost) every $t \in K$ a linear base of unit vectors $e^j(t) \in E^j(t),$ $j = 1, \ldots, d.$ The twisting condition means that if

$$(\phi_{\hat{p},\hat{z}} e^j(t)) = \sum_{j=1}^d a_{k_j}(t) e^j(t),$$

then

$$\lim_{n \to \infty} \frac{1}{n} \log |a_{k_j}(f_{\hat{p},\hat{z}}^n(t))| = 0.$$

We need to deduce the corresponding fact for the adjoint. For this, write

$$(\phi^*_p e^j_*(t)) = \sum_{j=1}^d \beta_{k_j}(t) e^j_*(t).$$

Hence

$$\beta_{k_j}(t) (e^j_*(t), e^j(t)) = \langle \phi^*_p e^j_*(t), e^j(t) \rangle$$

$$= \langle e^j_*(h(t)), \phi_{\hat{p},\hat{z}} e^j(h(t)) \rangle$$

$$= a_{k_j}(h(t)) \langle e^j_*(h(t)), e^j(h(t)) \rangle,$$

by definition $\langle e^j_*(t), e^j(t) \rangle = \cos(\alpha^j(t))$ for every $1 \leq i \leq d,$ where

$$\alpha^j(t) = \angle \left( e^j_*(t), e^j(t) \right)$$

$$= \frac{\pi}{2} - \angle \left( e^j(t), E^1 \oplus \cdots \oplus E^j \oplus \cdots \oplus E^k \right),$$

so, by the Oseledets theorem,

$$\lim_{n \to \infty} \frac{1}{n} \log |\langle e^j_*(f_{\hat{p},\hat{z}}^n(t)), e^j(f_{\hat{p},\hat{z}}^n(t)) \rangle| = 0\hat{\mu}_{\hat{p}} \text{-almost everywhere.}$$
Then
\[
\lim_{n \to \infty} \frac{1}{n} \log |\beta_{kj}(f_p^n(t))| = \lim_{n \to \infty} \frac{1}{n} \log |\alpha_{kj}(f_p^n(h(t)))|
\]
also as \( h : K \to K \) preserves \( \mu^e_x \), this is true for \( \mu^e_x \)-almost everywhere.

For \( l > 1 \) the proof is just the same, using the inner product induced on \( \Lambda'(K^d) \) by \( \langle \cdot, \cdot \rangle \), that is,
\[
\langle v_1 \wedge \cdots \wedge v_1, w_1 \wedge \cdots \wedge w_l \rangle_{\Lambda'(K^d)} = \det \left( \langle v_i, w_j \rangle \right).
\]

Thus, we have shown that \( A \) is twisting if and only if \( A_e \) is twisting. \( \square \)

Applying proposition 8.1 to the adjoint cocycle we get:

**Corollary 8.5.** There exists a section \( \xi^* : \hat{M} \to \text{Grass}(I, d) \) which is invariant under the cocycle \( F_{A_e} \) and the unstable linear holonomies of \( A_e \).

Moreover, given any \( c > 0 \) there exists \( \hat{M}_c \subset \hat{M} \) with \( \hat{\mu}(\hat{M}_c) > c \) such that, restricted to the sub-sequence of iterates \( k \) such that \( f^k(p) \) in \( \hat{M}_c \), the eccentricity \( E(I, A_e^k(f^k(p))) = E(I, A_e^k(p)) \) goes to infinity and the image \( \Lambda^k_e(f^k(p)) \xi^*_{\mu} \) of the most expanded l-subspace tends to \( \xi^*(p) \) as \( k \to \infty \).

The next lemma relates the invariant sections of the two cocycles, \( F \) and \( F_{A_e} \):

**Lemma 8.6.** For \( \hat{\mu} \)-almost every \( \hat{x}, t \subset \hat{M} \), the subspace \( \xi(\hat{x}, t) \) is transverse to the orthogonal complement of \( \xi^*(\hat{x}, t) \).

