# A volume preserving flow with essential coexistence of zero and non-zero Lyapunov exponents 

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#### Abstract

We demonstrate essential coexistence of hyperbolic and non-hyperbolic behavior in the continuous-time case by constructing a smooth volume preserving flow on a five-dimensional compact smooth manifold that has non-zero Lyapunov exponents almost everywhere on an open and dense subset of positive but not full volume and is ergodic on this subset while having zero Lyapunov exponents on its complement. The latter is a union of three-dimensional invariant submanifolds, and on each of these submanifolds the flow is linear with Diophantine frequency vector.


## 1. Introduction

The goal of this paper is to extend the main result in [10] to dynamical systems with continuous time thus demonstrating coexistence of regular and chaotic dynamics in an 'essential' way.

THEOREM 1.1. (Main theorem) There exists a compact smooth Riemannian manifold $\mathcal{M}$ of dimension five and a $C^{\infty}$ flow $h^{t}: \mathcal{M} \rightarrow \mathcal{M}$ such that the following hold.
(1) $h^{t}$ preserves the Riemannian volume $m$ on $\mathcal{M}$.
(2) $h^{t}(t \neq 0)$ has non-zero Lyapunov exponents (except for the exponent in the flow direction) almost everywhere on an open, dense and connected subset $\mathcal{U} \subset \mathcal{M}$; moreover, $h^{t} \mid \mathcal{U}$ is an ergodic flow.
(3) The complement $\mathcal{U}^{c}$ has a positive volume and is a union of three-dimensional invariant submanifolds; $h^{t}$ is a non-identity linear flow with Diophantine frequency vector on each invariant submanifold, and $h^{t}$ has zero Lyapunov exponents on $\mathcal{U}^{c}$.

We stress that each of the three-dimensional invariant submanifolds is, in turn, a union of two-dimensional invariant tori on which $h^{t}$ is a linear flow with Diophantine frequency vector (see $\S 3$ for details). This fact makes our construction non-trivial $\dagger$.

We emphasize the requirement that the open set $\mathcal{U}$ is everywhere dense $\ddagger$. Donnay [8] constructed an example of a surface on which the geodesic flow exhibits the coexistence phenomenon. It is obtained by inserting a light-bulb cap into a negatively curved surface. In this example the set of geodesics, which are trapped in the cap, is invariant, and has a positive volume, and almost every point in this set has zero Lyapunov exponents. Since it has non-empty interior, the stochastic sea in this example, i.e., the analog of the set $\mathcal{U}$ in our case, is not dense.

While this paper deals only with dynamical systems with continuous time, it is worth mentioning that, in the discrete-time case, the essential coexistence of chaotic and regular behavior has been demonstrated in various situations by Przytycki and Liverani for area preserving diffeomorphisms and by Bunimovich for billiards; see the paper [5], which surveys recent results on essential coexistence, and references therein.

We split the proof of Theorem 1.1 into several steps. In §2, we present some background information and basic notations from the theory of partial hyperbolicity and, in particular, introduce the notion of pointwise partial hyperbolicity on open sets for flows. In §3, we construct the manifold $\mathcal{M}$ and the open set $\mathcal{U}$, and introduce the 'start-up' flow $f^{t}$ that satisfies statement (3) of the theorem. In $\S 4$, we construct a volume preserving flow $g^{t}$, which is a small perturbation of $f^{t}$ and does not affect the action of $f^{t}$ on the set $\mathcal{U}^{c}$. The flow $g^{t}$ has non-zero Lyapunov exponents on a subset of positive volume in $\mathcal{U}$. Finally, in $\S 5$, we construct the desired flow $h^{t}$ as a small perturbation of the flow $g^{t}$. The proof of the main theorem is given in §6.

In our construction of the flows $g^{t}$ and $h^{t}$ we use the perturbation techniques developed in $[\mathbf{1 0}]$ for the case of diffeomorphisms (these techniques originated in $[\mathbf{6}, \mathbf{7}, \mathbf{1 5}]$ ). However, there is a crucial difference between the discrete-time and continuous-time cases. To effect our construction we ought to make perturbations of the vector fields that generate the required flows, but we want these perturbations to produce similar effects on the time-1 maps of the flows, as in [10]. This is made possible due to the crucial fact that the 'start-up' flow $f^{t}$ has a global cross-section, and in our construction we ensure that both perturbation flows $g^{t}$ and $h^{t}$ preserve this cross-section. We achieve this by using specific formulae for perturbations of the vector fields. At the core of our construction lies a new concept of pointwise partially hyperbolic flows on open sets. In §2, we introduce such flows, and we study their ergodicity.
$\dagger$ Indeed, consider the $C^{\infty}$ volume preserving diffeomorphism $P$ constructed in [10]. It has non-zero Lyapunov exponents almost everywhere on an open, dense and connected subset $\mathcal{U} \subset \mathcal{M}$, and $P \mid \mathcal{U}$ is ergodic. Furthermore, the complement $\mathcal{U}^{c}$ has positive volume, and $P \mid \mathcal{U}^{c}$ is the identity map and has zero Lyapunov exponents. A special flow $P^{t}$ over $P$ has non-zero Lyapunov exponents (except for the exponent in the flow direction) almost everywhere on an open, dense and connected subset $\mathcal{U} \times[0,1] / \sim$ (where $\sim$ means that the points $(x, 1)$ and $(P x, 0)$ are identified) which is not of full volume. However, $P^{t}$ is a periodic flow on its complement $\mathcal{U}^{c} \times[0,1] / \sim$.
$\ddagger$ To some extent this justifies the view of our example as 'KAM-type', despite the fact that the flow $h^{t}$ is not close to a completely integrable one.

## 2. Preliminaries

See $[\mathbf{2}, \mathbf{1 0}, \mathbf{1 3}]$ for more details.
Consider a diffeomorphism $f$ acting on a compact smooth Riemannian manifold $\mathcal{M}$. It is called uniformly partially hyperbolic on a compact invariant subset $\Lambda \subset \mathcal{M}$ if the following hold.
(1) For every $x \in \Lambda$ the tangent space at $x$ admits an invariant splitting

$$
\begin{equation*}
T_{x} \mathcal{M}=E^{s}(x) \oplus E^{c}(x) \oplus E^{u}(x) \tag{2.1}
\end{equation*}
$$

into stable $E^{s}(x)=E_{f}^{s}(x)$, central $E^{c}(x)=E_{f}^{c}(x)$ and unstable $E^{u}(x)=E_{f}^{u}(x)$ subspaces.
(2) There are numbers $0<\lambda<\tilde{\lambda} \leq 1 \leq \tilde{\mu}<\mu$ such that for every $t=1,2, \ldots$

$$
\begin{align*}
\left\|d f^{t} v\right\| \leq \lambda^{t}\|v\|, & v \in E^{s}(x), \\
\tilde{\lambda}^{t}\|v\| \leq\left\|d f^{t} v\right\| \leq \tilde{\mu}^{t}\|v\|, & v \in E^{c}(x),  \tag{2.2}\\
\mu^{t}\|v\| \leq\left\|d f^{t} v\right\|, & v \in E^{u}(x) .
\end{align*}
$$

In this paper we need a weaker property than uniform partial hyperbolicity. Let $\mathcal{S} \subset \mathcal{M}$ be an invariant open subset. We say that a diffeomorphism $F$ is pointwise partially hyperbolic on $\mathcal{S}$ if for every $x \in \mathcal{S}$ the tangent space at $x$ admits an invariant splitting (2.1) and there are continuous functions $0<\lambda(x)<\tilde{\lambda}(x) \leq 1 \leq \tilde{\mu}(x)<\mu(x)$ such that (2.2) holds with constants $\lambda, \tilde{\lambda}, \tilde{\mu}$ and $\mu$ replaced with these functions. Pointwise partially hyperbolic diffeomorphisms on compact manifolds were introduced in [4] where their ergodic properties were studied. If a diffeomorphism is pointwise partially hyperbolic on an open subset of a compact manifold then it could fail to have 'nice' properties and, in particular, could be not ergodic (see the discussion below).

We now consider a smooth flow $f^{t}$ on $\mathcal{M}$ which is generated by the vector field $\mathcal{X}_{f}(x)=\left.(d / d t) f^{t}(x)\right|_{t=0}$. We say that the flow is uniformly partially hyperbolic on a compact invariant subset $\Lambda \subset \mathcal{M}$ if for every $x \in \mathcal{M}$ the tangent space at $x$ admits an invariant splitting (2.1) into stable $E^{s}(x)=E_{f}^{s}(x)$, central $E^{c}(x)=E_{f}^{c}(x)$ and unstable $E^{u}(x)=E_{f}^{u}(x)$ subspaces such that the vector field $\mathcal{X}_{f}(x)$ is contained in the central subspace $E^{c}(x)$ and there are numbers $0<\lambda<\tilde{\lambda} \leq 1 \leq \tilde{\mu}<\mu$ such that (2.2) holds for all $t \in[0,1]$. Note that if a flow $f^{t}$ is uniformly partially hyperbolic, then for every $t \neq 0$ the time- $t$ map is uniformly partially hyperbolic with the same invariant splitting.

Given an invariant open subset $\mathcal{S} \subset \mathcal{M}$ we call a flow $f^{t}$ pointwise partially hyperbolic on $\mathcal{S}$ if its time-1 map $f^{1}$ is pointwise partially hyperbolic on $\mathcal{S}$.

Given $\delta>0$, we say that a flow $g^{t}$ is $\left(C^{1}, \delta\right)$-close to $f^{t}$ on an invariant set $\Lambda$ if $\mathcal{X}_{g}=\mathcal{X}_{f}$ outside $\Lambda$ and $\left\|\mathcal{X}_{g}-\mathcal{X}_{f}\right\| \leq \delta$. Uniformly partially hyperbolic flows form an open set in the $C^{1}$ topology (see Lemma B. 8 in Appendix B).

Given a an open subset $\mathcal{S} \subset \mathcal{M}$, we call a partition $\mathcal{P}$ of $\mathcal{S}$ a $(\delta, q)$-foliation with smooth leaves if there exist continuous functions $\delta=\delta(x)>0, q=q(x)>0$, and an integer $k>0$ such that for each $x \in \mathcal{S}$ the following hold.
(1) There exists a smooth immersed $k$-dimensional manifold $W(x)$ containing $x$ for which $\mathcal{P}(x)=W(x)$, where $\mathcal{P}(x)$ is the element of the partition $\mathcal{P}$ containing $x$. The manifold $W(x)$ is called the global leaf of the foliation at $x$; the connected
component of the intersection $W(x) \cap B(x, \delta(x))$ that contains $x$ is called the local leaf at $x$ and is denoted by $V(x)$.
(2) There exists a continuous map $\phi_{x}: B(x, q(x)) \rightarrow C^{1}(D, \mathcal{M})$ (where $D$ is the unit ball) such that $V(y)$ is the image of the map $\phi_{x}(y): D \rightarrow \mathcal{M}$ for each $y \in$ $B(x, q(x))$; the number $q(x)$ is called the size of $V(x)$.
We say that a foliation with smooth leaves is absolutely continuous if for almost every $x \in \mathcal{S}$ and almost every $y \in B(x, q(x))$ the conditional measure on the local leaf $V(y)$, generated by the volume $m$ on $\mathcal{M}$ and the partition of $B(x, q(x))$ by the local leaves, is absolutely continuous with respect to the leaf volume $m_{V(y)}$ on $V(y) \dagger$.

Let $W_{1}$ and $W_{2}$ be two foliations of $\mathcal{S}$ with smooth leaves that are transversal to each other at every point $z \in \mathcal{S}$. Let also $\mathcal{S}_{1} \subset \mathcal{S}$ be an open subset. We say that the pair $W_{1}$ and $W_{2}$ has the accessibility property on $\mathcal{S}_{1}$ if any two points $z, z^{\prime} \in \mathcal{S}_{1}$ are accessible via a $(u, s)$-path in $\mathcal{S}$; that is, the following hold.
(1) There exists a collection of points $z_{1}, \ldots, z_{n} \in \mathcal{S}$ such that $z=z_{1}, z^{\prime}=z_{n}$ and $z_{k} \in W_{i}\left(z_{k-1}\right)$ for $i=1,2$ and $k=2, \ldots, n$.
(2) The points $z_{k-1}$ and $z_{k}$ can be connected by a smooth curve $\gamma_{k} \subset V_{i}\left(z_{k-1}\right)$ in $\mathcal{S}$ for $i=1$ or 2 and $k=2, \ldots, n \ddagger$.
The collection of the leaf-wise paths $\gamma_{k}$ is called a ( $u, s$ )-path and is denoted by $\left[z_{1}, \ldots, z_{n}\right]$.

For a uniformly partially hyperbolic flow $f^{t}$, one can construct stable and unstable local manifolds of uniform size at every point in $\Lambda$. This may not be true for a flow that is pointwise partially hyperbolic on an open set $\mathcal{S}$. However, all pointwise partially hyperbolic flows that we consider in this paper will have global stable and unstable transverse foliations with smooth leaves. We denote these foliations by $W^{s}=W_{f}^{s}$ and $W^{u}=W_{f}^{u}$ respectively.

More precisely, let $f^{t}$ be a flow that is pointwise partially hyperbolic on an open set $\mathcal{S}$ and let $g^{t}$ be a sufficiently small perturbation of $f^{t}$ in the $C^{1}$ topology.

Definition 2.1. We call the perturbation $g^{t}$ gentle if there exists an open set $\mathcal{U} \subset \mathcal{S}$ such that $\overline{\mathcal{U}} \subset \mathcal{S}, \mathcal{U}$ is invariant under both $f^{t}$ and $g^{t}$ and $f^{t}\left|\mathcal{U}^{c}=g^{t}\right| \mathcal{U}^{c}$.

THEOREM 2.2. Assume that the strongly stable and unstable subspaces $E_{f^{t}}^{s}$ and $E_{f^{t}}^{u}$ for $f^{t}$ are integrable to continuous strongly stable and unstable foliations $W_{f^{t}}^{s}$ and $W_{f^{t}}^{u}$, respectively, with smooth leaves, and that these foliations are transverse. Then, for any sufficiently small gentle perturbation $g^{t}$ in the $C^{1}$ topology, the strongly stable and unstable subspaces, $E_{g^{t}}^{s}$ and $E_{g^{t}}^{u}$ for $g^{t}$ are integrable to continuous strongly stable and unstable foliations $W_{g^{t}}^{s}$ and $W_{g^{t}}^{t}$, respectively, with smooth leaves, and these foliations are transverse.

We call a flow $f^{t}$ that is pointwise partially hyperbolic on an open set $\mathcal{S}$ dynamically coherent if the subbundles $E^{c u}=E^{c} \oplus E^{u}, E^{c}$ and $E^{c s}=E^{c} \oplus E^{s}$, are integrable to
$\dagger$ The leaf volume $m_{V(y)}$ is generated by the restriction of the Riemannian metric on $\mathcal{M}$ to the smooth submanifold $V(y)$.
$\ddagger$ We stress that $V_{i}\left(z_{k-1}\right)$ is the local leaf of $W_{i}$ at $z_{i}$. In particular, the length of the curve $\gamma_{k}$ (the leg of the path) does not exceed $\delta\left(z_{k-1}\right)$.
continuous foliations with smooth leaves $W^{c u}, W^{c}$ and $W^{c s}$, called respectively the centerunstable, center and center-stable foliations. Furthermore, the foliations $W^{c}$ and $W^{u}$ are subfoliations of $W^{c u}$, while $W^{c}$ and $W^{s}$ are subfoliations of $W^{c s}$.

The following result is an extension of the classical result in $[\mathbf{9}, \mathbf{1 4}]$ to the case of flows that are pointwise partially hyperbolic on an open subset. It shows that dynamical coherence is a robust property within the class of gentle perturbations. The proof of this result is a simple modification of the argument in [9].

Theorem 2.3. Suppose that $f^{t}$ is a flow that is pointwise partially hyperbolic on an open set $\mathcal{S}$. Assume that $f^{t}$ possesses transverse strongly stable and unstable foliations with smooth leaves. Assume also that the center distribution is integrable to a smooth center foliation $W^{c}$. Then $f^{t}$ is dynamically coherent. Moreover, any flow that is close to $f^{t}$ in the $C^{1}$ topology and is a gentle perturbation of $f^{t}$ is dynamically coherent.

Since both subbundles $E^{c u}$ and $E^{c s}$ vary continuously with the map, so does $E^{c}$ and the corresponding center foliation $W^{c}$.

Given a smooth flow $f^{t}$, we denote by

$$
\lambda(x, v)=\lambda\left(x, v, f^{t}\right)=\limsup _{t \rightarrow \infty} \frac{1}{t} \log \left\|d f^{t} v\right\|
$$

the Lyapunov exponent of a non-zero vector $v$ at $x \in \mathcal{M}$ and by $\lambda_{i}(x)=\lambda_{i}\left(x, f^{t}\right)$, $i=1, \ldots, \operatorname{dim} \mathcal{M}$, the values of the Lyapunov exponents at $x$ in decreasing order. We also write

$$
\begin{equation*}
L_{k}\left(f^{t}\right)=\int_{\mathcal{M}} \sum_{i=1}^{k} \lambda_{i}\left(x, f^{t}\right) d m(z) \tag{2.3}
\end{equation*}
$$

where $m$ is the Riemannian volume. We call this number the $k$ th average Lyapunov exponent of $f^{t}$.

Consider a $C^{2}$ flow $f^{t}$ of a compact smooth manifold $\mathcal{M}$ that is pointwise partially hyperbolic on an open invariant set $\mathcal{S}$. Assume that $f^{t}$ preserves a smooth measure on $\mathcal{M}$. We say that $f^{t}$ has positive central exponents if there is an invariant set $\mathcal{A} \subset \mathcal{S}$ of positive measure such that for every $x \in \mathcal{A}$ and every $v \in E^{c}(x) \backslash \operatorname{Span}\left\{\mathcal{X}_{f}(x)\right\}$ the Lyapunov exponent $\lambda(x, v)$ is positive. The following theorem plays an important role in the proof of our main theorem.

THEOREM 2.4. Assume that the following conditions hold.
(1) $f^{t}$ has strongly stable and unstable $(\delta, q)$-foliations $W^{s}$ and $W^{u}$, where $\delta=\delta(x)$ and $q=q(x)$ are continuous functions on $\mathcal{S}$.
(2) The foliations $W^{s}$ and $W^{u}$ are absolutely continuous.
(3) $f^{t}$ has the accessibility property via the foliations $W^{s}$ and $W^{u}$ on $\mathcal{S}$.
(4) $f^{t}$ has positive central exponents.
(5) The Lyapunov exponents in the stable subspace $E^{s}(x)$ are all negative and the Lyapunov exponents in the unstable subspace $E^{u}(x)$ are all positive for almost every $x$.

Then $f^{t}$ has positive central exponents at almost every point $x \in \mathcal{S}$, and $f^{t} \mid \mathcal{S}$ is an ergodic flow.

Proof. Following [10, Theorem 2.2], we first show that $f^{t}$ has positive central exponents at almost every point $x \in \mathcal{S}$ and that $f^{t} \mid \mathcal{S}$ is an ergodic flow.