**Proof.** Recall that the stable linear holonomies of \( \hat{A} \) are trivial. Thus, the same is true for the unstable linear holonomies of \( A_e \). So, the fact that \( \xi^* \) is invariant under unstable linear holonomies means that it is constant on local stable sets of \( f \). Then the same is true about his orthogonal complement \( \eta(\hat{x}, t) = \xi^*(\hat{x}, t)^\perp \), which means that it only depends on \( \eta(\hat{x}, t) = \eta(x, t) \), where \( x = P(\hat{x}) \). Recall that the graph of \( \eta(x, \cdot) \) over \( K \) has zero \( m_c \)-measure, by proposition 6.1:
\[
m_c(\text{graph}\eta_x) = \int \int \delta_{\xi^*}(\eta(x, t))d\mu^e(t)d\mu^e(x) = \mu^e \times \mu^e(\{x : \xi^*(x, t) \in \eta(x, t)\}) = 0
\]
for \( \nu \)-almost every \( x \in \Sigma \). Hence \( \hat{\mu}(\{\hat{x}, t : \xi(\hat{x}, t) \in \eta(x, t)\}) = 0 \), which proves the lemma. \( \square \)

9. **Proof of theorem A**

Denote by \( \eta(\hat{x}, t) \in \text{Grass}(d - I, d) \) the orthogonal complement of \( \xi^*(\hat{x}, t) \) at each \( (\hat{x}, t) \subset \hat{M} \).

Recall that \( \xi^* \) was defined in corollary 8.5 and is invariant under \( A_e \):
\[
\hat{A}_e(\hat{x}, t)\xi^*(\hat{x}, t) = \xi^*(\hat{f}^{-1}(\hat{x}, t)) \text{ for } \hat{\mu} \text{-almost every } (\hat{x}, t).
\]

Consequently, \( \eta \) is invariant under \( A_e \).

According to lemma 8.6, we have that \( K^d = \xi(\hat{x}, t) \oplus \eta(\hat{x}, t) \) at \( \hat{\mu} \)-almost every point. To prove theorem A we are going to show that the Lyapunov exponents of \( A \) along \( \xi \) are strictly greater than those along \( \eta \). For that, let
\[ \xi(\hat{x}, t) = \xi^i(\hat{x}, t) \oplus \cdots \oplus \xi^{u}(\hat{x}, t) \quad \text{and} \quad \eta(\hat{x}, t) = \eta^i(\hat{x}, t) \oplus \cdots \oplus \eta^{u}(\hat{x}, t) \]

be the Oseledets decomposition of \( A \) restricted to the two invariant sub-bundles, where \( \xi^u \) corresponds to the smallest Lyapunov exponent among \( \xi^i \) and \( \eta^i \) the largest among all \( \eta^i \).

Denote \( d_u = \dim \xi^u \) and \( d_i = \dim \eta^i \), and then let \( \lambda_u \) and \( \lambda_i \) be the Lyapunov exponents associated to these two sub-bundles, respectively. Define

\[ \Delta^u(\hat{x}, t) = \frac{\det (A^u(\hat{x}, t), \xi^u(\hat{x}, t))}{\det (A^u(\hat{x}, t), W(\hat{x}, t))}^{\frac{1}{\Delta}} \]

where \( W(\hat{x}, t) = \xi^u(\hat{x}, t) \oplus \eta^u(\hat{x}, t) \). By the Oseledets theorem

\[ \lim_{n \to \infty} \frac{1}{n} \log \Delta^u(\hat{x}, t) = \frac{d_u}{d_u + d_i} (\lambda_u - \lambda_i) . \]

The proof of the following proposition is identical to the proof of proposition 7.3 in [2]:

**Proposition 9.1.** For every \( 0 < c < 1 \) there exist a set \( \hat{M} \subset \hat{M} \) with \( \bar{\mu}(\hat{M}) > c \) such that for \( \bar{\mu} \)-almost every \( (\hat{x}, t) \in \hat{M} \)

\[ \lim_{n \to \infty} \Delta^u(\hat{x}, t) = \infty \]

restricted to the sub-sequence of values \( n \) for which \( \hat{f}^n(\hat{x}, t) \in \hat{M} \).