There is a set $\mathcal{A} \subset \mathcal{S}$ of positive measure such that the flow $f^{t} \mid \mathcal{A}$ has non-zero Lyapunov exponents except along the flow direction. Hence, it has at most countably many ergodic components of positive measure in $\mathcal{A}$ (see [2]). Each such component contains the set

$$
A(x)=\bigcup_{y \in V^{c u}(x)} V^{s}(y)
$$

where $x$ is the density point of $\mathcal{A}$ and $V^{c u}(x)$ is a center-unstable local manifold at $x$. Since the strong stable foliation $W^{s}$ is continuous, the set $A(x)$ is open $(\bmod 0)$ in $\mathcal{M} \dagger$ and hence the set $\mathcal{A}$ is itself open $(\bmod 0)$. We shall show that the trajectory of almost every point in $\mathcal{S}$ is dense, which yields that $\mathcal{A}=\mathcal{S}(\bmod 0)$ and that $f^{t} \mid S$ is ergodic.

The proof of this claim for partially hyperbolic diffeomorphisms is given in [1] and can be extended to our case literally. We present the argument here for the reader's convenience. We call a point $p$ good for a given open set $U$ if $p$ has a neighborhood in which the orbit of almost every point enters $U$. It suffices to show that an arbitrary point $p$ is good. Since $f^{t}$ is accessible, there is a $(u, s)$-path $\left[z_{0}, \ldots, z_{k}\right]$ with $z_{0} \in U$ and $z_{k}=p$. We will show by induction on $j$ that each point $z_{j}$ is good. This is obvious for $j=0$. Now suppose that $z_{j}$ is good; then $z_{j}$ has a neighborhood $N$ such that $\operatorname{Orb}(x) \cap U \neq \emptyset$ for almost every $x \in N$. Let $B$ be the subset of $N$ consisting of points with this property that are also both forward and backward recurrent. It follows from the Poincaré recurrence theorem that $B$ has full measure in $N$. If $x \in B$, any point $y \in W^{s}(x) \cup W^{u}(x)$ has the property that $\operatorname{Orb}(y) \cap U \neq \emptyset$. The absolute continuity of the foliations $W^{s}$ and $W^{u}$ means that the set

$$
\bigcup_{x \in B} W^{s}(x) \cup W^{u}(x)
$$

has full measure in the set

$$
\bigcup_{x \in N} W^{s}(x) \cup W^{u}(x) .
$$

The latter is a neighborhood of $z_{j+1}$. Hence, $z_{j+1}$ is good.
3. Construction of the 'start-up' flow $f^{t}$

Let $A$ be an Anosov automorphism of the 2-torus $X=\mathbb{T}^{2}$ with expanding rate $\eta>1$ along the unstable direction. Consider the suspension flow $S^{t}$ on the suspension manifold $\mathcal{N}=X \times \mathbb{R} / \sim$, with the identification $(x, \tau+1) \sim(A x, \tau)$. The action of the suspension flow on $\mathcal{N}$ is exactly $S^{t}(x, \tau)=(x, \tau+t)$. See Appendix A for more details of the geometric structure of $\mathcal{N}$.

Given $\alpha \in \mathbb{T}^{2}$, let $T_{\alpha}^{t}: \mathcal{N} \rightarrow \mathcal{N}$ be a linear flow defined by $(x, \tau) \mapsto(x+t \alpha, \tau)$. It preserves each level set $X \times\{\tau\}$.
$\dagger$ That is, there is an open set $\mathcal{V} \subset \mathcal{M}$ such that $A(x)=\mathcal{V}(\bmod 0)$ with respect to the volume $m$.

Set $Y=\mathbb{T}^{2}$ and $\mathcal{M}=\mathcal{N} \times Y$. To effect our construction we choose:
(A1) a Cantor set $C \subset Y$ of positive area whose complement $U=Y \backslash C$ is a non-empty open and connected set;
(A2) an open square $U_{0}$ such that $\overline{U_{0}} \subset U$;
(A3) an open neighborhood $U_{1}$ of $\overline{U_{0}}$ such that $\overline{U_{1}} \subset U$, whose choice will be specified in §5.2(1).
We also choose a $C^{\infty}$ function $\kappa: Y \rightarrow \mathbb{R}$ such that:
( $\kappa 1) ~ \kappa(y)>0$ for $y \in U$ and $\kappa(y)=0$ for $y \in C$;
(к2) $\kappa(y)=1$ for $y \in U_{1}$;
(к3) $\|\kappa\|_{C^{1}} \leq 1$.
(The existence of such a function $\kappa$ follows from the specific construction of the Cantor set and the choice of the set $U_{1}$ in §5.1.) We also choose a $C^{\infty} \operatorname{map} \alpha: Y \rightarrow \mathbb{R}^{2}$ such that:
$(\alpha 1) \alpha(y)=0$ for $y \in U_{1}$;
( $\alpha 2$ ) $\alpha(y)=\alpha_{0}>0$ for all $y \in C$ where $\alpha_{0}$ is a Diophantine vector;
( $\alpha 3$ ) $\sup _{y \in Y}\|\alpha(y)\| \leq \bar{\alpha}$, where $\bar{\alpha}$ is a positive number determined in §5.2.
We now set $\mathcal{U}=\mathcal{N} \times U$ and $\mathcal{U}^{c}=\mathcal{N} \times C$ and define the flow $f^{t}$ on $\mathcal{M}$ by the formula

$$
\begin{equation*}
f^{t}((x, \tau), y)=((x+t \alpha(y), \tau+t \kappa(y)), y) \tag{3.1}
\end{equation*}
$$

where $(x, \tau) \in \mathcal{N}$ and $y \in Y$. The following proposition describes the properties of the flow $f^{t}$ and its proof follows immediately from the definitions.
Proposition 3.1. The following statements hold.
(1) $f^{t}$ is a $C^{\infty}$ volume preserving flow.
(2) $f^{t}$ preserves each fiber $\mathcal{N} \times\{y\}$, on which $f^{t}$ is the composition of the scaled suspension flow $S^{t \kappa(y)}$ and the linear flow $T_{\alpha(y)}^{t}$; in particular, $f^{t}$ is exactly the suspension flow $S^{t}$ on $\mathcal{N} \times\{y\}$ for $y \in \overline{U_{1}}$.
(3) $f^{t}$ is pointwise partially hyperbolic on $\mathcal{U}$ with one-dimensional stable $E_{f}^{s}$, onedimensional unstable $E_{f}^{u}$ and three-dimensional center $E_{f}^{c}$ subbundles; $E_{f}^{s}$ and $E_{f}^{u}$ are integrable to strongly stable and unstable foliations $W_{f}^{s}$ and $W_{f}^{u}$ with smooth leaves, which are absolutely continuous, uniformly transversal and have local leaves of uniform size.
(4) $f^{t}$ is uniformly partially hyperbolic on $\mathcal{N} \times A$ where $A \subset U$ is a subset, and hence $f^{t}$ is dynamically coherent with the central foliation $W_{f}^{c}=W_{S^{t}}^{c} \times Y$.
(5) $f^{t}$ preserves every two-dimensional torus $X \times\{\tau\} \times\{y\}(\tau \in[0,1], y \in C$ are fixed) and acts on it as a linear flow with a Diophantine frequency vector; moreover, $f^{t} \mid \mathcal{U}^{c}$ has all zero Lyapunov exponents on $\mathcal{U}^{c}$.
(6) For every $z=((x, \tau), y) \in \mathcal{M}$ the Lyapunov exponents of $f^{t}$ are as follows:

$$
\begin{aligned}
\lambda_{1}\left(z, f^{t}\right)=\kappa(y) \log \eta \geq 0 & =\lambda_{2}\left(z, f^{t}\right)=\lambda_{3}\left(z, f^{t}\right)=\lambda_{4}\left(z, f^{t}\right) \\
& \geq \lambda_{5}\left(z, f^{t}\right)=-\kappa(y) \log \eta
\end{aligned}
$$

Moreover, if $z \in \mathcal{U}$, then $\lambda_{1}\left(z, f^{t}\right)=\lambda^{u}\left(z, f^{t}\right)>0$ is the Lyapunov exponent in the $E_{f}^{u}(z)$ subspace, $\lambda_{5}\left(z, f^{t}\right)=\lambda^{s}\left(z, f^{t}\right)<0$ is the Lyapunov exponent in the $E_{f}^{s}(z)$ subspace, and $\lambda_{2}\left(z, f^{t}\right), \lambda_{3}\left(z, f^{t}\right)$ and $\lambda_{4}\left(z, f^{t}\right)$ are Lyapunov exponents in the flow direction and two directions in $Y$ respectively.

## 4. Removing zero exponents

In this section we will construct a gentle perturbation $g^{t}$ of the original flow $f^{t}$ with positive central Lyapunov exponents on a set of positive volume but not necessarily ergodic. Then we will perturb $g^{t}$ to the desired flow $h^{t}$ of the main theorem in $\S 5$.

Given $z \in \mathcal{M}$, there is a local Cartesian coordinate system ( $u, s, \tau, a, b$ ) (see Appendix A) such that

$$
F^{u}(z):=\frac{\partial}{\partial u}=E_{f}^{u}(z), \quad F^{s}(z):=\frac{\partial}{\partial s}=E_{f}^{s}(z), \quad F^{\tau}(z):=\frac{\partial}{\partial \tau}=E_{f}^{\tau}(z)
$$

are the unstable, stable and flow directions of $f^{t}$ respectively, and

$$
F^{a}(z):=\frac{\partial}{\partial a}=E_{f}^{a}(z), \quad F^{b}(z):=\frac{\partial}{\partial b}=E_{f}^{b}(z)
$$

are the other two central directions tangent to $Y$.
The following statement describes properties of the flow $g^{t}$.
Proposition 4.1. Given $\delta_{g}>0$, there is a $C^{\infty}$ volume preserving flow $g^{t}$ on $\mathcal{M}$ such that the following hold.
(1) $g^{t}$ is $\left(C^{1}, \delta_{g}\right)$-close to $f^{t}$; i.e., $\left\|\mathcal{X}_{f}-\mathcal{X}_{g}\right\| \leq \delta_{g}$, where $\mathcal{X}_{f}$ and $\mathcal{X}_{g}$ are the vector fields of the flows $f^{t}$ and $g^{t}$ respectively.
(2) $g^{t}=f^{t}$ outside $\mathcal{N} \times U_{0}$, and hence $g^{t}$ is a gentle perturbation of $f^{t}$ and satisfies statements (3)-(5) of Proposition 3.1.
(3) $g^{t}$ preserves the subbundles $E_{f}^{\omega}, \omega=u a b, u a b \tau$; moreover,

$$
\begin{equation*}
\operatorname{det}\left(d g^{t} \mid E_{f}^{\omega}(z)\right)=\operatorname{det}\left(d f^{t} \mid E_{f}^{\omega}(z)\right) \quad \text { for all } z \in \mathcal{M} \tag{4.1}
\end{equation*}
$$

(4) The average Lyapunov exponents of $g^{t}$ satisfy

$$
\begin{equation*}
L_{5}\left(g^{t}\right)=0<L_{1}\left(g^{t}\right)<L_{2}\left(g^{t}\right)<L_{3}\left(g^{t}\right)=L_{4}\left(g^{t}\right) . \tag{4.2}
\end{equation*}
$$

To prove this proposition, we extend the approach in [10] to the case of flows and obtain the flow $g^{t}$ as a result of two consecutive perturbations. First, we perturb the start-up flow $f^{t}$ to a flow $\widetilde{g}^{t}$ by adding a rotational vector field $\widetilde{\mathcal{X}}_{R}$ to the vector field $\mathcal{X}_{f} \dagger$. This produces two positive average Lyapunov exponents for the flow $\widetilde{g}^{t}$ in the $E_{f}^{u a}$ subbundle, i.e., $L_{1}\left(\widetilde{g}^{t}\right)<L_{2}\left(\widetilde{g}^{t}\right)$. Next, we perturb $\widetilde{g}^{t}$ to the desired flow $g^{t}$ by adding another rotational vector field $\mathcal{X}_{R}$ to the vector field $\mathcal{X}_{\widetilde{g}}$ for the flow $\widetilde{g}^{t}$. As a result the flow $g^{t}$ has three positive average Lyapunov exponents in the $E_{f}^{u a b}$ subbundle, i.e., $L_{1}\left(g^{t}\right)<L_{2}\left(g^{t}\right)<L_{3}\left(g^{t}\right)$.

The vector fields $\widetilde{\mathcal{X}}_{R}$ and $\mathcal{X}_{R}$ are chosen to be supported on disjoint open subsets $\widetilde{\Omega}_{R}$ and $\Omega_{R}$ of $\mathcal{N} \times U_{0}$ respectively such that $\widetilde{\mathcal{X}}_{R}=0$ outside $\widetilde{\Omega}_{R}, \mathcal{X}_{R}=0$ outside $\Omega_{R}$ and $\left\|\widetilde{\mathcal{X}}_{R}\right\|_{C^{1}},\left\|\mathcal{X}_{R}\right\|_{C^{1}}<\delta_{g} / 2$. Since $\mathcal{N} \times U_{0}$ is invariant under $f^{t}$, we have that $g^{t}=\widetilde{g}^{t}=f^{t}$ outside $\mathcal{N} \times U_{0}$.

Our construction utilizes the following crucial feature of the flow $f^{t}$ : the set

$$
\begin{equation*}
\Pi_{0}=X \times\{0\} \times U_{0} \tag{4.3}
\end{equation*}
$$

$\dagger$ That is, the flow generated by $\widetilde{\mathcal{X}}_{R}$ is concentrated in a small neighborhood of a point in $\mathcal{M}$ where it acts as a small rotation around this point; see (4.6).
is a global cross-section of $f^{t} \mid \mathcal{N} \times U_{0}$, and the time-1 map restricted to $\Pi_{0}$ is exactly the Poincaré return map of $f^{t}$ to $\Pi_{0}$. Furthermore, we make the construction of vector fields $\tilde{\mathcal{X}}_{R}$ and $\mathcal{X}_{R}$ in such a way that $\Pi_{0}$ is also a global cross-section for both flows $\widetilde{g}^{t}$ and $g^{t}$ with the time-1 maps to be the Poincaré return map to $\Pi_{0}$. This fact allows us to apply arguments similar to those in [10] to our flow case by focusing on the time- 1 maps.
4.1. Construction of the flow $\widetilde{g}^{t}$. In this section we construct the flow $\widetilde{g}^{t}$ by perturbing the vector field $\mathcal{X}_{f}$ inside the set $\mathcal{N} \times U_{0}$.

To effect our construction we choose distinct periodic points $q, p^{a}, p^{b}$ and $p^{\tau}$ of the Anosov automorphism $A$ of $X$, which are close to each other. Let $V_{A}^{s}(q), V_{A}^{u}(q), V_{A}^{s}\left(p^{i}\right)$ and $V_{A}^{u}\left(p^{i}\right), i=a, b, \tau$, be the stable and unstable local manifolds at these periodic points. We may assume that each intersection $V_{A}^{u}(q) \cap V_{A}^{s}\left(p^{i}\right)$ and $V_{A}^{u}\left(p^{i}\right) \cap V_{A}^{S}(q)$ consists of a single point, which we denote by $\left[q, p^{i}\right]$ and $\left[p^{i}, q\right]$ respectively. Let $\gamma^{i}$ denote the closed quadrilateral path with the collection of points $q,\left[q, p^{i}\right], p^{i},\left[p^{i}, q\right]$ and $q$, and let

$$
\gamma(q)=V_{A}^{u}(q) \cup V_{A}^{s}(q), \quad \gamma\left(p^{i}\right)=V_{A}^{u}\left(p^{i}\right) \cup V_{A}^{s}\left(p^{i}\right) .
$$

Choose $v>0$ and set, for $i=a, b, \tau$,

$$
\begin{align*}
\Omega_{0}^{i}(v) & =\left(\bigcup_{t \in\left[0, \iota\left(p^{i}\right)\right]} B_{\mathcal{N}}\left(f^{t}\left(p^{i}, 0\right), v\right)\right) \times U_{0}, \\
\widehat{\Omega}_{0}^{i}(v) & =\left(\bigcup_{(x, \tau) \in(\gamma(q) \times[0, \iota(q)]) \cup\left(\gamma\left(p^{i}\right) \times\left[0, \iota\left(p^{i}\right)\right]\right)} B_{\mathcal{N}}((x, \tau), v)\right) \times U_{0},  \tag{4.4}\\
\Omega_{0} & =\Omega_{0}(\nu)=\left(\bigcup_{i=a, b, \tau} \Omega^{i}(\nu)\right) \cup\left(\bigcup_{i=a, b, \tau} \widehat{\Omega}^{i}(v)\right),
\end{align*}
$$

where $\iota(q)$ and $\iota\left(p^{i}\right)$ are the periods of $q$ and $p^{i}$ respectively, and $B_{\mathcal{N}}((x, \tau), r)$ is the ball in $\mathcal{N}$ of radius $r$ centered at the point $(x, \tau) \in \mathcal{N}$. We choose a sufficiently small number $v$ to ensure that

$$
m\left(\operatorname{Proj}_{\Pi_{0}} \Omega_{0}\right) \leq 0.05 m\left(\Pi_{0}\right)
$$

where $\Pi_{0}$ is given by (4.3), and $\operatorname{Proj}_{\Pi_{0}}: \mathcal{N} \times U_{0} \rightarrow \Pi_{0}$ is the natural projection onto $\Pi_{0}$ given by the formula

$$
\operatorname{Proj}_{\Pi_{0}}((x, \tau), y)=((x, 0), y)
$$

To construct the vector field $\widetilde{\mathcal{X}}_{R}$ we choose a $C^{\infty}$ function $\psi: \mathbb{R} \rightarrow[0,1]$ such that:
(1) $\psi=1$ on $(-0.9,0.9)$;
(2) $\psi>0$ on $(-1,1)$ and $\psi=0$ outside $(-1,1)$;
(3) $\|\psi\|_{C^{1}} \leq 10$.

Observe that $\mathcal{N} \times U_{0}$ is invariant under the flow $f^{t}$ and that

$$
f^{t}((x, \tau), y)=((x, \tau+t), y)
$$

for $((x, \tau), y) \in \mathcal{N} \times U_{0}$. In other words, $f^{t} \mid \mathcal{N} \times U_{0}$ is the product of the suspension flow on $\mathcal{N}$ and the identity map on $U_{0}$. It follows that $\Pi_{0}$ is a global cross-section for the flow $f^{t} \mid \mathcal{N} \times U_{0}$, and the time-1 map restricted to $\Pi_{0}$ is $f^{1}=A \times \mathrm{Id}$ and is
exactly the Poincaré return map of $f^{t}$ to $\Pi_{0}$. Given a subset $\Pi \subset \Pi_{0}$, we call a set $\Pi \times\left[\tau_{1}, \tau_{2}\right] \subset \Pi_{0} \times \mathbb{R} / \sim=\mathcal{N} \times U_{0}$ a tube if

$$
f^{t}\left(\Pi \times\left\{\tau_{1}\right\}\right) \cap\left(\Pi \times\left\{\tau_{1}\right\}\right)=\emptyset \quad \text { for all } t \in\left[0, \tau_{2}-\tau_{1}\right]
$$

It is easy to check that $\Pi \times\left[\tau_{1}, \tau_{2}\right]$ is a tube if and only if the sets $\Pi, f^{1}(\Pi), f^{2}(\Pi), \ldots, f^{l}(\Pi)$ are pairwise disjoint, where $l=\left\lfloor\tau_{2}-\tau_{1}\right\rfloor$. Choose a nonperiodic point $z_{0}=\left(x_{0}, 0, y_{0}\right) \in \Pi_{0} \backslash \operatorname{Proj}_{\Pi_{0}} \Omega_{0}$, where $x_{0}$ is a non-periodic point of the Anosov automorphism $A, y_{0}$ is the center of the square $U_{0}$ and $\Omega_{0}$ is the set given by (4.4).

In what follows we will use the local Cartesian coordinate system in a neighborhood of $\Pi_{0}$ originated at $z_{0}$ given by $(u, s, t, a, b)$, where $(u, s)$ are the coordinates in $X$ along the stable and unstable directions of the hyperbolic diffeomorphism $A, t$ is the coordinate along the time direction and $(a, b)$ are coordinates in $Y$. In this coordinate system a point $z \in \Pi_{0}$ is given as $z=(u, s, 0, a, b)$. We will also use the $u a$-cylindrical coordinates $(r, \theta, \tau, s, b)$, where $u=r \cos \theta, a=r \sin \theta$. Given $\varepsilon>0$, one can choose a $u a$-cylinder $B \subset \Pi_{0}$ centered at $z_{0}$ of size $\varepsilon$, i.e.,

$$
B=\{(r, \theta, 0, s, b): r \leq \varepsilon,|s| \leq \varepsilon,|b| \leq \varepsilon\} .
$$

Given a sufficiently large $N_{0} \geq 20 k_{0}$ (the number $k_{0}$ is defined below in Lemma 4.5), we can choose $\varepsilon$ so small that $f^{i}(B) \cap B=\emptyset$ for $i=1, \ldots, N_{0}$. Consider the tube

$$
\begin{equation*}
\widetilde{\Omega}_{R}=B \times[0,1 / 2] . \tag{4.5}
\end{equation*}
$$

Since $z_{0} \notin \operatorname{Proj}_{\Pi_{0}} \Omega_{0}$, we can further reduce $\varepsilon$ to ensure that $B \cap \operatorname{Proj}_{\Pi_{0}}\left(\Omega_{0}\right)=\emptyset$. Hence, $\widetilde{\Omega}_{R} \cap \Omega_{0}=\emptyset$.

Given $\beta>0$, define a $C^{\infty}$ rotational vector field $\tilde{\mathcal{X}}_{R}=\tilde{\mathcal{X}}_{R, \beta}$ on $\mathcal{M}$ as follows:

$$
\tilde{\mathcal{X}}_{R, \beta}(z)= \begin{cases}\beta \tilde{\psi}(z) \frac{\partial}{\partial \theta}, & z \in \widetilde{\Omega}_{R}  \tag{4.6}\\ 0, & z \in \mathcal{M} \backslash \widetilde{\Omega}_{R}\end{cases}
$$

where

$$
\tilde{\psi}(z)=\widetilde{\psi}(r, \theta, \tau, s, b)=\psi\left(\frac{r^{2}}{\varepsilon^{2}}\right) \psi\left(\frac{s}{\varepsilon}\right) \psi\left(\frac{b}{\varepsilon}\right) \psi\left(\frac{\tau-(1 / 4)}{1 / 4}\right) .
$$

It is easy to see that $\|\tilde{\psi}(\partial / \partial \theta)\| \leq c$ where $c>0$ is a constant, which is independent of $\varepsilon$. Hence, $\left\|\tilde{\mathcal{X}}_{R, \beta}\right\| \rightarrow 0$ as $\beta \rightarrow 0$. Furthermore, $\tilde{\mathcal{X}}_{R, \beta}$ is divergence free. Let $\widetilde{g}^{t}=\widetilde{g}_{\beta}^{t}$ be the flow generated by the vector field

$$
\mathcal{X}_{\tilde{g}}=\mathcal{X}_{\widetilde{g}, \beta}=\mathcal{X}_{f}+\widetilde{\mathcal{X}}_{R, \beta}
$$

Proposition 4.2. There exists $\beta>0$ such that $\tilde{g}^{t}=\tilde{g}_{\beta}^{t}$ is a $C^{\infty}$ volume preserving flow with the following properties.
(1) $\widetilde{g}^{t}$ is $\left(C^{1}, \delta_{g} / 2\right)$-close to $f^{t}$; i.e., $\left\|\mathcal{X}_{f}-\mathcal{X}_{\widetilde{g}}\right\|_{C^{1}} \leq \delta_{g} / 2$, where $\mathcal{X}_{f}$ and $\mathcal{X}_{\widetilde{g}}$ are the vector fields corresponding to the flows $f^{t}$ and $\widetilde{g}^{t}$ respectively.
(2) $\tilde{g}^{t}=f^{t}$ outside $\mathcal{N} \times U_{0}$, and hence $\tilde{g}^{t}$ is a gentle perturbation of $f^{t}$ and satisfies statements (3)-(5) of Proposition 3.1.
(3) $\tilde{g}^{t}$ preserves the subbundles $E_{f}^{\omega}, \omega=u a, u a b, u a b \tau$; moreover,

$$
\begin{equation*}
\operatorname{det}\left(d \widetilde{g}^{t} \mid E_{f}^{\omega}(z)\right)=\operatorname{det}\left(d f^{t} \mid E_{f}^{\omega}(z)\right) \quad \text { for all } z \in \mathcal{M} \tag{4.7}
\end{equation*}
$$

(4) The average Lyapunov exponents of $\widetilde{g}^{t}$ satisfy

$$
\begin{equation*}
L_{5}\left(\widetilde{g}^{t}\right)=0<L_{1}\left(\widetilde{g}^{t}\right)<L_{2}\left(\widetilde{g}^{t}\right)=L_{3}\left(\widetilde{g}^{t}\right)=L_{4}\left(\widetilde{g}^{t}\right) \tag{4.8}
\end{equation*}
$$

(5) $\quad \Pi_{0}$ is a global cross-section of the flow $\widetilde{g}^{t} \mid \mathcal{N} \times U_{0}$, and the time- 1 map $\widetilde{g}^{1}$ is the Poincaré return map of $g^{t}$ to $\Pi_{0}$; furthermore, there exist $\lambda>0$ and a $\widetilde{g}^{1}$-invariant subset $\Pi \subset \Pi_{0}$ such that

$$
m(\Pi) \geq 20 k_{0} m(\Pi \cap B)>0
$$

and for any $z \in \Pi$ the flow $\tilde{g}^{t}$ has two positive Lyapunov exponents $\lambda_{1}\left(z, \widetilde{g}^{t}\right)>$ $\lambda_{2}\left(z, \tilde{g}^{t}\right) \geq \lambda$ along the $E_{f}^{u a}$ subbundle.
Proof. Statements (1) and (2) are easy corollaries of the construction of the flow $\widetilde{g}^{t}$. To prove statement (3) we will first show that $d \widetilde{g}^{t}$ preserves the subbundle $E_{f}^{u a}$. It suffices to check that for any smooth vector field $\mathcal{X} \in E_{f}^{u a}$ and any $z \in \widetilde{\Omega}_{R}$, the Lie bracket $\left[\mathcal{X}_{\widetilde{g}}(z), \mathcal{X}(z)\right] \in E_{f}^{u a}(z)$. Indeed, we have

$$
\mathcal{X}_{\widetilde{g}}(z)=\frac{\partial}{\partial \tau}+\beta \widetilde{\psi}(z) \frac{\partial}{\partial \theta}=\frac{\partial}{\partial \tau}+\beta \widetilde{\psi}(z)\left(-a \frac{\partial}{\partial u}+u \frac{\partial}{\partial a}\right),
$$

and the direct calculation yields

$$
\left[\mathcal{X}_{\widetilde{g}}, \frac{\partial}{\partial \omega}\right]=\beta\left(\frac{\partial(a \tilde{\psi})}{\partial \omega} \frac{\partial}{\partial u}-\frac{\partial(u \widetilde{\psi})}{\partial \omega} \frac{\partial}{\partial a}\right) \in E_{f}^{u a}, \quad \omega=u, a
$$

Since $d \widetilde{g}^{t}$ preserves the subbundle $E_{f}^{u a}$, it also preserves the subbundles $E_{f}^{u a b}$ and $E_{f}^{u a b \tau}$. Next, consider the variational differential equations

$$
\frac{d}{d t} f^{t}=D \mathcal{X}_{f} d f^{t}, \quad \frac{d}{d t} \widetilde{g}^{t}=D \mathcal{X}_{\widetilde{g}} d \widetilde{g}^{t}
$$

The determinants along $E_{f}^{\omega}$ with $\omega=u a$, $u a b, u a b \tau$, satisfy

$$
\begin{align*}
\frac{d}{d t} \operatorname{det}\left(d f^{t} \mid E_{f}^{\omega}\right) & =\operatorname{div}\left(\mathcal{X}_{f} \mid E_{f}^{\omega}\right) \operatorname{det}\left(d f^{t} \mid E_{f}^{\omega}\right) \\
\frac{d}{d t} \operatorname{det}\left(d \widetilde{g}^{t} \mid E_{f}^{\omega}\right) & =\operatorname{div}\left(\mathcal{X}_{\widetilde{g}} \mid E_{f}^{\omega}\right) \operatorname{det}\left(d \widetilde{g}^{t} \mid E_{f}^{\omega}\right) \tag{4.9}
\end{align*}
$$

Direct calculations show that $\tilde{\mathcal{X}}_{R}=\mathcal{X}_{\widetilde{g}}-\mathcal{X}_{f}$ is divergence free along $E_{f}^{\omega}$ and thus $\operatorname{div}\left(\mathcal{X}_{f} \mid E_{f}^{\omega}\right)=\operatorname{div}\left(\mathcal{X}_{\widetilde{g}} \mid E_{f}^{\omega}\right)$. Therefore, using (4.9) and the fact that $\operatorname{det}\left(d g^{0} \mid E_{f}^{\omega}\right)=$ $\operatorname{det}\left(d f^{0} \mid E_{f}^{\omega}\right)=1$ we find that

$$
\operatorname{det}\left(d \widetilde{g}^{t} \mid E_{f}^{\omega}(z)\right)=\operatorname{det}\left(d f^{t} \mid E_{f}^{\omega}(z)\right)
$$

and statement (3) follows.
It remains to prove statements (4) and (5). We need the following lemma showing that $\Pi_{0}$ is a global cross-section for the flow $\widetilde{g}^{t} \mid \mathcal{N} \times U_{0}$.

Lemma 4.3. Given $z \in B$, the $r$-, $s$-, $b$ - and $\tau$-coordinates of $\widetilde{g}^{t}(z)$ and $f^{t}(z)$ are the same for $t \in[0,1 / 2]$. Consequently, $\tilde{g}^{\frac{1}{2}}(B)=f^{\frac{1}{2}}(B)$ and $\Pi_{0}$ is a global cross-section for the flow $\widetilde{g}^{t} \mid \mathcal{N} \times U_{0}$.

Proof. Let us compare the orbit segments of $\widetilde{g}^{t}(z)$ and $f^{t}(z)$ for $t \in[0,1 / 2]$. Note that for any smooth function $\varphi$ and any vector field $\mathcal{X}$ we have that

$$
\frac{d}{d t} \varphi\left(F^{t}(z)\right)=\left.L_{\mathcal{X}} \varphi\right|_{F^{t}(z)}
$$

where $L \mathcal{X}(\cdot)$ is the Lie derivative and $F^{t}$ is the flow that is generated by $\mathcal{X}$. This implies that, for $\omega=r, s, b$,

$$
\begin{gathered}
\frac{d}{d t} \tau\left(\widetilde{g}^{t}(z)\right)=L_{\mathcal{X}_{\widetilde{g}}} \tau=\frac{\partial \tau}{\partial \tau}+\beta \widetilde{\psi}\left(\widetilde{g}^{t}(z)\right) \frac{\partial \tau}{\partial \theta}=1=L_{\mathcal{X}_{f}} \tau=\frac{d}{d t} \tau\left(f^{t}(z)\right) \\
\frac{d}{d t} \omega\left(\widetilde{g}^{t}(z)\right)=L_{\mathcal{X}_{\widetilde{g}}} \omega=\frac{\partial \omega}{\partial \tau}+\beta \widetilde{\psi}\left(\widetilde{g}^{t}(z)\right) \frac{\partial \omega}{\partial \theta}=0=L_{\mathcal{X}_{f}} \omega=\frac{d}{d t} \omega\left(f^{t}(z)\right)
\end{gathered}
$$

Under the same initial condition at $t=0$, we get that the $r$-, $s$-, $b$ - and $\tau$-coordinates of $\widetilde{g}^{t}(z)$ and $f^{t}(z)$ are the same. Since $B$ has the cylindrical structure, we obtain that $\widetilde{g}^{t}(z) \in B \times \tau\left(f^{t}(z)\right)=B \times\{t\}$. In particular, $\tilde{g}^{\frac{1}{2}}(B)=f^{\frac{1}{2}}(B)=B \times\left\{\frac{1}{2}\right\}$. Since $\mathcal{X}_{\widetilde{g}}=\mathcal{X}_{f}$ outside $\widetilde{\Omega}_{R}=B \times[0,1 / 2]$, we have that $\widetilde{g}^{1}\left(\Pi_{0}\right)=f^{1}\left(\Pi_{0}\right)=\Pi_{0}$. In other words, $\Pi_{0}$ is also a global cross-section for $\tilde{g}^{t} \mid \mathcal{N} \times U_{0}$.

It follows from the lemma that $\Pi_{0}$ is a global cross-section for the flow $\widetilde{g}^{t} \mid \mathcal{N} \times U_{0}$ and the time-1 map $\widetilde{g}^{1}$ restricted to $\Pi_{0}$ is the Poincaré return map of $\widetilde{g}^{t}$ on $\Pi_{0}$. Therefore, (4.8) is equivalent to

$$
\begin{equation*}
L_{4}(\widetilde{G})=0<L_{1}(\widetilde{G})<L_{2}(\widetilde{G})=L_{3}(\widetilde{G}) \tag{4.10}
\end{equation*}
$$

where $\widetilde{G}=\widetilde{G}_{\beta}=\widetilde{g}^{1} \mid \Pi_{0}$. In fact, by (4.7), we have that $L_{k}(\widetilde{G})=L_{k}\left(f^{1} \mid \Pi_{0}\right)$ for $k=$ $2,3,4$, and thus we only need to show that

$$
\begin{equation*}
L_{1}(\widetilde{G})<L_{1}\left(f^{1} \mid \Pi_{0}\right) . \tag{4.11}
\end{equation*}
$$

To this end we will use the following lemma.
Lemma 4.4. For all $z \in \Pi_{0}$ the derivative of $\widetilde{G}=\widetilde{g}^{1} \mid \Pi_{0}$ along $E_{f}^{u a}$ has the form

$$
d \widetilde{G}_{\beta}(z) \left\lvert\, E_{f}^{u a}(z)=\left(\begin{array}{cc}
\eta A(\beta, z) & \eta B(\beta, z)  \tag{4.12}\\
C(\beta, z) & D(\beta, z)
\end{array}\right)\right.,
$$

where

$$
\begin{gathered}
A=A(\beta, z)=1-\beta r \sigma_{r} \sin \theta \cos \theta-\frac{\beta^{2} \sigma^{2}}{2}-\beta^{2} r \sigma \sigma_{r} \cos ^{2} \theta+O\left(\beta^{3}\right) \\
B=B(\beta, z)=-\beta \sigma-\beta r \sigma_{r} \sin ^{2} \theta-\beta^{2} r \sigma \sigma_{r} \sin \theta \cos \theta+O\left(\beta^{3}\right) \\
C=C(\beta, z)=\beta \sigma+\beta r \sigma_{r} \cos ^{2} \theta-\beta^{2} r \sigma \sigma_{r} \sin \theta \cos \theta+O\left(\beta^{3}\right) \\
D=D(\beta, z)=1+\beta r \sigma_{r} \sin \theta \cos \theta-\frac{\beta^{2} \sigma^{2}}{2}-\beta^{2} r \sigma \sigma_{r} \sin ^{2} \theta+O\left(\beta^{3}\right)
\end{gathered}
$$

Proof. The desired relation (4.12) is apparent for $z \in \Pi_{0} \backslash B$ since $\widetilde{G}=f^{1}$ and $\sigma=0$ on $\Pi_{0} \backslash B$. Given $z=(r, \theta, 0, s, b) \in B$ in the $u a$-cylindrical coordinate, by Lemma 4.3, we have that $\tilde{g}^{t}(z)=(r, \theta+\theta(t), t, s, b)$, where $\theta(t)=\beta \int_{0}^{t} \widetilde{\psi}\left(\widetilde{g}^{\tau}(z)\right) d \tau$ for $0 \leq t \leq 1 / 2$. In particular, the coordinate of $\widetilde{g}^{\frac{1}{2}}(z)$ is $(r, \theta+\beta \sigma, 1 / 2, s, b)$, where

$$
\sigma=\sigma(r, s, b)=\frac{1}{4} \psi\left(\frac{r^{2}}{\varepsilon^{2}}\right) \psi\left(\frac{b}{\varepsilon}\right) \psi\left(\frac{s}{\varepsilon}\right) \int_{-1}^{1} \psi(u) d u .
$$

Back in the Cartesian coordinate system ( $u, a, \tau, s, b$ ), we obtain that

$$
\begin{aligned}
\tilde{g}^{\frac{1}{2}}(z) & =\left(u_{1}, a_{1}, 1 / 2, s, b\right) \\
& :=(u \cos (\beta \sigma)-a \sin (\beta \sigma), u \sin (\beta \sigma)+a \cos (\beta \sigma), 1 / 2, s, b),
\end{aligned}
$$

and hence

$$
\widetilde{G}(z)=\widetilde{g}^{1}(z)=f^{\frac{1}{2}} \widetilde{g}^{\frac{1}{2}}(z)=\left(u_{1}, a_{1}, 1, s, b\right)=\left(\eta u_{1}, a_{1}, 0, \eta^{-1} s, b\right) .
$$

The last equality follows from (A.1) and the fact that $\widetilde{g}^{1}\left(z_{0}\right)=f^{1}\left(z_{0}\right)$, where $z_{0}$ is the center of $B$. Since $\widetilde{G}$ preserves the $E_{f}^{u a}$ subbundle, we have that

$$
d \widetilde{G}_{\beta}(z) \left\lvert\, E_{f}^{u a}(z)=\left(\begin{array}{cc}
\eta A(\beta, z) & \eta B(\beta, z) \\
C(\beta, z) & D(\beta, z)
\end{array}\right)\right.,
$$

where

$$
A(\beta, z)=\frac{\partial u_{1}}{\partial u}, \quad B(\beta, z)=\frac{\partial u_{1}}{\partial a}, \quad C(\beta, z)=\frac{\partial a_{1}}{\partial u}, \quad D(\beta, z)=\frac{\partial a_{1}}{\partial a} .
$$

Then

$$
\begin{aligned}
A=\frac{\partial u_{1}}{\partial u} & =\cos (\beta \sigma)+[-u \sin (\beta \sigma)-a \cos (\beta \sigma)] \beta \sigma_{u} \\
& =\cos (\beta \sigma)-\beta r \sin (\theta+\beta \sigma) \sigma_{r} \cos \theta \\
& =1-\beta r \sigma_{r} \sin \theta \cos \theta-\frac{\beta^{2} \sigma^{2}}{2}-\beta^{2} r \sigma \sigma_{r} \cos ^{2} \theta+O\left(\beta^{3}\right) .
\end{aligned}
$$

Similarly, we can obtain the formulae for $B, C$ and $D$.
Lemma 4.4 allows us to follow the line of argument in [10, proof of Lemma 4.1] to establish (4.11). For the reader's convenience we outline the argument here.