So now fix some \( 0 < c < 1 \) and \( \hat{M} \), given by proposition 9.1. Let \( g : \hat{M} \to \hat{M} \) be the first return map:

\[ g(\hat{x}, t) = \hat{f}^p(\hat{x}, t) . \]

Then we can define the induced cocycle \( G : \hat{M} \times \mathbb{R}^d \to \hat{M} \times \mathbb{R}^d \)

\[ G((\hat{x}, t), v) = (g(\hat{x}, t), Dg(\hat{x}, t)v) , \]

where \( Dg(\hat{x}, t) = A^u(\hat{x}, t) \). It is well known (see [23, proposition 4.18]) that the Lyapunov exponents of \( G \) with respect to \( \frac{1}{\bar{\mu}(\hat{M})} \bar{\mu} \) are the products of the exponents of \( \hat{F}_\lambda \) by the average return time \( 1/\bar{\mu}(\hat{M}) \). Thus, to show that \( \lambda_u > \lambda_i \) it suffices to prove the corresponding fact for \( G \).

Define

\[ \hat{\Delta}^k((\hat{x}, t)) = \frac{\det (D^k(\hat{x}, t), \xi^u(\hat{x}, t))^{\frac{1}{\Delta}}}{\det (D^k(\hat{x}, t), W(\hat{x}, t))}. \]

Then \( \hat{\Delta}^k((\hat{x}, t)) \) is a sub-sequence of \( \Delta^u((\hat{x}, t)) \) such that \( \hat{f}^n(\hat{x}, t) \in \hat{M} \). So, using proposition 9.1 we conclude that

\[ \lim_{n \to \infty} \sum_{j=0}^{k-1} \log \hat{\Delta} \left( g^j(\hat{x}, t) \right) = \lim_{n \to \infty} \log \hat{\Delta} \left( p \right) = \infty \]

for \( \bar{\mu} \)-almost every \((\hat{x}, t) \in \hat{M}\).

We need the following classical fact (see [17, corollary 6.10]):

**Lemma 9.2.** Let \( T : X \to X \) be a measurable transformation preserving a probability measure \( \nu \) in \( X \), and \( \varphi : X \to \mathbb{R} \) be a \( \nu \) integrable function such that \( \lim_{n \to \infty} \sum_{j=0}^{n-1} \left( \varphi \circ T^j \right) = +\infty \) at \( \nu \) almost every point. Then \( \int \varphi \, d\nu > 0 \).
Applying the lemma to $T = g$ and $\varphi = \log \hat{\Delta}$ we find that
\[
\lim_{k \to \infty} \frac{1}{k} \log \hat{\Delta}^k(\hat{x}, t) = \lim_{k \to \infty} \frac{1}{k} \sum_{j=0}^{k-1} \log \hat{\Delta} \left( g^j(\hat{x}, t) \right) = \int \log \frac{d\hat{\mu}}{\hat{\mu}(\hat{M}_t)} > 0
\]
at $\hat{\mu}$-almost every point. On the other hand the relation between Lyapunov exponents gives that
\[
\lim_{k \to \infty} \frac{1}{k} \log \hat{\Delta}^k(\hat{x}, t) = \frac{d_s}{d_u + d_s} \left( \lambda_u - \lambda_s \right) \frac{1}{\hat{\mu}(\hat{M}_t)}
\]
this means that $\lambda_u > \lambda_s$, so there is a gap between the first $l$ Lyapunov exponents and the remaining $d - l$ ones. Since this applies for every $1 \leq l \leq d$, we conclude that the Lyapunov spectrum is simple.

What is left to prove is that a simple cocycle is also a continuity point for the Lyapunov exponents.

**Proposition 9.3.** If $A$ is simple, then, for every $1 \leq i \leq d$, the functions $\lambda_i : H^\alpha(\hat{M}) \to \mathbb{R}$ are continuous in $A$.

**Proof.** Take $\Phi_k : \hat{M} \times \mathbb{R}^d \to \mathbb{R}$
\[
\Phi_k(\hat{x}, \nu) = \frac{\log \left\| A_k(\hat{x}) \nu \right\|}{\|\nu\|},
\]
then for every $k \in \mathbb{N}$, there exists an $F_{\hat{A}_k}$-invariant $u$-state $\hat{m}_k^u$ such that $\lambda_1(\hat{A}_k) = \int \Phi_k d\hat{m}_k^u$.

Passing to a subsequence if necessary we can suppose that $\hat{m}_k^u$ converges in the weak* topology to some $F_{\hat{A}}$-invariant $u$-state $\hat{m}_k^u$. By theorem 7.1 $\hat{m}_k^u = \int \delta_{\hat{m}_k^u} d\hat{\mu}(\hat{x})$, this implies that
\[
\lambda_1(\hat{A}_k) \to \int \Phi(\hat{x}, \nu) d\hat{m}_k^u = \lambda_1(\hat{A}).
\]

Now, using the same argument for every $i$-dimensional Grassmannian we get that $\lambda_1 + \cdots + \lambda_i$ is also continuous, concluding the proof.