Denote by $e_{\beta}(z)$ the unique number such that the vector $v_{\beta}(z)=\left(1, e_{\beta}(z)\right)^{t} \in E_{\widetilde{G}_{\beta}}^{u}(z)$ for all $z \in \Pi_{0}$. One can show that

$$
L_{\beta}=L_{1}\left(\widetilde{G}_{\beta}\right)=\int_{\Pi_{0}} \log \eta d m(z)-\int_{\Pi_{0}} \log \left[D(\beta, z)-\eta B(\beta, z) e_{\beta}\left(\widetilde{G}_{\beta}(z)\right)\right] d m(z)
$$

Note that $L_{0}=L_{1}\left(f^{1} \mid \Pi_{0}\right)$; we will show that

$$
\begin{equation*}
\left.\frac{d L_{\beta}}{d \beta}\right|_{\beta=0}=0,\left.\quad \frac{d^{2} L_{\beta}}{d \beta^{2}}\right|_{\beta=0}<0, \tag{4.13}
\end{equation*}
$$

which immediately implies that (4.11) holds for all sufficiently small $\beta>0$.
To show (4.13) observe that

$$
\left.\frac{d L_{\beta}}{d \beta}\right|_{\beta=0}=-\left.\int_{\Pi_{0}} D_{\beta}\right|_{\beta=0} d m(z)=0
$$

thus proving the first relation in (4.13). To prove the second relation, note that

$$
\left.\frac{d^{2} L_{\beta}}{d \beta^{2}}\right|_{\beta=0}=\int_{\Pi_{0}}\left[\left(D_{\beta}\right)^{2}-D_{\beta \beta}+2 \eta B_{\beta} \frac{\partial}{\partial \beta}\left(e_{\beta}\left(\widetilde{G}_{\beta}(z)\right)\right)\right]_{\beta=0} d m(z) .
$$

This integral can be written as

$$
\begin{align*}
& \int_{\Pi_{0}}\left[\left(D_{\beta}(0, z)\right)^{2}-D_{\beta \beta}(0, z)+2 \eta B_{\beta}(0, z) C_{\beta}(0, z)\right] d m(z) \\
& \quad+\int_{\Pi_{0}} \sum_{i=1}^{\infty} \frac{1}{\eta^{i}} 2 B_{\beta}(0, z) C_{\beta}\left(0, f^{-i}(z)\right) d m(z) \tag{4.14}
\end{align*}
$$

The first term in (4.14) is bounded from above by

$$
-(1-\varepsilon) \int_{\Pi_{0}} \sigma^{2} d m(z)-\frac{1}{8} \int_{\Pi_{0}} r^{2} \sigma_{r}^{2} d m(z)
$$

To estimate the second term in (4.14) note that

$$
\int_{\Pi_{0}} 2 B_{\beta}(0, z) C_{\beta}\left(0, F^{-i}(z)\right) d m(z) \leq 4 \int_{\Pi_{0}}\left(\sigma^{2}+r^{2} \sigma_{r}^{2}\right) d m(z)
$$

and that $B_{\beta}(0, z) C_{\beta}\left(0, f^{-i}(z)\right)=0$ for all $z \in \Pi_{0} \backslash B$ and all i. Moreover, $B_{\beta}(0, z) C_{\beta}\left(0, f^{-i}(z)\right)=0$ for every $z \in B$ and $i=1, \ldots, N_{0}-1$ since $f^{i}(B) \cap B=\emptyset$. This allows us to take $N_{0}>0$ large enough to ensure that the second term is bounded by

$$
\frac{1}{10} \int_{B}\left(\sigma^{2}+r^{2} \sigma_{r}^{2}\right) d m(z)
$$

Hence,

$$
\left.\frac{d^{2} L_{\beta}}{d \beta^{2}}\right|_{\beta=0} \leq-\left(\frac{9}{10}-\varepsilon\right) \int_{\Pi_{0}} \sigma^{2} d m(z)-\frac{1}{40} \int_{\Pi_{0}} r^{2} \sigma_{r}^{2} d m(z)<0 .
$$

This completes the proof of the inequality (4.11) thus guaranteeing that for any sufficiently small $\lambda>0$ the level set

$$
\Pi=\left\{z \in \Pi_{0}: \lambda_{1}(z, \widetilde{G}) \geq \lambda_{2}(z, \widetilde{G})>\lambda\right\}
$$

has positive volume. It is also invariant under $\widetilde{G}$. Since $f^{i}(B) \cap B=\emptyset$ for $i=1, \ldots, N_{0}$, we obtain that the sets $\widetilde{g}^{i}(\Pi \cap B)=\Pi \cap \widetilde{g}^{i}(B)=\Pi \cap f^{i}(B)$ corresponding to different $i$ are pairwise disjoint subsets of $\Pi$. This implies that

$$
m(\Pi) \geq N_{0} m(\Pi \cap B) \geq 20 k_{0} m(\Pi \cap B)>0,
$$

thus completing the proof of Proposition 4.2.
4.2. Construction of the flow $g^{t}$. We perturb the flow $\widetilde{g}^{t}$ to a flow $g^{t}$ by adding a vector field $\mathcal{X}_{R}$ to the vector field $\mathcal{X}_{\widetilde{g}}$. We obtain $\mathcal{X}_{R}$ as a sum of rotational vector fields in the $a b$-direction along several pairwise disjoint tubes so that the total rotation along an orbit that passes through these tubes is $\pi / 2$. This ensures positive Lyapunov exponents along the $E_{f}^{u a b}$ subbundle for the flow $g^{t}$.

Note that there is $M_{0}>0$ such that for any flow $F^{t}$ that is sufficiently $C^{1}$-close to the flow $f^{t}$

$$
\begin{equation*}
\left\|F^{1}-f^{1}\right\|_{C^{1}} \leq M_{0}\left\|\mathcal{X}_{F}-\mathcal{X}_{f}\right\|_{C^{1}} \tag{4.15}
\end{equation*}
$$

According to Lemma B.5, $M_{0}$ depends only on the Riemannian metric and the start-up flow $f^{t}$.

Let the number $\lambda>0$ and the subset $\Pi \subset \Pi_{0}$ be as in statement (5) of Proposition 4.2. Given $K>0$, let

$$
\begin{align*}
\Lambda^{\prime}=\Lambda^{\prime}(K)= & \left\{z \in \Pi:\left|\frac{1}{k} \log \left\|d \widetilde{g}^{k}(z, v)\right\|-\lambda\right| \leq 0.1 \lambda\right.  \tag{4.16}\\
& \text { for all } \left.v \in E_{f}^{u a}(z),\|v\|=1 \text { and all }|k| \geq 0.5 K\right\}
\end{align*}
$$

and

$$
\begin{equation*}
\Lambda=\Lambda(K)=\bigcap_{i=0}^{k_{0}-1} \tilde{g}^{-i} \Lambda^{\prime}(K) \tag{4.17}
\end{equation*}
$$

where $k_{0}>0$ is a number which will be defined later (see Lemma 4.5). Since $m\left(\Lambda^{\prime}(K)\right) \rightarrow$ $m(\Pi)$ as $K \rightarrow \infty$ and hence $m(\Lambda(K)) \rightarrow m(\Pi)$ as $K \rightarrow \infty$, one can choose $K$ so large that

$$
\begin{gather*}
K \lambda \geq \max \left\{5 k_{0} \lambda, 10 \log 2,-10 k_{0} \log \left(1-M_{0} \delta_{g}\right)\right\},  \tag{4.18}\\
\lambda m(\Pi)+40 \log \left(1-M_{0} \delta_{g}\right) m(\Pi \backslash \Lambda)>0,  \tag{4.19}\\
20 m(\Pi \backslash \Lambda) \leq m(\Pi) . \tag{4.20}
\end{gather*}
$$

Set

$$
\begin{equation*}
\Lambda^{*}=\Lambda \backslash \bigcup_{i=0}^{k_{0}-1} \tilde{g}^{-i}\left(\operatorname{Proj}_{\Pi_{0}}\left(\Omega_{0} \cup \widetilde{\Omega}_{R}\right)\right) \tag{4.21}
\end{equation*}
$$

where $\Omega_{0}$ and $\widetilde{\Omega}_{R}$ are given by (4.4) and (4.5) respectively. If the number $v$ is chosen small enough, statement (5) of Proposition 4.2 allows us to assume that

$$
\begin{gather*}
m\left(\operatorname{Proj}_{\Pi_{0}} \Omega_{0} \cap \Pi\right) \leq m(\Pi) / 20 k_{0}  \tag{4.22}\\
m\left(\operatorname{Proj}_{\Pi_{0}} \widetilde{\Omega}_{R} \cap \Pi\right)=m(B \cap \Pi) \leq m(\Pi) / 20 k_{0}
\end{gather*}
$$

Combining (4.20)-(4.22), we find that

$$
\begin{equation*}
m\left(\Lambda^{*}\right) \geq 0.8 m(\Pi) \tag{4.23}
\end{equation*}
$$

By statement (5) of Proposition 4.2, the set $\Pi$ is invariant under the time-1 map of the flow $\tilde{g}^{t}$.

We will approximate the set $\Pi$ by constructing an appropriate Rokhlin-Halmos tower (see [11]) for the map $\widetilde{g}^{1}$. More precisely, we choose a measurable subset $\Gamma^{\prime} \subset \Pi$ such that the sets $\widetilde{g}^{i}\left(\Gamma^{\prime}\right)$ are pairwise disjoint for $-K \leq i \leq 6 K+k_{0}-1$ and

$$
\begin{equation*}
m\left(\bigcup_{i=-K}^{6 K+k_{0}-1} \widetilde{g}^{i} \Gamma^{\prime}\right) \geq 0.9 m(\Pi) \tag{4.24}
\end{equation*}
$$

Consider the set $\Gamma_{0}$ of first entries of orbits $\left\{\widetilde{g}^{i}(z)\right\}_{i=0}^{5 K-1}$ (with $z \in \Gamma^{\prime}$ ) to the set $\Lambda^{*}$. More precisely, set

$$
\Gamma_{0}=\left\{\widetilde{g}^{j}(z): z \in \Gamma^{\prime}, 0 \leq j \leq 5 K-1, \widetilde{g}^{j}(z) \in \Lambda^{*}, \widetilde{g}^{i}(z) \notin \Lambda^{*} \text { for } i<j\right\}
$$

and let

$$
\begin{equation*}
\Gamma_{i}=\widetilde{g}^{i}\left(\Gamma_{0}\right), \quad \Gamma=\bigcup_{i=-K}^{K+k_{0}-1} \Gamma_{i} . \tag{4.25}
\end{equation*}
$$

Clearly, the sets $\left\{\Gamma_{i}\right\}$ are pairwise disjoint for $-K \leq i \leq K+k_{0}-1$. We then approximate $\Gamma_{0}$ by finitely many disjoint $a b$-cylinders $B_{0 j}$ of the form

$$
\begin{aligned}
B_{0 j} & =B^{u}\left(z_{j}, r_{j}^{\prime}\right) \times B^{s}\left(z_{j}, r_{j}^{\prime \prime}\right) \times B^{a b}\left(z_{j}, r_{j}\right) \\
& =\left\{\left(u_{j}, s_{j}, a_{j}, b_{j}\right):\left|u_{j}\right| \leq r_{j}^{\prime},\left|s_{j}\right| \leq r_{j}^{\prime \prime}, a_{j}^{2}+b_{j}^{2} \leq r_{j}^{2}\right\} \\
& =\left\{\left(u_{j}, s_{j}, \rho_{j}, \varphi_{j}\right):\left|u_{j}\right| \leq r_{j}^{\prime},\left|s_{j}\right| \leq r_{j}^{\prime \prime}, \rho_{j} \leq r_{j}\right\},
\end{aligned}
$$

where $r_{j}^{\prime}, r_{j}^{\prime \prime}, r_{j}>0$ for $j=1, \ldots, J$ and $z_{j}=\left(u_{j}, s_{j}, a_{j}, b_{j}\right)=\left(u_{j}, s_{j}, \rho_{j}, \varphi_{j}\right) \in \Pi_{0}$ is the center of $B_{0 j}$. For $i=-K, \ldots, K+k_{0}-1$ set

$$
\begin{equation*}
B_{i j}=\widetilde{g}^{i}\left(B_{0 j}\right), \quad \Delta_{i}=\bigcup_{j=1}^{J} B_{i j} . \tag{4.26}
\end{equation*}
$$

We can choose the sets $B_{0 j}$ in such a way that:
(1) $B_{i j} \cap B_{k l}=\emptyset$ for $(i, j) \neq(k, l)$ with $-K \leq i, k \leq K+k_{0}-1$ and $1 \leq j, l \leq J$;
(2) for each $i=0,1, \ldots, k_{0}$

$$
\begin{equation*}
m\left(\Gamma_{i} \Delta \Delta_{i}\right) \leq 0.05 \max \left\{m\left(\Gamma_{i}\right), m\left(\Delta_{i}\right)\right\} ; \tag{4.27}
\end{equation*}
$$

(3) $B_{i j} \cap \operatorname{Proj}_{\Pi_{0}}\left(\Omega_{0} \cup \widetilde{\Omega}_{R}\right)=\emptyset$ for $0 \leq i \leq k_{0}-1,1 \leq j \leq J$.

The last property implies that the set $B_{i j}=\widetilde{g}^{i}\left(B_{0 j}\right)=f^{i}\left(B_{0 j}\right)$ lies in a neighborhood around $f^{i}\left(z_{j}\right)$ and hence is still an $a b$-cylinder if the numbers $r_{j}, r_{j}^{\prime}, r_{j}^{\prime \prime}$ are chosen small enough.

We need the following lemma.
Lemma 4.5. Given $\delta>0$, there is $\theta_{0}=\theta_{0}(\delta)>0$ such that for any $\theta \in\left[0, \theta_{0}\right]$ and any tube $T=C \times[0,1 / 2]$, where $C \subset \Pi_{0}$ is an ab-cylinder of the form

$$
C=B^{u}\left(z, r^{\prime}\right) \times B^{s}\left(z, r^{\prime \prime}\right) \times B^{a b}(z, r),
$$

there exist a subtube $T^{\prime}=C^{\prime} \times[1 / 40,19 / 40] \subset T$, where $C^{\prime} \subset C$ is a cylinder of the form

$$
C^{\prime}=B^{u}\left(z, r_{0}^{\prime}\right) \times B^{s}\left(z, r_{0}^{\prime \prime}\right) \times B^{a b}\left(z, r_{0}\right),
$$

and a $C^{\infty}$ vector field $\mathcal{X}=\mathcal{X}_{T, \theta}$ on $\mathcal{M}$ such that the following hold.
(1) $\mathcal{X}$ is a rotation vector field with speed $\theta$ in the ab-plane; i.e.,

$$
\mathcal{X}(z)=\mathcal{X}(u, s, a, b, \tau)=\theta\left(-b \frac{\partial}{\partial a}+a \frac{\partial}{\partial b}\right), \quad z \in T^{\prime} .
$$

(2) $\mathcal{X}=0$ outside $T$.
(3) $m\left(C^{\prime}\right) / m(C) \geq 0.75$.
(4) $r_{0} / r, r_{0}^{\prime} / r^{\prime}, r_{0}^{\prime \prime} / r^{\prime \prime} \geq 0.9$.
(5) $\|\mathcal{X}\|_{C^{1}}<\delta$.

Moreover, let $k_{0}>0$ be such that

$$
\begin{equation*}
\bar{\theta}:=\frac{2 \pi}{k_{0} \int_{-1}^{1} \psi(t) d t}<\theta_{0}\left(\delta_{g} / 2\right), \tag{4.28}
\end{equation*}
$$

where $\psi(t)$ is the function in $\S 4.1$ and $\delta_{g}$ is given by Proposition 4.1. Then $\left\|\mathcal{X}_{T, \bar{\theta}}\right\| \leq \delta_{g} / 2$.
Proof. Given $0<\alpha<1$, we define a subcylinder

$$
C_{\alpha}=B^{u}\left(z, \alpha r^{\prime}\right) \times B^{s}\left(z, \alpha r^{\prime \prime}\right) \times B^{a b}(z, \alpha r) \subset C
$$

By (A.2), the volume of $C_{\alpha}$ and $C$ is induced by the flat metric $d u^{2}+d s^{2}+d a^{2}+d b^{2}$, and hence the ratio $m\left(C_{\alpha}\right) / m(C)$ depends only on $\alpha$ but not on the cylinder $C$. It follows that $m\left(C_{\alpha}\right) / m(C) \rightarrow 1$ as $\alpha \rightarrow 1$.

Fix $\alpha>0.9$ such that $m\left(C_{\alpha}\right) / m(C)>0.75$, and set $C^{\prime}=C_{\alpha}$. Let us choose a $C^{\infty}$ function $\xi: \mathbb{R} \rightarrow[0,1]$ satisfying:
(1) $\xi=1$ on $(-\alpha, \alpha)$;
(2) $\xi>0$ on $(-1,1)$ and $\xi=0$ outside $(-1,1)$;
(3) $\|\xi\|_{C^{1}} \leq 2 /(1-\alpha)$.

We introduce the $a b$-cylindrical coordinate ( $u, s, \rho, \varphi$ ), and define a $C^{\infty}$ rotational vector field $\mathcal{X}=\mathcal{X}_{T, \theta}$ by the formula

$$
\mathcal{X}_{T, \theta}(z)= \begin{cases}\theta \tilde{\xi}(z) \frac{\partial}{\partial \varphi}, & z \in T  \tag{4.29}\\ 0, & z \in \mathcal{M} \backslash T\end{cases}
$$

where

$$
\widetilde{\xi}(z)=\tilde{\xi}(u, s, \rho, \xi, \tau)=\xi\left(\frac{u}{r^{\prime}}\right) \xi\left(\frac{s}{r^{\prime \prime}}\right) \xi\left(\frac{\rho}{r}\right) \psi\left(\frac{\tau-1 / 4}{1 / 4}\right)
$$

and $\psi$ is the smooth function in $\S 4.1$. Note that $\|\tilde{\xi}(\partial / \partial \varphi)\| \leq c$ where $c>0$ depends only on $\alpha$ but not on the choice of the cylinder $C$. Thus for any $\delta>0$, there is $\theta_{0}=\theta_{0}(\delta)>0$ such that $\|\mathcal{X}\|_{C^{1}}<\delta$ for any $\theta \in\left[0, \theta_{0}\right]$.

Consider the tubes $T_{i j}=B_{i j} \times[0,1 / 2]$. Applying Lemma 4.5 with $T=T_{i j}$, we obtain a vector field $\mathcal{X}_{R, i j}=\mathcal{X}_{T_{i j}, \bar{\theta}}$ such that $\left\|\mathcal{X}_{R, i j}\right\| \leq \delta_{g} / 2$, where $\bar{\theta}$ is given by (4.28). Moreover, there is a sub-cylinder $B_{i j}^{\prime} \subset B_{i j}$ such that $m\left(B_{i j}^{\prime}\right) / m\left(B_{i j}\right) \geq 0.75$. Furthermore, by Lemma 4.5, we may assume that $\widetilde{g}^{i}\left(B_{0 j}^{\prime}\right)=B_{i j}^{\prime}$ for $i=1, \ldots, k_{0}$. Finally, let

$$
\begin{equation*}
\Delta_{i}^{\prime}=\bigcup_{j=1}^{J} B_{i j}^{\prime}, \quad \Omega_{R}=\bigcup_{i=0}^{k_{0}-1} \bigcup_{j=1}^{J} T_{i j} \tag{4.30}
\end{equation*}
$$

and define the vector field $\mathcal{X}_{R}$ by

$$
\begin{equation*}
\mathcal{X}_{R}=\sum_{i=0}^{k_{0}-1} \sum_{j=1}^{J} \mathcal{X}_{R, i j} \tag{4.31}
\end{equation*}
$$

We obtain a new flow $g^{t}$ generated by the vector field $\mathcal{X}_{g}=\mathcal{X}_{\widetilde{g}}+\mathcal{X}_{R}$. Clearly, $g^{t}$ is a $C^{\infty}$ volume preserving flow since $\mathcal{X}_{R}$ is divergence free. We will show that the flow $g^{t}$ has all the desired properties stated in Proposition 4.1.

Proof of Proposition 4.1. Statements (1) and (2) follow immediately from the construction of the flow $g^{t}$, and statement (3) can be proved in the same way as statement (3) of Proposition 4.2.

We will prove statement (4). We need the following statement whose proof is very similar to the proof of Lemma 4.3.

Lemma 4.6. Given $z \in B_{i j}$, the $\rho_{j^{-}}, u_{j}$-, $s_{j}$ - and $\tau$-coordinates of $g^{t}(z)$ and $f^{t}(z)$ are the same for $t \in[0,1 / 2]$. Consequently, $g^{\frac{1}{2}}\left(B_{i j}\right)=f^{\frac{1}{2}}\left(B_{i j}\right)$, and hence $\Pi_{0}$ is a global cross-section for the flow $g^{t} \mid \mathcal{N} \times U_{0}$.

By the lemma, the time- 1 map $g^{1}$ restricted to $\Pi_{0}$ is the Poincaré return map of $g^{t}$ on $\Pi_{0}$. Therefore, (4.2) is equivalent to

$$
\begin{equation*}
L_{4}(G)=0<L_{1}(G)<L_{2}(G)<L_{3}(G), \tag{4.32}
\end{equation*}
$$

where $G=g^{1} \mid \Pi_{0}$. In fact, by (4.1) and (4.7) we have for $k=3,4$ that

$$
L_{k}(G)=L_{k}\left(f^{1} \mid \Pi_{0}\right)=L_{k}\left(\widetilde{g}^{1} \mid \Pi_{0}\right)=L_{k}(\widetilde{G})
$$

Hence, we only need to show that

$$
\begin{equation*}
L_{2}(G)<L_{3}(G) . \tag{4.33}
\end{equation*}
$$

We follow the argument in [10, §4.2], and give a sketch of the proof of (4.33).
Set $\Delta_{0}^{*}=\Delta_{0}^{\prime} \cap \Lambda$, where $\Delta_{0}^{\prime}$ and $\Lambda$ are given by (4.30) and (4.17) respectively, and

$$
\begin{gathered}
U_{1}=G^{-K} \Delta_{0}^{*}, \quad U_{2}=\Delta_{0} \backslash \Delta_{0}^{*}, \\
U_{3}=G^{k_{0}}\left(\left(\Delta_{0} \cap \Lambda\right) \backslash \Delta_{0}^{*}\right), \quad U_{4}=G^{k_{0}}\left(\Delta_{0} \backslash \Lambda\right) .
\end{gathered}
$$

Consider the first return map $\bar{G}=G^{R}$ on the set

$$
U=U_{1} \cup U_{2} \cup U_{3} \cup U_{4} \subset \Pi_{0},
$$

where $R=R(z)$ is the first return time of $z \in U$ to $U$ under $G$. Note that the flow $g^{t}$ preserves the $E_{f}^{u a b}$-subbundle, and so does $G$.

We intend to show that

$$
\begin{equation*}
\int_{U}\left(\log \left\|\wedge^{3}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\|-\log \left\|\wedge^{2}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\|\right) d m(z)>0 \tag{4.34}
\end{equation*}
$$

where

$$
\wedge^{k}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right): \wedge^{k}\left(E_{f}^{u a b}(z)\right) \rightarrow \wedge^{k}\left(E_{f}^{u a b}(z)\right)
$$

is the $k$ th exterior power of $d \bar{G} \mid E_{f}^{u a b}(z)$. Indeed, assuming that (4.34) holds, consider the $G$-invariant set

$$
\Pi^{\prime}=\bigcup_{i=-\infty}^{\infty} G^{i}(U) \subset \Pi_{0}
$$

For $k=2$, 3 we have that

$$
\begin{aligned}
\int_{U} \log \left\|\wedge^{k}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\| d m(z) & =\int_{\Pi^{\prime}} \log \left\|\wedge^{k}\left(d G \mid E_{f}^{u a b}(z)\right)\right\| d m(z) \\
& =\int_{\Pi^{\prime}} \sum_{i=1}^{k} \lambda_{i}(z, G) d m(z)=L_{k}\left(G \mid \Pi^{\prime}\right)
\end{aligned}
$$

and hence (4.34) implies that $L_{2}\left(G \mid \Pi^{\prime}\right)<L_{3}\left(G \mid \Pi^{\prime}\right)$. Since $G=\widetilde{G}$ outside $\Pi^{\prime}$, we obtain that $L_{2}(G)<L_{3}(G)$.

To show (4.34) we split the integral over $U$ into four integrals $I_{1}, I_{2}, I_{3}$ and $I_{4}$ over the domains $U_{1}, U_{2}, U_{3}$ and $U_{4}$ respectively, and we obtain lower bounds for each of them. Namely, we will show that

$$
\begin{gather*}
I_{1} \geq 0.85 K \lambda \cdot 0.7 m\left(\Delta_{0}\right), \quad I_{2} \geq k_{0} \log \left(1-M_{0} \delta_{g}\right) \cdot 0.25 m\left(\Delta_{0}\right),  \tag{4.35}\\
I_{3} \geq 0, \quad I_{4} \geq 2 \log \left(1-M_{0} \delta_{g}\right) m(\Pi \backslash \Lambda) .
\end{gather*}
$$

The lower bounds for $I_{2}, I_{3}$ and $I_{4}$ can be obtained using arguments in $[\mathbf{1 0}$, proof of Lemma 4.2]. However, the proof of the lower bound for $I_{1}$ in our continuous-time case requires substantial changes and we shall present it here. We need the following lemma.
Lemma 4.7. Let $z \in U_{1}=G^{-K}\left(\Delta_{0}^{*}\right)$. Then for any $v \in E_{f}^{u a b}(z)$, we have

$$
\begin{equation*}
\left\|d_{z} \bar{G}(v)\right\| \geq \frac{\sqrt{2}}{2}\|v\| e^{0.9 K \lambda} \tag{4.36}
\end{equation*}
$$

Proof. Note that for any $z \in G^{-K}\left(\Delta_{0}^{*}\right)$, the first return time $R(z)$ is at least $2 K+k_{0}$. Set

$$
z_{1}=G^{K}(z), \quad z_{2}=G^{k_{0}}\left(z_{1}\right)=G^{K+k_{0}}(z), \quad z_{3}=\bar{G}(z)=G^{R(z)}(z) .
$$

Since the orbit segments $\left\{g^{t}(z)\right\}_{0 \leq t \leq K}$ from $z$ to $z_{1}$ and $\left\{g^{t}\left(z_{2}\right)\right\}_{0 \leq t \leq R(z)-K-k_{0}}$ from $z_{2}$ to $z_{3}$ are outside the set $\Omega_{R}$, we have that

$$
z_{1}=G^{K}(z)=\widetilde{G}^{K}(z), \quad z_{3}=G^{R(z)-K-k_{0}}\left(z_{2}\right)=\widetilde{G}^{R(z)-K-k_{0}}\left(z_{2}\right) .
$$

On the other hand, since $z_{1} \in \Delta_{0}^{*}=\Delta_{0}^{\prime} \cap \Lambda$, we can assume that $z \in B_{0 j}^{\prime}$ for some $j$, and, by our construction, we have $G^{i}\left(z_{1}\right) \in B_{i j}^{\prime}$ for $i=1, \ldots, k_{0}-1$, and every cylinder $B_{i j}^{\prime}$ is inside a local coordinate neighborhood of its center. Therefore, we write $z_{1}=$ $(u, s, a, b, 0) \in B_{0 j}^{\prime}$, and, applying similar arguments to those in the proof of Lemma 4.4, we have that

$$
\begin{aligned}
G\left(z_{1}\right)=f^{\frac{1}{2}} g^{\frac{1}{2}}\left(z_{1}\right) & =f^{\frac{1}{2}}(u, s, a \cos \phi-b \sin \phi, a \sin \phi+b \cos \phi, 1 / 2) \\
& =(u, s, a \cos \phi-b \sin \phi, a \sin \phi+b \cos \phi, 1) \\
& =\left(\eta u, \eta^{-1} s, a \cos \phi-b \sin \phi, a \sin \phi+b \cos \phi, 0\right)
\end{aligned}
$$

where

$$
\phi=\frac{1}{4} \bar{\theta} \int_{-1}^{1} \psi(t) d t=\frac{\pi}{2 k_{0}} .
$$

Repeating this calculation for $G^{1}\left(z_{1}\right), G^{2}\left(z_{1}\right), \ldots, G^{k_{0}-1}\left(z_{1}\right)$ and observing that $k_{0} \phi=$ $\pi / 2$, we obtain that

$$
\begin{aligned}
z_{2}=G^{k_{0}}\left(z_{1}\right) & =\left(\eta^{k_{0}} u, \eta^{-k_{0}} s, a \cos \left(k_{0} \phi\right)-b \sin \left(k_{0} \phi\right), a \sin \left(k_{0} \phi\right)+b \cos \left(k_{0} \phi\right), 0\right) \\
& =\left(\eta^{k_{0}} u, \eta^{-k_{0}} s,-b, a, 0\right)
\end{aligned}
$$

This formula means that $d_{z_{1}} G^{k_{0}}$ is non-contracting along the $E_{f}^{u a b}$ subbundle and rotates the vector in $E_{f}^{a b}$ by the angle $\pi / 2$.

To obtain (4.36), we write $v=v^{u a}+v^{b} \in E_{f}^{u a}(z) \oplus E_{f}^{b}(z)$ and consider the following two cases.
(1) If $\left\|v^{b}\right\| \leq(\sqrt{2} / 2)\|v\|$, since $d_{z} G^{K}=d_{z} \widetilde{G}^{K}$ and $z \in G^{-K} \Delta_{0}^{\prime} \subset \widetilde{G}^{-K} \Lambda^{\prime}$, by (4.16) and (4.17), we find that

$$
\left\|d_{z} G^{K} v\right\|=\left\|d_{z} \widetilde{G}^{K} v\right\| \geq\left\|d_{z} \widetilde{G}^{K} v^{u a}\right\| \geq\left\|v v^{u a}\right\| e^{0.9 K \lambda} \geq \frac{\sqrt{2}}{2}\|v\| e^{0.9 K \lambda},
$$

and hence

$$
\left\|d_{z} \bar{G} v\right\|=\left\|d_{z^{2}} \widetilde{G}^{R(z)-K-k_{0}} d_{z_{1}} G^{k_{0}} d_{z} \widetilde{G}^{K} v\right\| \geq\left\|d_{z} \widetilde{G}^{K} v\right\| \geq \frac{\sqrt{2}}{2}\|v\| e^{0.9 K \lambda} .
$$

(2) If $\left\|v^{b}\right\| \geq(\sqrt{2} / 2)\|v\|$, since $d_{z_{1}} G^{k_{0}}$ rotates the vector in $E_{f}^{a b}$ by the angle $\pi / 2$, we have

$$
d_{z} G^{K+k_{0}} v^{b}=d_{z_{1}} G^{k_{0}}\left(d_{z} \widetilde{G}^{K} v^{b}\right) \in E_{f}^{u a}\left(z_{2}\right) .
$$

Since $z_{2} \in \Lambda^{\prime}$, by (4.16) we obtain

$$
\begin{aligned}
\left\|d_{z} \bar{G} v\right\| \geq\left\|d_{z} \bar{G} v^{b}\right\| & =\left\|d_{z_{2}} \widetilde{G}^{R(z)-K-k_{0}} d_{z} G^{K+k_{0}} v^{b}\right\| \\
& \geq\left\|d_{z} G^{K+k_{0}} v^{b}\right\| e^{0.9 K \lambda} \geq \frac{\sqrt{2}}{2}\|v\| e^{0.9 K \lambda} .
\end{aligned}
$$

By the lemma and (4.18), we find that

$$
\log \left\|d_{z} \bar{G}(v)\right\| \geq 0.9 K \lambda-0.5 \log 2+\log \|v\| \geq 0.85 K \lambda+\log \|v\|,
$$

for any $z \in U_{1}$ and $v \in E_{f}^{u a b}(z)$. Hence,

$$
\log \left\|\wedge^{3}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\|-\log \left\|\wedge^{2}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\| \geq 0.85 K \lambda .
$$

On the other hand, it is proved in $[\mathbf{1 0}]$ that $m\left(U_{1}\right) \geq 0.7 m\left(\Delta_{0}\right)$. Therefore,

$$
\begin{aligned}
I_{1} & =\int_{U_{1}}\left(\log \left\|\wedge^{3}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\|-\log \left\|\wedge^{2}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\|\right) d m(z) \\
& \geq 0.85 K \lambda \cdot 0.7 m\left(\Delta_{0}\right)
\end{aligned}
$$

It follows from (4.35) that

$$
\begin{aligned}
& \int_{U}\left(\log \left\|\wedge^{3}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\|-\log \left\|\wedge^{2}\left(d \bar{G} \mid E_{f}^{u a b}(z)\right)\right\|\right) \\
& \quad \geq 0.595 \lambda K m\left(\Delta_{0}\right)+0.25 k_{0} \log \left(1-M_{0} \delta_{g}\right) m\left(\Delta_{0}\right)+2 \log \left(1-M_{0} \delta_{g}\right) m(\Pi \backslash \Lambda) \\
& \quad \geq 0.57 \lambda K m\left(\Delta_{0}\right)+2 \log \left(1-M_{0} \delta_{g}\right) m(\Pi \backslash \Lambda) \\
& \quad \geq 0.0627 \lambda m(\Pi)-0.05 \lambda m(\Pi)=0.0127 \lambda m(\Pi)>0 .
\end{aligned}
$$

The last two inequalities follow from (4.18), (4.19) and [10, Sublemma 4.4], which states that $m\left(\Delta_{0}\right) \geq 0.11 K^{-1} m(\Pi)$. This completes the proof of statement (4) of Proposition 4.1.

## 5. Accessibility

Notice that the flow $g^{t}$ has positive central exponents on a set of positive volume but is not necessarily ergodic. We will perturb $g^{t}$ to a flow $h^{t}$ that is pointwise partially hyperbolic on the open set $\mathcal{U}$ and still has positive central exponents. Furthermore, we will ensure that the flow $h^{t}$ possesses two transversal strongly stable and unstable foliations, $W_{h}^{s}$ and $W_{h}^{u}$, of $\mathcal{U}$ and satisfies the accessibility property on $\mathcal{U}$ via these two foliations. In view of Theorem 2.4, $h^{t}$ is indeed the desired flow in our main theorem.

We will follow the arguments in [10] and make some necessary modifications for the flow case. We choose two sequences of open subsets $U_{n}, \widetilde{U}_{n} \subset U, n=1,2, \ldots$ such that: (A4) $U_{0} \subset \widetilde{U}_{1}$;
(A5) $\widetilde{U}_{n} \subset \overline{\widetilde{U}_{n}} \subset U_{n} \subset \overline{U_{n}} \subset U$ and $\bigcup_{n \geq 1} U_{n}=U$;
(A6) $\tilde{U}_{n}$ and $U_{n}$ are connected sets for any $n \geq 1$.
Set

$$
\begin{equation*}
\mathcal{U}_{n}=\mathcal{N} \times U_{n}, \quad \tilde{\mathcal{U}}_{n}=\mathcal{N} \times \tilde{U}_{n} \tag{5.1}
\end{equation*}
$$

We will construct a sequence of flows $\left\{h_{n}^{t}\right\}_{n \geq 0}$, whose limit is the desired flow $h^{t}$. The goal of this section is to prove the following statement.