The simplicity plus the continuity implies that there exists a neighbourhood of $A$ with simple Lyapunov spectrum. This completes the proof of theorem A.

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**Appendix. Continuous maps are dense in $L^1(M,N)$**

Let $M$ be a normal topological space and $N$ be a geodesically convex separable metric space (section 5.2). Denote by $\mathcal{F}$ the set of measurable maps $f : M \to N$. Given any regular $\sigma$-finite Borel measure $\mu$ on $M$, fix any point $\hat{0} \in N$ and define $L^1_{\mu}(M, N) =$
\( \{ f \in \mathcal{F} : \int \text{dist}_N \left( f(x), \hat{0} \right) \, d\mu(x) < \infty \} \). When \( \mu \) is a finite measure, the choice of \( \hat{0} \in N \) is irrelevant: different choices yield the same space \( L^1_\mu(M, N) \).

The function \( \text{dist}_{L^1_\mu(M, N)} : L^1_\mu(M, N) \times L^1_\mu(M, N) \to \mathbb{R} \) defined by

\[
\text{dist}_{L^1_\mu(M, N)}(f, g) = \int d_N \left( f(x), \hat{0} \right) \, d\mu(x)
\]

is a distance in \( L^1_\mu(M, N) \). The special case \( N = \mathbb{R} \) of the next proposition is well known, but here we need the following more general statement:

**Proposition A.1.** The subset of continuous maps \( f : M \to N \) is dense in the space \( L^1_\mu(M, N) \).

We call \( s : M \to N \) a simple map if there exist points \( v_1, \ldots, v_k \in N \) pairwise disjoint measurable sets \( A_1, \ldots, A_k \subset M \) with finite \( \mu \)-measure such that

\[
s(x) = \begin{cases} 
v_i & \text{if } x \in A_i \\
\hat{0} & \text{if } x \notin \bigcup_{i=1}^k A_i
\end{cases}
\]

Proposition A.1 is an immediate consequence of lemmas A.2 and A.3 below.

**Lemma A.2.** The set \( S \) of simple functions is dense in \( L^1_\mu(M, N) \).

**Proof.** Consider any \( f \in L^1_\mu(M, N) \). Given \( \epsilon > 0 \), fix a set \( K_0 \subset M \) with finite \( \mu \)-measure and such that

\[
\int_{M \setminus K_0} \text{dist}_N \left( f(x), \hat{0} \right) \, d\mu(x) \leq \frac{\epsilon}{4}.
\]

Let \( \{v_1, \ldots, v_i, \ldots \} \) be a countable dense subset of \( N \). The family

\[
\{ B(v_i, \frac{\epsilon}{2\mu(K_0)}) : i \in \mathbb{N} \}
\]

covers \( N \) and, consequently,

\[
B_i = B \left( v_i, \frac{\epsilon}{2\mu(K_0)} \right) \setminus \bigcup_{j < i} B \left( v_i, \frac{\epsilon}{2\mu(K_0)} \right), \quad i \in \mathbb{N}
\]

is a partition of \( N \). Then \( A_i = K_0 \cap f^{-1}(B_i), i \in \mathbb{N} \) is a partition of \( K_0 \) into measurable sets. Fix \( k \in \mathbb{N} \) large enough that

\[
\int_{K_0 \setminus \bigcup_{i=1}^k A_i} \text{dist}_N \left( f(x), \hat{0} \right) \, d\mu(x) \leq \frac{\epsilon}{4}.
\]

Now define \( s : M \to N \) by

\[
s(x) = \begin{cases} 
v_i & \text{if } x \in A_i \text{ for } i = 1, \ldots, k \\
\hat{0} & \text{if } x \notin \bigcup_{i=1}^k A_i
\end{cases}
\]
Then
\[
\int_{M\setminus \cup_{i=1}^{k} A_i} \text{dist}_N(f(x), s(x)) \, d\mu(x) = \int_{M\setminus \cup_{i=1}^{k} A_i} \text{dist}_N(f(x), \hat{0}) \, d\mu(x) \leq \frac{\epsilon}{2}
\]
and
\[
\int_{\cup_{i=1}^{k} A_i} \text{dist}_N(f(x), s(x)) \, d\mu(x) \leq \mu(\bigcup_{i=1}^{k} A_i) \frac{\epsilon}{\mu(K_0)} \leq \frac{\epsilon}{2}.
\]
Thus \(\text{dist}_{L_\tau(M,N)}(f, s) < \epsilon\), which proves the lemma. \(\Box\)