Proposition 5.1. Given $\delta_{h}>0$, one can find a sequence of positive numbers $\left\{\delta_{n}\right\}$ with $\delta_{n} \leq \min \left\{\delta_{h} / 2^{n}, d\left(C, U_{n}\right)^{2}\right\}$ as well as a sequence of $C^{\infty}$ divergence-free vector fields $\mathcal{X}_{n}$ on $\mathcal{M}$, generating a sequence of volume preserving flows $h_{n}^{t}$, such that for $n \geq 0$ the following hold.
(1) $\mathcal{X}_{0}=\mathcal{X}_{g}$, and hence $h_{0}^{t}=g^{t}$.
(2) $\left\|\mathcal{X}_{n+1}-\mathcal{X}_{n}\right\|_{C^{n+1}} \leq \delta_{n}$.
(3) $\mathcal{X}_{n}=\mathcal{X}_{f}$ on $\mathcal{M} \backslash \mathcal{U}_{n}$, and $\mathcal{X}_{n}=\mathcal{X}_{n-1}$ on $\mathcal{U}_{n-2}$; in particular, each flow $h_{n}^{t}$ is a gentle perturbation of $f^{t}$, and hence satisfies statements (3)-(5) of Proposition 3.1.
(4) For every $z \in \mathcal{M}$, we have

$$
E_{h_{n}}^{u a b \tau}(z)=E_{g}^{u a b \tau}(z), \quad \operatorname{det}\left(d h_{n}^{t} \mid E_{h_{n}}^{u a b \tau}(z)\right)=\operatorname{det}\left(d g^{t} \mid E_{g}^{u a b \tau}(z)\right)
$$

(5) For all $z \in \mathcal{U}_{j}, j=1, \ldots, n$, and $\omega=u, s, c$,

$$
\angle\left(E_{h_{n}}^{\omega}(z), E_{h_{n-1}}^{\omega}(z)\right) \leq \delta_{j} / 2^{n-j}
$$

(6) If the number $\delta_{g}$ in Proposition 4.1 is sufficiently small, then each flow $h_{n}^{t}$ is stably accessible in the following sense: let a flow $\tilde{h}^{t}$ be a gentle perturbation of the flow $f^{t}$, and assume that $L\left(E_{\tilde{h}}^{\omega}(z), E_{h_{n}}^{\omega}(z)\right) \leq \delta_{n}$ for all $z \in \mathcal{U}_{n}$ and $\omega=u, s, c$. Then any two points $z_{1}, z_{2} \in \tilde{\mathcal{U}}_{n}$ are accessible via a $(u, s)_{\tilde{h}^{t}}$-path in $\mathcal{U}$. In particular, $h_{n}^{t}$ has the accessibility property in $\tilde{\mathcal{U}}_{n}$.
Statements (1)-(3) imply that the limit vector field $\mathcal{X}_{h}=\lim _{n \rightarrow \infty} \mathcal{X}_{n}$ exists. Moreover,

$$
\left\|\mathcal{X}_{n}-\mathcal{X}_{k}\right\|_{C^{k+1}} \leq \sum_{j=k}^{n-1}\left\|\mathcal{X}_{j+1}-\mathcal{X}_{j}\right\|_{C^{j+1}} \leq \sum_{j=k}^{n-1} \delta_{j} \leq \delta_{h} / 2^{k-1}
$$

for any $n \geq k \geq 0$. It follows that $\mathcal{X}_{n}$ converges to $\mathcal{X}_{h}$ uniformly in the $C^{k+1}$ topology. Since $k$ is arbitrary, $\mathcal{X}_{h}$ is a $C^{\infty}$ vector field. In the following section we will show that the flow $h^{t}$ generated by $\mathcal{X}_{h}$ has all the desired properties.
5.1. Construction of the sets $U_{n}$ and $\widetilde{U}_{n}$. We view the 2 -torus $Y$ as the square $[0,8] \times$ $[0,8]$ whose opposite sides are identified. For each $n \geq 1$, consider the partition of $Y$ into squares

$$
\widehat{Z}_{i j}^{(n)}=\left[\frac{i}{2^{n}}, \frac{i+1}{2^{n}}\right] \times\left[\frac{j}{2^{n}}, \frac{j+1}{2^{n}}\right], \quad i, j=0,1, \ldots, 2^{n+3}-1 .
$$

Without loss of generality we will assume that the square $U_{0}$, constructed in $\S 3$, is contained in some $\widehat{Z}_{i_{0} j_{0}}^{(1)}$ so that

$$
d\left(U_{0}, \widehat{Z}_{i_{0} j_{0}}^{(1)}\right) \geq 1 / 2^{4} \quad \text { and } \quad d\left(C, \widehat{Z}_{i_{0} j_{0}}^{(1)}\right)>2,
$$

where $C$ is the Cantor set constructed in $\S 3$. Consider the open squares

$$
\begin{aligned}
& Z_{i j}^{(n)}=\left(\frac{i}{2^{n}}-\frac{1}{2^{n+2}}, \frac{i+1}{2^{n}}+\frac{1}{2^{n+2}}\right) \times\left(\frac{j}{2^{n}}-\frac{1}{2^{n+2}}, \frac{j+1}{2^{n}}+\frac{1}{2^{n+2}}\right), \\
& \widetilde{Z}_{i j}^{(n)}=\left(\frac{i}{2^{n}}-\frac{1}{2^{n+5}}, \frac{i+1}{2^{n}}+\frac{1}{2^{n+5}}\right) \times\left(\frac{j}{2^{n}}-\frac{1}{2^{n+5}}, \frac{j+1}{2^{n}}+\frac{1}{2^{n+5}}\right) .
\end{aligned}
$$

Clearly, these squares have the same center as $\widehat{Z}_{i j}^{(n)}$ and $\widehat{Z}_{i j}^{(n)} \subset \widetilde{Z}_{i j}^{(n)} \subset Z_{i j}^{(n)}$. For $n \geq 1$ consider the set

$$
Y_{n}=\left\{y \in Y: d(y, C) \geq 1 / 2^{n-2}\right\} .
$$

Since $U_{0} \subset Y_{1}$, we let $Y_{n}^{\prime}$ be the connected component of $Y_{n}$ that contains $U_{0}$. Finally, consider the sets

$$
\widehat{U}_{1}=\widehat{Z}_{i_{0} j_{0}}^{(1)}, \quad U_{1}=Z_{i_{0} j_{0}}^{(1)} \quad \text { and } \quad \widetilde{U}_{1}=\widetilde{Z}_{i_{0} j_{0}}^{(1)}
$$

and, for $n>1$,

$$
\widehat{U}_{n}=\bigcup_{\widehat{Z}_{i j}^{(n)} \cap Y_{n}^{\prime} \neq \emptyset} \widehat{Z}_{i j}^{(n)}, \quad U_{n}=\bigcup_{\widehat{Z}_{i j}^{(n)} \cap Y_{n}^{\prime} \neq \emptyset} Z_{i j}^{(n)}, \quad \widetilde{U}_{n}=\bigcup_{\widehat{Z}_{i j}^{(n)} \cap Y_{n}^{\prime} \neq \emptyset} \widetilde{Z}_{i j}^{(n)} .
$$

It is clear that the sets $U_{n}$ and $\tilde{U}_{n}$ satisfy conditions (A4)-(A6).
Let $\widehat{\mathcal{Z}}_{n}=\left\{\widehat{Z}_{i j}^{(n)}: \widehat{Z}_{i j}^{(n)} \subset \widehat{U}_{n} \backslash \widehat{U}_{n-1}\right\}$ and $\mathcal{Z}_{n}=\left\{Z_{i j}^{(n)}: \widehat{Z}_{i j}^{(n)} \subset \widehat{\mathcal{Z}}_{n}\right\}$. Relabeling elements of $\mathcal{Z}_{n}$, we will denote them by $Z_{1}^{(n)}, \ldots, Z_{k_{n}}^{(n)}$, and we will use the notations $\widehat{Z}_{l}^{(n)}$ and $\widetilde{Z}_{l}^{(n)}$ for the corresponding squares contained in $Z_{l}^{(n)}$. Thus we have

$$
U_{n}=U_{n-1} \cup\left(\bigcup_{l=1}^{k_{n}} Z_{l}^{(n)}\right)
$$

Clearly the collection of sets $\left\{\widehat{Z}_{l}^{(n)}: n=1,2, \ldots, l=1, \ldots, k_{n}\right\}$ forms a countable partition of $U$ up to a set of zero volume while the collection of sets $\left\{Z_{l}^{(n)}: n=\right.$ $\left.1,2, \ldots, l=1, \ldots, k_{n}\right\}$ forms a cover of $U$ of multiplicity at most 4. The following lemma is proved in [10].
LEMmA 5.2. There is a labeling of the squares $\left\{Z_{l}^{(n)}\right\}$ by integers from 1 to 8 such that for any $y \in U$, the labels of the squares $Z_{l}^{(n)}$ containing $y$ are all different. In particular, $Z_{1}^{(1)}$ can be labeled by 1.
5.2. Construction of the vector fields $\mathcal{X}_{n}$. The construction is similar to the one in $[\mathbf{1 0}$, §5.2], with a slightly modification on the choice of the collection of periodic points. We need the following preparations before we construct the vector fields $\mathcal{X}_{n}$.

Let $q_{j}, j=1, \ldots, 8$, be eight periodic points of the Anosov automorphism $A$ whose orbits are pairwise disjoint. There is $\epsilon_{0}>0$ such that

$$
B_{X}\left(A^{i} q_{j}, \epsilon_{0}\right) \cap B_{X}\left(A^{i} q_{j^{\prime}}, \epsilon_{0}\right)=\emptyset
$$

whenever $j \neq j^{\prime}$ and $i=-1,0,1$. For each $q_{j}$ we choose three periodic points $p_{j}^{i} \in$ $B_{X}\left(A^{i} q_{j}, \epsilon_{0} / 3\right)$ for $A, i=a, b, \tau$, whose orbits are pairwise disjoint. Write $\left[q_{j}, p_{j}^{i}\right]=$ $V_{A}^{u}\left(q_{j}\right) \cap V_{A}^{s}\left(p_{j}^{i}\right), i=a, b, \tau$, where $V_{A}^{s}$ and $V_{A}^{u}$ are the stable and unstable local manifolds respectively. For $i=a, b, \tau$ and $j=1, \ldots, 8$, consider the closed quadrilateral $(u, s)_{A}$-path $\gamma_{j}^{i}$ with the collection of points $q_{j},\left[q_{j}, p_{j}^{i}\right], p_{j}^{i},\left[p_{j}^{i}, q_{j}\right]$, and $q_{j}$. Without loss of generality, we will assume $q_{1}=q, p_{1}^{i}=p^{i}$ and $\gamma_{1}^{i}=\gamma^{i}$ for $i=a, b, \tau$, where $q$, $p^{i}$ and $\gamma^{i}$ are chosen as in the beginning of $\S 4.1$.

For $j=1, \ldots, 8$ and $i=a, b, \tau$, we have

$$
A^{\iota\left(q_{j}\right)}\left(q_{j}\right)=q_{j}, \quad A^{\iota\left(p_{j}^{i}\right)}\left(p_{j}^{i}\right)=p_{j}^{i},
$$

where $\iota\left(q_{j}\right)$ and $\iota\left(p_{j}^{i}\right)$ are periods of $q_{j}$ and $p_{j}^{i}$ respectively. There exists $\bar{\alpha}(j, i)>0$ such that for any $\alpha \in Y=\mathbb{T}^{2}$ with $\|\alpha\| \leq \bar{\alpha}(j, i)$, the Anosov affine map $A+\alpha$ has a $\iota\left(q_{j}\right)$ periodic point $q_{j}(\alpha)$ close to $q_{j}$ and a $\iota\left(p_{j}^{i}\right)$-periodic point $p_{j}^{i}(\alpha)$ close to $p_{j}^{i}$. Moreover, we can choose the number $\bar{\alpha}$ (in condition ( $\alpha 3$ ) at the beginning of $\S 3$ ) to be less than $\min \{\bar{\alpha}(j, i): j=1, \ldots, 8, i=a, b, \tau\}$ such that any two points from the set of periodic points

$$
\left\{q_{j}(\alpha), p_{j}^{i}(\alpha): j=1, \ldots, 8, i=a, b, \tau,\|\alpha\| \leq \bar{\alpha}\right\}
$$

are disjoint.
Given $n \geq 1$ and $l=1, \ldots, k_{n}$, let $j$ be the label of $Z_{l}^{(n)}$ in Lemma 5.2, and $y_{0}(n, l)=$ $\left(a_{0}(n, l), b_{0}(n, l)\right)$ the center of $Z_{l}^{(n)}$. We take the points associated to $Z_{l}^{(n)}$ as follows:

$$
\begin{equation*}
q(n, l)=q_{j}\left(\alpha\left(y_{0}(n, l)\right)\right), \quad p^{i}(n, l)=p_{j}^{i}\left(\alpha\left(y_{0}(n, l)\right)\right), \tag{5.2}
\end{equation*}
$$

where $i=a, b, \tau$. Recall that $\eta$ is the expanding rate of $A$ along its unstable direction, and the function $\kappa: Y \rightarrow \mathbb{R}$ is given in the beginning of $\S 3$. For $n \geq 1$ let us choose a square $Z_{l}^{(n)} \in \mathcal{Z}_{n}$. In the case $n>1$, we write for simplicity $q=q(n, l)$ and $p^{i}=p^{i}(n, l)$, and we let $\eta_{-}(n, l)=\min \left\{\eta^{\kappa(y)}: y \in Z_{l}^{(n)}\right\}$. Define the numbers

$$
\begin{align*}
\alpha_{u}^{i} & =\alpha_{u}^{i}(n, l) \\
\alpha_{s}^{i} & =d\left(p^{i},\left[p^{i}, q\right]\right),  \tag{5.3}\\
\breve{\alpha}_{u}^{i} & =\breve{\alpha}_{u}^{i}(n, l) \\
\breve{\alpha}_{s}^{i} & =\breve{\alpha}_{s}^{i}\left(n, l\left(p^{i},\left[q, p^{i}(n, l) / \eta_{-}(n, l),\right.\right.\right. \\
& =\alpha_{s}^{i}(n, l) / \eta_{-}(n, l)
\end{align*}
$$

and the rectangles in $X$

$$
\begin{aligned}
& \Pi^{i}(n, l)=B_{F^{u}}\left(p^{i}, \alpha_{u}^{i}\right) \times B_{F^{s}}\left(p^{i}, \alpha_{s}^{i}\right), \\
& \breve{\Pi}^{i}(n, l)=B_{F^{u}}\left(p^{i}, \breve{\alpha}_{u}^{i}\right) \times B_{F^{s}}\left(p^{i}, \breve{\alpha}_{s}^{i}\right) .
\end{aligned}
$$

We will assume that the rectangles $\Pi^{i}(n, l), n \geq 1, l=1, \ldots, k_{n}$ and $i=a, b$, $\tau$, are pairwise disjoint if the number $\bar{\alpha}$ is chosen sufficiently small. Finally, we let

$$
\begin{align*}
& \epsilon_{\tau}=\epsilon_{\tau}(n, l)=\min \left\{\kappa(y) / 2: y \in Z_{l}^{(n)}\right\},  \tag{5.4}\\
& \breve{\epsilon}_{\tau}=\breve{\epsilon}_{\tau}(n, l)=5 \epsilon_{\tau}(n, l) / 6 .
\end{align*}
$$

In the case $n=1$, we have $Z_{1}^{(1)}=U_{1}$ and $q(1,1)=q_{1}, p^{i}(1,1)=p_{1}^{i}$ since the function $\alpha=0$ on $U_{1}$. Choose $l_{u}^{i}$ and $l_{s}^{i}$ such that

$$
A^{-l_{u}^{i}}\left(\left[p_{1}^{i}, q_{1}\right]\right) \in B_{X}\left(p_{1}^{i}, \nu / 2\right), \quad A^{l_{s}^{i}}\left(\left[q_{1}, p_{1}^{i}\right]\right) \in B_{X}\left(p_{1}^{i}, \nu / 2\right),
$$

where $v$ is given in (4.4). Then we set

$$
\alpha_{u}^{i}=\alpha_{u}^{i}(1,1)=d\left(p_{1}^{i}, A^{-l_{u}^{i}}\left[p_{1}^{i}, q_{1}\right]\right), \quad \alpha_{s}^{i}=\alpha_{s}^{i}(1,1)=d\left(p_{1}^{i}, A_{s}^{l_{s}^{i}}\left[q_{1}, p_{1}^{i}\right]\right)
$$

with other quantities and sets to be defined in a similar way.
In addition to the squares $\widehat{Z}_{i j}^{(n)}, \widetilde{Z}_{i j}^{(n)}$ and $Z_{i j}^{(n)}$ constructed in the previous subsection, we need to consider the squares,

$$
\begin{aligned}
& \breve{Z}_{i j}^{(n)}=\left(\frac{i}{2^{n}}-\frac{1}{2^{n+3}}, \frac{i+1}{2^{n}}+\frac{1}{2^{n+3}}\right) \times\left(\frac{j}{2^{n}}-\frac{1}{2^{n+3}}, \frac{j+1}{2^{n}}+\frac{1}{2^{n+3}}\right), \\
& \bar{Z}_{i j}^{(n)}=\left(\frac{i}{2^{n}}-\frac{1}{2^{n+4}}, \frac{i+1}{2^{n}}+\frac{1}{2^{n+4}}\right) \times\left(\frac{j}{2^{n}}-\frac{1}{2^{n+4}}, \frac{j+1}{2^{n}}+\frac{1}{2^{n+4}}\right),
\end{aligned}
$$

as well as the intervals,

$$
I_{n}=J_{n}=\left(-\frac{3}{2^{n+2}}, \frac{3}{2^{n+2}}\right), \quad \breve{I}_{n}=\breve{J}_{n}=\left(-\frac{5}{2^{n+3}}, \frac{5}{2^{n+3}}\right),
$$

and

$$
K=(-1 / 4,1+1 / 4), \quad \breve{K}=(-1 / 8,1+1 / 8), \quad \bar{K}=(-1 / 16,1+1 / 16) .
$$

Note that we have

$$
\widehat{Z}_{i j}^{(n)} \subset \widetilde{Z}_{i j}^{(n)} \subset \bar{Z}_{i j}^{(n)} \subset \breve{Z}_{i j}^{(n)} \subset Z_{i j}^{(n)}
$$

and similar relations for $I_{n}$ and $J_{n}$.
Fix $n \geq 1$ and $l=1, \ldots, k_{n}$, and write $\alpha_{\omega}^{i}=\alpha_{\omega}^{i}(n, l), \breve{\alpha}_{\omega}^{i}=\breve{\alpha}_{\omega}^{i}(n, l)$ for $i=a, b, \tau$, $\omega=u, s$, and $\epsilon_{\tau}=\epsilon_{\tau}(n, l), \breve{\epsilon}_{\tau}=\breve{\epsilon}_{\tau}(n, l)$. We choose functions as follows.
(1) $\phi^{i}$ and $\psi^{i}$ are $C^{\infty}$ functions on $\mathbb{R}$ such that:

- $\quad \phi^{i}=$ constant on $\left(-\breve{\alpha}_{u}^{i}, \breve{\alpha}_{u}^{i}\right)$ and $\psi^{i}=$ constant on $\left(-\breve{\alpha}_{s}^{i}, \breve{\alpha}_{s}^{i}\right)$;
- $\quad \phi^{i}(r)=0$ for $|r| \geq \alpha_{u}^{i}, \psi^{i}(r)=0$ for $|r| \geq \alpha_{s}^{i}$;
- $\quad \int_{0}^{ \pm \alpha_{u}^{i}} \phi^{i}(r) d r=0$, and $\psi^{i}(r)>0$ for any $|r|<\alpha_{s}^{i}$;
- $\quad\left\|\phi^{i}\right\|_{C^{n}},\left\|\psi^{i}\right\|_{C^{n}} \leq 1$.
(2) $\xi_{\tau}$ and $\xi_{Y}$ are $C^{\infty}$ functions supported on $K$ and $I_{n}$ respectively such that:
- $\quad \xi_{\tau}=$ constant on $\breve{K}$, and $\xi_{Y}=$ constant on $\breve{I}_{n}$;
- $\quad \xi_{\tau}(r)>0$ for $r \in K$, and $\xi_{Y}(r)>0$ for $r \in I_{n}$;
- $\quad \xi_{\tau}(r)=0$ for $r \notin K$, and $\xi_{Y}(r)=0$ for $r \notin I_{n}$;
- $\quad\left\|\xi_{\tau}\right\|_{C^{n}},\left\|\xi_{Y}\right\|_{C^{n}} \leq 1$.
(3) $\zeta_{\tau}$ and $\zeta_{Y}$ are $C^{\infty}$ functions supported on $\left(-\epsilon_{\tau}, \epsilon_{\tau}\right)$ and $I_{n}$ respectively such that:
- $\quad \zeta_{\tau}=$ constant on $\left(-\breve{\epsilon}_{\tau}, \breve{\epsilon}_{\tau}\right)$, and $\zeta_{Y}=$ constant on $\breve{I}_{n}$;
- $\quad \zeta_{\tau}(r)>0$ for $r \in\left(-\epsilon_{\tau}, \epsilon_{\tau}\right)$, and $\zeta_{Y}(r)>0$ for $r \in I_{n}$;
- $\quad \zeta_{\tau}(r)=0$ for $r \notin\left(-\epsilon_{\tau}, \epsilon_{\tau}\right)$, and $\zeta_{Y}(r)=0$ for $r \notin I_{n}$;
- $\quad\left\|\zeta_{\tau}\right\|_{C^{n}},\left\|\zeta_{Y}\right\|_{C^{n}} \leq 1$.

Now we are ready to construct the sequence of vector fields $\mathcal{X}_{n}$. Given $n \geq 1$, $l=1, \ldots, k_{n}$ and $i=a, b, \tau$, take the Cartesian coordinate system $z=(u, s, \tau, a, b)=$ $(x, \tau, a, b)$ with the origin at $\left(p^{i}(n, l), 1 / 2, y_{0}(n, l)\right)$. In this coordinate system the interval $K$ is in the symmetric form $(-3 / 4,3 / 4)$. Take the boxes for $i=a, b$,

$$
\Omega^{i}=\Omega_{n, l}^{i}=\left\{(x, \tau, y): x \in \Pi^{i}(n, l),|\tau| \leq \epsilon_{\tau}(n, l), y \in Z_{l}^{(n)}\right\}
$$

and

$$
\Omega^{\tau}=\Omega_{n, l}^{\tau}=\left\{(x, \tau, y): x \in \Pi^{\tau}(n, l), \tau \in K, y \in Z_{l}^{(n)}\right\} .
$$

By the construction of the rectangles $\Pi^{i}(n, l)$, we have that $\Omega^{i}(n, l) \cap \Omega^{i^{\prime}}\left(n^{\prime}, l^{\prime}\right)=\emptyset$ if $(i, n, l) \neq\left(i^{\prime}, n^{\prime}, l^{\prime}\right)$. Similarly, we can choose $\breve{\Omega}^{i}, i=a, b, \tau$, by taking $\breve{\Pi}^{i}, \breve{\epsilon}_{\tau}, \breve{K}$ and $\breve{Z}_{l}^{(n)}$. Next we define three divergence-free vector field

$$
\begin{aligned}
& \mathcal{X}^{a}=\mathcal{X}_{n, l}^{a}=\zeta_{Y}(b) \zeta_{\tau}(\tau) \psi^{a}(s)\left(-\xi_{Y}^{\prime}(a) \int_{0}^{u} \phi^{a}(r) d r, 0,0, \xi_{Y}(a) \phi^{a}(u), 0\right), \\
& \mathcal{X}^{b}=\mathcal{X}_{n, l}^{b}=\zeta_{Y}(a) \zeta_{\tau}(\tau) \psi^{b}(s)\left(-\xi_{Y}^{\prime}(b) \int_{0}^{u} \phi^{b}(r) d r, 0,0,0, \xi_{Y}(b) \phi^{b}(u)\right), \\
& \mathcal{X}^{\tau}=\mathcal{X}_{n, l}^{\tau}=\zeta_{Y}(a) \zeta_{Y}(b) \psi^{\tau}(s)\left(-\xi_{\tau}^{\prime}(\tau) \int_{0}^{u} \phi^{\tau}(r) d r, 0, \xi_{\tau}(\tau) \phi^{\tau}(u), 0,0\right) .
\end{aligned}
$$

Clearly each vector field $\mathcal{X}_{n, l}^{i}$ vanishes outside the corresponding box $\Omega_{n, l}^{i}$, and it is constant on the smaller box $\breve{\Omega}_{n, l}^{i}$. Finally, we set

$$
\begin{equation*}
\widehat{\mathcal{X}}_{n}=\sum_{l=1}^{k_{n}}\left(\mathcal{X}_{n, l}^{a}+\mathcal{X}_{n, l}^{b}+\mathcal{X}_{n, l}^{\tau}\right), \quad \mathcal{X}_{n}=\mathcal{X}_{g}+\sum_{k=1}^{n} \beta_{k} \widehat{\mathcal{X}}_{k}, \tag{5.5}
\end{equation*}
$$

where the sequence of small positive numbers $\left\{\beta_{n}\right\}$ is determined inductively to ensure statements (2) and (5) of Proposition 5.1. Let $h_{n}^{t}$ be the flow on $\mathcal{M}$ generated by the vector field $\mathcal{X}_{n}$.
5.3. Proof of Proposition 5.1. Statements (1)-(4) follow directly from our construction. It remains to show how to choose the sequence of positive numbers $\delta_{n}$ such that $h_{n}^{t}$ satisfies statements (5) and (6) of the proposition. Note that these two statements only concern those invariant subbundles $E^{\omega}$ and foliations $W^{\omega}, \omega=u, s, c, c s, c u$, which are the same for the flow and its time-1 map. Therefore, the choice of $\delta_{n}$ and related arguments are similar to the diffeomorphism case in [10]. We will outline the proof here.

For any gentle perturbation $h_{\natural}^{t}$ of $f^{t}$ (see Definition 2.1), we denote by $W_{h_{\sharp}}^{c}(z)$ the center manifold of $h_{\square}^{t}$ at the point $z \in \mathcal{M}$. Given a square $Z_{l}^{(n)}$ with the center $y_{0}(n, l)$, let $q(n, l), p^{i}(n, l), i=a, b, \tau$, be the associated periodic points given by (5.2), and $z_{0}=z_{0}(n, l)=\left(q(n, l), 1 / 2, y_{0}(n, l)\right)$. We denote by $W_{h_{घ}}^{c}\left(z_{0}, K, Z_{l}^{(n)}\right)$ the connected
component of $W_{h_{\square}}^{c}\left(z_{0}\right) \cap\left(X \times K \times Z_{l}^{(n)}\right)$ that contains $z_{0}$. We will also use similar notations $W_{h_{घ}}^{c}\left(z_{0}, \breve{K}, \breve{Z}_{l}^{(n)}\right)$, etc.

Next we will introduce two important families of maps $\Theta$ and $\Psi$ for a gentle perturbation $h_{\natural}^{t}$ of $f^{t}$.

Fixing $n \geq 1$ and $l=1, \ldots, k_{n}$, we take the collection of points $q=q(n, l), p^{i}=$ $p^{i}(n, l), i=a, b, \tau$. Consider a quadrilateral $(u, s)_{h_{\natural}^{t}}$-path $\widehat{\gamma}^{i}=\left\{z_{1}, \ldots, z_{5}\right\}$ with initial point $z_{1}$ defined by

$$
\begin{align*}
& z_{2}=V_{h_{\natural}}^{u}\left(z_{1}\right) \cap V_{h_{\natural}}^{c s}\left(p^{i}, 1 / 2, a_{0}, b_{0}\right), \\
& z_{3}=V_{h_{\natural}}^{s}\left(z_{2}\right) \cap V_{h_{\natural}}^{c u}\left(p^{i}, 1 / 2, a_{0}, b_{0}\right), \\
& z_{4}=V_{h_{\natural}}^{u}\left(z_{3}\right) \cap V_{h_{\natural}}^{c s}\left(z_{1}\right),  \tag{5.6}\\
& z_{5}=V_{h_{\natural}}^{s}\left(z_{4}\right) \cap V_{h_{\natural}}^{c u}\left(z_{1}\right) .
\end{align*}
$$

This path defines a map $\Theta^{i}=\Theta_{n, l, h_{\sharp}}^{i}$ given by $\Theta^{i}\left(z_{1}\right)=z_{5}$. It is easy to see that $z_{5} \in V_{h_{\sharp}}^{c}\left(z_{1}\right)$, and $\Theta^{i}$ maps $W_{h_{\sharp}}^{c}\left(z_{0}, K, Z_{l}^{(n)}\right)$ into itself. Reparameterizing the curve on $V_{h_{\sharp}}^{u}\left(z_{1}\right)$ from $z_{1}$ to $z_{2}$ by $\sigma:[0,1] \rightarrow V_{h_{\sharp}}^{u}\left(z_{1}\right)$ so that $\sigma(0)=z_{1}$ and $\sigma(1)=z_{2}$, we obtain a parameterized family of quadrilaterals $\widehat{\gamma}^{i}(\vartheta)=\left\{z_{1}(\vartheta), \ldots, z_{5}(\vartheta)\right\}, \vartheta \in[0,1]$, where $z_{1}(\vartheta)=z_{1}, z_{2}(\vartheta)=\sigma(\vartheta)$, and $z_{k}(\vartheta), k=3,4,5$, are obtained in the way similar to (5.6). Then we define $\Theta_{\vartheta}^{i}=\Theta_{\vartheta, n, l, h_{\natural}}^{i}$ given by $\Theta_{\vartheta}^{i}\left(z_{1}\right)=z_{5}(\vartheta)$. Clearly $\Theta_{0}^{i}=\mathrm{Id}, \Theta_{1}^{i}=\Theta^{i}$, $\Theta_{\vartheta}^{i}$ maps $W_{h_{\sharp}}^{c}\left(z_{0}, K, Z_{l}^{(n)}\right)$ into $W_{h_{\sharp}}^{c}\left(z_{0}\right)$ and depends continuously on $\vartheta \in[0,1]$.

On the other hand, given $z=((x, \tau), y) \in \mathcal{U}$, there is a $(u, s)_{f^{t}}$-path $\gamma_{f}(z)$ connecting $z$ to $z^{\prime}=((q, \tau), y)$ whose length does not exceed $2 d(x, q)$. This generates a map $\Psi_{f}=\Psi_{f, n, l}$ from $\mathcal{U}$ to $\{q\} \times K \times G$ given by $\Psi_{f}(z)=z^{\prime}$. Furthermore, given a gentle perturbation $h_{\natural}^{t}$ of $f^{t}$ and a point $z \in Z_{l}^{(n)}$, we can find a $(u, s)_{h_{\sharp}}$-path $\gamma_{h_{\natural}}(z)$ that is close to $\gamma_{f}(z)$ and connects $z$ to a point $z^{\prime}=z^{\prime}\left(h_{\natural}^{t}\right) \in W_{h_{\natural}}^{c}\left(z_{0}, K, Z_{l}^{(n)}\right)$. We can then define $\Psi_{h_{\sharp}}=\Psi_{h_{\sharp}, n, l}$ by $\Psi_{h_{\sharp}}(z)=z^{\prime}\left(h_{\sharp}^{t}\right)$.

Note that the maps $\Psi_{h_{\natural}, n, l}, \Theta_{n, l, h_{\natural}}^{i}$ and $\Theta_{\vartheta, n, l, h_{\natural}}^{i}, i=a, b, \tau$, depend continuously on $h_{\natural}^{t}$ as long as $h_{\natural}^{t}$ is a gentle perturbation of $f^{t}$ with $h_{\natural}^{t}=f^{t}$ outside some fixed $\mathcal{U}_{n}$ and with $\angle\left(E_{h_{घ}}^{\omega}(z), E_{f}^{\omega}(z)\right)$ sufficiently small for all $z \in \mathcal{U}_{n}$ and $\omega=u, s, c$. Moreover, the continuity is uniform with respect to $z$.

Given a set $\Gamma \subset \mathcal{M}$ and a gentle perturbation $h_{\natural}^{t}$ of $f^{t}$, set

$$
\begin{aligned}
\mathcal{A}_{h_{\natural}}(\Gamma)= & \{z \in \mathcal{M} \text { : there exists } y \in \Gamma \text { such that } \\
& \left.y \text { is accessible to } z \text { via a }(u, s)_{h_{\natural}^{t}} \text {-path }\right\} .
\end{aligned}
$$

For $n \geq 1$ denote by $\epsilon_{n}=\min \left\{1 / 2^{n+5}, \breve{\epsilon}_{\tau}(n, l), l=1, \ldots, k_{n}\right\}$, where $\breve{\epsilon}_{\tau}(n, l)$ is defined by (5.4).

We now briefly describe how to choose the sequence $\left\{\delta_{n}\right\}$. See $[10]$ for more details. Recall that $U_{1}=Z_{1}^{(1)}, \widetilde{U}_{1}=\widetilde{Z}_{1}^{(1)}$, and $\mathcal{U}_{1}=\mathcal{N} \times U_{1}, \widetilde{\mathcal{U}}_{1}=\mathcal{N} \times \widetilde{U}_{1}$. Choose $\theta_{0}>0$ such that the families of maps $\Psi_{h_{घ}}$ and $\Theta_{h_{\natural}}^{i}$ are well defined for any gentle perturbation $h_{\sharp}^{t}$ of $f^{t}$ with $\angle\left(E_{h_{\sharp}}^{\omega}(z), E_{f}^{\omega}(z)\right) \leq 2 \theta_{0}$ for $\omega=u, s, c$. We assume that the number $\delta_{g}$ in Proposition 4.1 is so small that the flow $h_{0}^{t}=g^{t}$ satisfies $\angle\left(E_{h_{0}}^{\omega}(z), E_{f}^{\omega}(z)\right) \leq \theta_{0}$ and $d\left(\Theta_{\vartheta, 1,1, h_{0}}^{i}(z), z\right) \leq \epsilon_{1} / 4$ for $z \in \mathcal{N} \times U_{0}, \vartheta \in[0,1]$ and $i=a, b, \tau$.

Now choose $\theta_{1}^{\prime}$ with $0<\theta_{1}^{\prime} \leq \theta_{0} / 2$ such that $d\left(\Psi_{h_{घ}}(z), \Psi_{h_{0}}(z)\right) \leq 1 / 2^{8}$ if $\angle\left(E_{h_{\natural}}^{\omega}(z)\right.$, $\left.E_{h_{0}}^{\omega}(z)\right) \leq 2 \theta_{1}^{\prime}$ for all $z \in \mathcal{N} \times Z_{1}^{(1)}$. Also choose $\delta_{1}^{\prime}>0$ such that if $\left\|\mathcal{X}_{1}-\mathcal{X}_{0}\right\| \leq \delta_{1}^{\prime}$, then $\angle\left(E_{h_{1}}^{\omega}(z), E_{h_{0}}^{\omega}(z)\right) \leq \theta_{1}^{\prime}$. Finally set $\theta_{1}=\min \left\{\theta_{1}^{\prime}, \theta_{1}^{\prime \prime}\right\}$ and $\delta_{1}=\min \left\{\delta_{1}^{\prime}, \delta_{1}^{\prime \prime}, \theta_{1}\right\}$, where $\delta_{1}^{\prime \prime}$ and $\theta_{1}^{\prime \prime}$ are given by Lemma 5.3 below. We can show:
(1) $d\left(\Psi_{h_{1}}(z), \Psi_{h_{0}}(z)\right) \leq 1 / 2^{8}$ for all $z \in \mathcal{N} \times Z_{1}^{(1)}$;
(2) $d\left(\Theta_{\vartheta, 2, l, h_{1}}^{i}(z), z\right) \leq \epsilon_{2} / 4$ for all $z \in W_{h_{1}}^{c}\left(z_{0}(2, l), K, Z_{l}^{(2)}\right), i=a, b, \tau, \vartheta \in[0,1]$ and $l=1, \ldots, k_{2}$;
(3) $\mathcal{A}_{h_{\natural}}\left(z_{0}(1,1)\right) \supset W_{h_{\natural}}^{c}\left(z_{0}(1,1), \bar{K}, \bar{Z}_{1}^{(1)}\right)$ for any gentle perturbation $h_{\natural}^{t}$ of $f^{t}$, close to $h_{1}^{t}$, with $\angle\left(E_{h_{\sharp}}^{\omega}(z), E_{h_{1}}^{\omega}(z)\right) \leq \theta_{1}^{\prime}, \omega=u, s, c$ and $z \in \mathcal{N} \times Z_{1}^{(1)}$.
Moreover, the above statements imply that $\mathcal{A}_{h_{घ}}\left(z_{0}(1,1)\right) \supset \mathcal{N} \times \widetilde{Z}_{1}^{(1)}$; in particular, $h_{1}^{t}$ has the accessibility property on $\mathcal{N} \times \widetilde{Z}_{1}^{(1)}$.

Proceeding inductively, we can choose $\delta_{n}$ such that statements (5) and (6) of Proposition 5.1 hold. Furthermore, we have for $i=a, b, \tau, \vartheta \in[0,1]$ and $l=$ $1, \ldots, k_{n+1}$ :
(1) $d\left(\Psi_{h_{n}}(z), \Psi_{h_{n-1}}(z)\right) \leq 1 / 2^{n+7}$ for all $z \in \mathcal{N} \times Z_{l}^{(n)}$;
(2) $d\left(\Theta_{\vartheta, n+1, l, h_{n}}^{i}(z), z\right) \leq \epsilon_{n+1} / 4$ for all $z \in W_{h_{n}}^{c}\left(z_{0}(n+1, l), K, Z_{l}^{(n+1)}\right)$;
(3) $\mathcal{A}_{h_{\natural}}\left(z_{0}(n, l)\right) \supset W_{h_{\natural}}^{c}\left(z_{0}(n, l), \bar{K}, \bar{Z}_{l}^{(n)}\right)$ for any gentle perturbation $h_{\natural}^{t}$ of $f^{t}$, close to $h_{n}^{t}$, with $\angle\left(E_{h_{\sharp}}^{\omega}(z), E_{h_{n}}^{\omega}(z)\right) \leq \delta_{n}, \omega=u, s, c$ and $z \in \mathcal{N} \times Z_{l}^{(n)}$.
Therefore, $\mathcal{A}_{h_{\natural}}\left(z_{0}(n, l)\right) \supset \mathcal{N} \times \widetilde{Z}_{l}^{(n)}$ for all $l=1, \ldots, k_{n+1}$. In other words, $h_{\natural}^{t}$ has the accessibility property on $\mathcal{N} \times \widetilde{Z}_{l}^{(n)}$.

Note that

$$
\tilde{\mathcal{U}}_{n}=\tilde{\mathcal{U}}_{n-1} \cup\left(\bigcup_{l=1}^{k_{n}} \mathcal{N} \times \widetilde{Z}_{l}^{(n)}\right)
$$

and the intersection of any two sets among $\tilde{\mathcal{U}}_{n-1}$ and $\mathcal{N} \times \widetilde{Z}_{l}^{(n)}, l=1, \ldots, k_{n}$, contains a non-empty open set whenever they intersect. Since $\tilde{\mathcal{U}}_{n}$ is connected, we obtain accessibility of $h_{\natural}^{t}$ on $\widetilde{\mathcal{U}}_{n}$. In particular, $h_{n}^{t}$ has the accessibility property on $\widetilde{\mathcal{U}}_{n}$ when we apply that $h_{\square}^{t}=h_{n}^{t}$.
5.4. A technical lemma. The proof of Proposition 5.1 heavily relies on the following technical statements.
Lemma 5.3. Suppose for some $n>0, d\left(\Theta_{\vartheta, n, l, h_{n-1}}^{i}(z), z\right) \leq \epsilon_{n} / 4$ for all $i=a, b, \tau$, $\vartheta \in[0,1], z \in W_{h_{n-1}}^{c}\left(z_{0}(n, l), K, Z_{l}^{(n)}\right), l=1, \ldots, k_{n}$. Then there are $\delta_{n}^{\prime \prime}, \theta_{n}^{\prime \prime}>0$ such that the following hold.
(1) If $\left\|\mathcal{X}_{n}-\mathcal{X}_{n-1}\right\|_{C^{n}} \leq \delta_{n}^{\prime \prime}$, then

$$
\begin{equation*}
d\left(\Theta_{\vartheta, n+1, l, h_{n}}^{i}(z), z\right) \leq \epsilon_{n+1} / 4 \quad \text { as } z \in W_{h_{n}}^{c}\left(z_{0}(n+1), K, Z_{l}^{(n+1)}\right) \text {, } \tag{5.7}
\end{equation*}
$$

for all $i=a, b, \tau, \vartheta \in[0,1]$, and $l=1, \ldots, k_{n+1}$.
(2) For any gentle perturbation $h_{\square}^{t}$ of $f^{t}$, close to $h_{n}^{t}$ and with

$$
\angle\left(E_{h_{\sharp}}^{\omega}(z), E_{h_{n}}^{\omega}(z)\right) \leq \theta_{n}^{\prime \prime} \quad \text { for all } z \in \mathcal{N} \times Z_{l}^{(n)}, \omega=u, s, c,
$$

we have

$$
\begin{equation*}
\mathcal{A}_{h_{\sharp}}\left(z_{0}(n, l)\right) \supset W_{h_{\sharp}}^{c}\left(z_{0}(n, l), \bar{K}, \bar{Z}_{l}^{(n)}\right) \quad \text { for all } l=1, \ldots, k_{n} . \tag{5.8}
\end{equation*}
$$

In particular, (5.8) holds with $h_{\natural}^{t}=h_{n}^{t}$.
This lemma is an adaptation of [10, Lemma 5.2] to the flow case. It can be proved in a similar fashion subject to the following sublemma.