**Lemma A.3.** For every \(s \in \mathcal{S}\) and \(\epsilon > 0\) there exists a continuous map \(f : M \to N\) such that \(\text{dist}_{L_\tau(M,N)}(f, s) \leq \epsilon\).

**Proof.** Let \(A_i\) and \(v_i, i = 1, \ldots, k\) be as in the definition of the simple map \(s\) and \(\tau \geq 1\) be as in (19). Denote \(L = \max\{d(v_i, 0) : i = 1, \ldots, k\}\). For each \(i = 1, \ldots, k\), consider a compact set \(K_i \subset A_i\) such that \(\mu(A_i \setminus K_i) < \epsilon/(4k\tau L)\). Since the \(K_i\) are pairwise disjoint, and \(M\) is assumed to be normal, there exist pairwise disjoint open sets \(B_i \supset K_i, i = 1, \ldots, k\) with \(\mu(B_i \setminus K_i) < \epsilon/(4k\tau L)\). In particular, we also have \(\mu(A_i \setminus B_i) < \epsilon/(4k\tau L)\).

By the Urysohn lemma, there are continuous functions \(\psi_i : M \to \mathbb{R}, i = 1, \ldots, k\) such that
\[
\psi_i(x) = \begin{cases} 1 & \text{if } x \in K_i \\ 0 & \text{if } x \notin B_i. \end{cases}
\]

Now we use the assumption that \(N\) is geodesically convex. For each \(i = 1, \ldots, k\), fix \(\lambda_i : [0, 1] \to N\) with \(\lambda_i(1) = v_i\) and \(\lambda_i(0) = \hat{0}\). Then define \(f : M \to N\) by
\[
f(x) = \begin{cases} \lambda_i(\psi_i(x)) & \text{if } x \in B_i \text{ with } i = 1, \ldots, k \\ \hat{0} & \text{if } x \notin \bigcup_{i=1}^{k} B_i. \end{cases}
\]

It is clear that \(f\) is continuous, because the \(B_i\) are open and pairwise disjoint. Moreover, \(f(x) = s(x)\) if
\[
\text{either } x \in \bigcup_{i=1}^{k} K_i \text{ or } x \in M \setminus \left( \bigcup_{i=1}^{k} A_i \cup \bigcup_{i=1}^{k} B_i \right).
\]
All the other values of \(x\) fall into some of the following cases:

1. \(x \in A_i \cap (B_j \setminus K_j)\) for some \(i\) and \(j\): then
   \[
d_N(f(x), s(x)) \leq d_N(\lambda_i(\psi_j(x)), \hat{0}) + d_N(\hat{0}, v_i) \leq \tau d_N(v_j, \hat{0}) + d_N(\hat{0}, v_i) \leq 2\tau L.
   \]
2. \(x \in A_i \setminus \bigcup_{j=1}^{k} B_j\) for some \(i\): then \(d_N(f(x), s(x)) = d_N(\hat{0}, v_i) \leq L\).
3. \(x \in B_j \setminus \bigcup_{i=1}^{k} A_i\) for some \(j\): then
   \[
d_N(f(x), s(x)) = d_N(\lambda_j(\psi_j(x)), \hat{0}) \leq \tau d_N(v_j, \hat{0}) \leq \tau L.
   \]

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In either case, $x$ belongs to the set
\[
\bigcup_{i=1}^{k} (A_i \setminus B_i) \cup \bigcup_{j=1}^{k} (B_j \setminus K_j)
\]
which, by construction, has $\mu$-measure bounded by $\epsilon/(2\tau L)$. So,
\[
\text{dist}_{L^1}(\mu(M), \mu(N))(f, s) \leq \frac{\epsilon}{2\tau L} \max \left\{ d_N(f(x), s(x)) : x \in M \right\} \leq \epsilon,
\]
as claimed. □

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