Sublemma 5.4. For each $n>0$, there exists $\delta_{n}^{\prime \prime}>0$ such that if $\left\|\mathcal{X}_{n}-\mathcal{X}_{n-1}\right\|_{C^{n}}=$ $\beta_{n}\left\|\widehat{\mathcal{X}}_{n}\right\|_{C^{n}} \leq \delta_{n}^{\prime \prime}$, then, for all $l=1, \ldots, k_{n}$ :
(1) $\Theta^{a}((q, 1 / 2, a, 0))=\left(q, 1 / 2, a^{\prime}, 0\right)$ with $a^{\prime}<a$ for any $a \in I_{n}$;
(2) $\Theta^{b}((q, 1 / 2, a, b))=\left(q, 1 / 2, a, b^{\prime}\right)$ with $b^{\prime}<b$ for any $a \in I_{n}, b \in J_{n}$;
(3) $\Theta^{\tau}((q, \tau, a, b))=\left(q, \tau^{\prime}, a, b\right)$ with $\tau^{\prime}<\tau$ for any $a \in I_{n}, b \in J_{n}$ and $\tau \in K$, where $q=q(n, l)$ is given by (5.2), and $\Theta^{i}=\Theta_{n, l, h_{n}}^{i}$ for $i=a, b, \tau$.

Proof. The proof is similar to that in [6] (see Lemma B.4; see also [10, Sublemma 5.3]) but is adapted to the language of vector fields.

We will prove the first statement. Consider the coordinate system ( $u, s, \tau, a, b$ ) in $\Omega_{n, l}^{a}$ with the origin at $\left(p^{a}(n, l), 1 / 2, y_{0}(n, l)\right)$. Write $q=q(n, l), p^{a}=p^{a}(n, l)$ and $y_{0}=y_{0}(n, l)$. We may assume that the square $Z_{l}^{(n)}$ is parameterized as $(a, b) \in I_{n} \times J_{n}$, that the center $y_{0}(n, l)$ of $Z_{l}^{(n)}$ is $(0,0)$ and that the local coordinates of the points $q,\left[q, p^{a}\right],\left[p^{a}, q\right]$ and $p^{a}$ are $\left(u_{0}, s_{0}\right),\left(0, s_{0}\right),\left(u_{0}, 0\right)$ and $(0,0)$ respectively, where $u_{0}=\alpha_{u}^{a}(n, l)$ and $s_{0}=\alpha_{s}^{a}(n, l)$ given by (5.3).

Consider the case $n>1$ first. Note that the vector field $\mathcal{X}_{n}$ inside $\Omega^{a}(n, l)$ is exactly $\mathcal{X}_{f}+\beta_{n} \mathcal{X}_{n, l}^{a}$. For any $a_{1}=a \in I_{n}, b \in J_{n}$, and $\tau \in\left(1 / 2-\epsilon_{\tau}, 1 / 2+\epsilon_{\tau}\right)$, choose the point $\underline{z}_{1}=\left(q, \tau, a_{1}, b\right)=\left(u_{0}, s_{0}, \tau, a_{1}, b\right)$. Note that under the original flow $f^{t}$, we have a closed quadrilateral $(u, s)_{f^{t}}$-path $\underline{\gamma}=\left\{\underline{z}_{1}, \underline{z}_{2}, \underline{z}_{3}, \underline{z}_{4}, \underline{z}_{5}\right\}$, where

$$
\begin{aligned}
& \underline{z}_{2}=\left(\left[q, p^{a}\right], \tau, \underline{a}_{2}, b\right)=\left(0, s_{0}, \tau, \underline{a}_{2}, b\right), \\
& \underline{z}_{3}=\left(p^{a}, \tau, \underline{a}_{3}, b\right)=\left(0,0, \tau, \underline{a}_{3}, b\right), \\
& \underline{z}_{4}=\left(\left[p^{a}, q\right], \tau, \underline{a}_{4}, b\right)=\left(u_{0}, 0, \tau, \underline{a}_{4}, b\right), \\
& \underline{z}_{5}=\left(\left[q, p^{a}\right], \tau, \underline{a}_{5}, b\right)=\left(u_{0}, s_{0}, \tau, \underline{a}_{5}, b\right)=\underline{z}_{1},
\end{aligned}
$$

and $\underline{a}_{k}=a_{1}=a$ for $k=1,2,3,4,5$.
Let us compare the vector field $\mathcal{X}_{n}=\mathcal{X}_{f}+\beta_{n} \mathcal{X}_{n, l}^{a}$ on each leg $\mathcal{L}_{k}=\left[\underline{z}_{k}, \underline{z}_{k+1}\right]$ for $k=1,2,3,4$. In fact, $\mathcal{X}_{n, l}^{a} \equiv 0$ on legs $\mathcal{L}_{1}$ and $\mathcal{L}_{4}$. Since the $u$-component of every point on the $\operatorname{leg} \mathcal{L}_{2}$ is 0 , the $u$-component of the vector field $\mathcal{X}_{n, l}^{a}$ is 0 , and the $a$-component does not depend on the $u$-coordinate. On the leg $\mathcal{L}_{3}=\left[\underline{z}_{3}, \underline{z}_{4}\right]$, the $u$-component of $\mathcal{X}_{n, l}^{a}$ is negative at the interior points and it is zero at two endpoints $\underline{z}_{3}$ and $\underline{z}_{4}$, while the $a$-component is positive, with the value smoothly changing from a constant to zero.

Now choose the point $z_{1}=\underline{z}_{1}$, and we have the quadrilateral $(u, s)_{h_{n}^{t}}$-path $\gamma=$ $\left\{z_{1}, z_{2}, z_{3}, z_{4}, z_{5}\right\}$. By the above comparison, the $\tau$ - and $b$-coordinate are the same for each $z_{k}, k=1,2,3,4,5$. By the construction of the vector fields $\mathcal{X}_{n}$, the image of the leg $\left[z_{3}, z_{4}\right]$ under the flow $h_{n}^{t}$ is contained in $\Omega^{a}$. Now let $a_{k}$ be the $a$-coordinate of $z_{k}, k=1,2,3,4,5$. Since the $a$-component of $\mathcal{X}_{n}$ is the same for all points on $\mathcal{L}_{1}, \mathcal{L}_{2}$
and $\mathcal{L}_{4}$, and it changes from a constant to zero along the unstable leg $\mathcal{L}_{3}$, then we have $a_{1}=a_{2}=a_{3}>a_{4}=a_{5}$. This shows statement (1) for the case $n>1$.

In the case $n=1$, similar arguments can be used with the following modification, and we will obtain $a_{1}=a_{2} \geq a_{3}>a_{4}=a_{5}$. This completes the proof of statement (1). Statements (2) and (3) can be proved in a similar way.

## 6. Proof of main theorem

Since each $\mathcal{X}_{n}$ is divergence free, so is $\mathcal{X}_{h}$, and hence $h^{t}$ is volume preserving. The first statement of the main theorem follows.

Note that $h^{t}=f^{t}$ on $\mathcal{U}^{c}$ and is of the form

$$
h^{t}((x, \tau), y)=\left(\left(x+t \alpha_{0}, \tau\right), y\right)
$$

for each $z=((x, \tau), y) \in \mathcal{U}^{c}=\mathcal{N} \times C$, where $\alpha_{0}$ is a Diophantine vector (see $\left.\S 3\right)$. Hence, $h^{t}$ preserves each three-dimensional submanifold $\mathcal{N} \times\{y\}, y \in C$, and $h^{t} \mid \mathcal{N} \times\{y\}$ is a non-identity linear flow since $\alpha(y) \neq 0$. Moreover, the frequency vector $\alpha(y)$ is Diophantine if $y \in C$. Thus statements (1) and (3) of the main theorem follow.

It remains to prove the second statement. By Proposition 4.1, each diffeomorphism $h_{n}^{t}$ is pointwise partially hyperbolic on $\mathcal{U}$ and uniformly partially hyperbolic on $\overline{\mathcal{U}}_{n}$. By Theorem B. 1 in Appendix B, if the sequence $\delta_{n}$ decreases sufficiently fast, the limit flow $h^{t}$ is pointwise partially hyperbolic on $\mathcal{U}$.

We now claim that the one-dimensional strongly stable $E_{h}^{s}$ and unstable $E_{h}^{u}$ subbundles are integrable to invariant strongly stable $W_{h}^{s}$ and unstable $W_{h}^{u}$ foliations with smooth leaves, which are transversal and absolutely continuous. Recall that the 'start-up' flow $f^{t}$ has strongly stable and unstable local manifolds $V_{f}^{s}(z)$ and $V_{f}^{u}(z)$ respectively at each $z \in \mathcal{U}$. Moreover, these local manifolds are of uniform size, larger, say, than a certain number $4 r>0$. By Proposition 5.1(3), $h_{n}^{t}\left|\mathcal{U}_{n}^{c}=f^{t}\right| \mathcal{U}_{n}^{c}$, and thus $V_{h_{n}}^{\omega}(z)=V_{f}^{\omega}(z)$ for all $z \in \mathcal{U} \backslash \mathcal{U}_{n}, \omega=s, u$. On the other hand, each $h_{n}^{t}$ is a perturbation of $h_{n-1}^{t}$ on the compact set $\mathcal{U}_{n}$, on which both $h_{n}^{t}$ and $h_{n-1}^{t}$ are uniformly partially hyperbolic if $\delta_{n}$ is sufficiently small. Furthermore, let $r_{n}$ be the size of $V_{h_{n}}^{\omega}(z)$ for $z \in \mathcal{U}_{n}$. One can have $r_{n} / r_{n-1} \geq 2^{-1 / 2^{n}}$, and thus by induction we have that the size of local manifolds of $h_{n}^{t} \mid \mathcal{U}_{n}$ is bigger than $r$. Therefore, given $z \in \mathcal{U}$, we obtain that the size of $V_{h_{n}}^{\omega}(z)$ has a lower bound $r>0$, which is independent of $z$ and $n$.

Write each $V_{h_{n}}^{s}(z)$ in the coordinate chart as

$$
V_{h_{n}}^{s}(z)=\exp _{z}\left\{\left(v, \psi_{h_{n}}^{s}(v)\right): v \in B^{s}\left(0, r_{n}\right)\right\},
$$

where $B^{s}\left(0, r_{n}\right) \subset E_{h_{n}}^{s}(z)$ is the ball centered at origin of radius $r_{n}$ and $\psi_{h_{n}}^{s}: B^{s}\left(0, r_{n}\right) \rightarrow$ $E_{h_{n}}^{c u}(z)$ is a $C^{1}$ map satisfying the following.
(1) $\psi_{h_{n}}^{s}(0)=0$ and $d \psi_{h_{n}}^{s}(0)=0$.
(2) If the numbers $\delta_{n}$ and $\theta_{n}$ decay sufficiently fast then there are $r>0$ and $\Delta>0$ such that $r_{n} \geq r$ and $\left\|\psi_{h_{n}}^{s}\right\|_{C^{1}} \leq \Delta$ for all $n \geq 0$.
This implies that $z \in V_{h_{n}}^{s}(z)$ and $T_{z} V_{h_{n}}^{s}(z)=E_{h_{n}}^{s}(z)$. Furthermore:
(1) $h_{n}^{t}\left(V_{h_{n}}^{s}(z)\right) \subset V_{h_{n}}^{s}\left(h_{n}^{t}(z)\right)$;
(2) $\quad d\left(h_{n}^{t}(z), h_{n}^{t}(y)\right) \leq \tilde{\lambda}(z) d(z, y)$ for each $y \in V_{h_{n}}^{s}(z)$ and some continuous function $\tilde{\lambda}(z)$ on $\mathcal{U}$ for which $0<\lambda(z) \leq \tilde{\lambda}(z)<\lambda^{\prime}(z)$ (where $\lambda^{\prime}(z)$ is the function in the definition of pointwise partial hyperbolicity).
The sequence of functions $\psi_{h_{n}}^{s}(v),\|v\| \leq r$, is compact in the $C^{1}$ topology, and hence there is a subsequence $\psi_{h_{n_{k}}}^{s}$ that converges to a $C^{1}$ function $\psi$ satisfying $\psi(0)=0, d \psi(0)=0$ and $\|\psi\|_{C^{1}} \leq \Delta$. Setting

$$
\begin{equation*}
V(z)=\exp _{z}\left\{(v, \psi(v)): v \in B^{s}(0, r)\right\} \tag{6.1}
\end{equation*}
$$

we have:
(1) $z \in V(z)$ and $T_{z} V(z)=E_{h}^{s}(z)$;
(2) $h^{t}(V(z)) \subset V\left(h^{t}(z)\right)$;
(3) $\quad d\left(h^{t}(z), h^{t}(y)\right) \leq \tilde{\lambda}(z) d(z, y)$ for each $y \in V(z)$.

This implies that if $m_{k}$ is any subsequence for which $\psi_{h_{m_{k}}}^{s}$ converges in the $C^{1}$ topology to a function $\tilde{\psi}$, then $\tilde{\psi}=\psi$. Thus the formula (6.1) determines uniquely a local strongly stable manifold through $z$ and the formula $W(z)=\bigcup_{n \geq 0} h^{-t}\left(V\left(h^{t}(z)\right)\right)$ defines the global strongly stable manifold through $z$. These manifolds form a continuous strongly stable foliation with smooth leaves for $h^{t}$. In a similar fashion we can obtain strongly unstable local manifolds and construct a strongly unstable foliation with smooth leaves for $h^{t}$. These two foliations are transverse at every point $z \in \mathcal{U}$.

We will now show that the Lyapunov exponent $\lambda_{h}^{s}(z)$ in the direction $E_{h}^{s}(z)$ is negative at almost every point $z \in \mathcal{U}$. Indeed, let $Z \subset \mathcal{U}$ be the set of points at which $\lambda_{h}^{s}(z)=0$. If $m(Z)>0$ then

$$
\begin{aligned}
0 & =\int_{Z} \lambda_{h}^{s}(z) d m=\int_{Z} \lim _{n \rightarrow \infty} \frac{1}{n} \log \prod_{i=0}^{n-1} \lambda_{h}\left(h^{i}(z)\right) d m(z) \\
& =\lim _{n \rightarrow \infty} \frac{1}{n} \int_{Z} \sum_{i=0}^{n-1} \log \lambda_{h}\left(h^{i}(z)\right) d m(z) \\
& =\int_{Z} \log \lambda_{h}(z) d m(z)<0
\end{aligned}
$$

(recall that $\lambda_{h}(z)$ is the contraction coefficient along $\left.E_{h}^{s}(z)\right)$. This contradiction proves our claim. Similarly, one can prove that the Lyapunov exponent $\lambda_{h}^{u}(z)$ in the direction $E_{h}^{u}(z)$ is positive at almost every point $z \in \mathcal{U}$.

Since $h^{t}$ is non-uniformly partially hyperbolic on $\mathcal{U}$, by [2, Theorem 8.6.1], we obtain that its strongly stable and unstable foliations are absolutely continuous.

Next we will show that the flow $h^{t}$ has the accessibility property on $\mathcal{U}$ via its invariant foliations $W_{h}^{s}$ and $W_{h}^{u}$. Indeed, by Proposition 5.1(5), for any $n>k, \omega=s, u, c$, and any $z \in \mathcal{U}_{k} \subset \mathcal{U}$,

$$
\angle\left(E_{h_{n}}^{\omega}(z), E_{h_{k}}^{\omega}(z)\right) \leq \delta_{k}\left(1-1 / 2^{n-k}\right)<\delta_{k} .
$$

Taking the limit as $n \rightarrow \infty$, we obtain that $\angle\left(E_{h}^{\omega}(z), E_{h_{k}}^{\omega}(z)\right) \leq \delta_{k}$ on $\mathcal{U}_{k}$. Hence, by Proposition 5.1(6), the flow $h^{t}$ has the accessibility property on each $\tilde{\mathcal{U}}_{k}$. Since $k$ is arbitrary, we obtain that the flow $h^{t}$ has the accessibility property on $\mathcal{U}$.

To show that the flow $h^{t}$ has positive central Lyapunov exponent, we first recall that the average Lyapunov exponents of the flow $g^{t}$ are arranged as in (4.2). Set $c=$ $L_{3}\left(g^{t}\right)-L_{2}\left(g^{t}\right)>0$. By the upper semicontinuity of $L_{i}(\cdot)$, we choose the number $\delta_{h}>0$ in Proposition 5.1 to be so small that $L_{2}\left(h^{t}\right)<L_{2}\left(g^{t}\right)+c / 2$. On the other hand, it follows from Proposition 5.1(4) that

$$
L_{4}\left(h_{n}^{t}\right)=\int_{\mathcal{M}} \operatorname{det}\left(d h_{n}^{t} \mid E_{h_{n}}^{u a b \tau}(z)\right) d m=\int_{\mathcal{M}} \operatorname{det}\left(d g^{t} \mid E_{g}^{u a b \tau}(z)\right) d m=L_{4}\left(g^{t}\right)
$$

Taking the limit as $n \rightarrow \infty$ we obtain $L_{4}\left(h^{t}\right)=L_{4}\left(g^{t}\right)=L_{3}\left(g^{t}\right)$. Therefore,

$$
L_{4}\left(h^{t}\right)-L_{2}\left(h^{t}\right)=\int_{\mathcal{M}}\left(\lambda_{3}\left(z, h^{t}\right)+\lambda_{4}\left(z, h^{t}\right)\right) d m(z) \geq c / 2>0
$$

Then there is a subset $\mathcal{A} \subset \mathcal{U}$ such that $\lambda_{3}\left(z, h^{t}\right)+\lambda_{4}\left(z, h^{t}\right)>0$ for all $z \in \mathcal{A}$, and thus $\lambda_{2}\left(z, h^{t}\right) \geq \lambda_{3}\left(z, h^{t}\right) \geq \frac{1}{2}\left[\lambda_{3}\left(z, h^{t}\right)+\lambda_{4}\left(z, h^{t}\right)\right]>0$ for all $z \in \mathcal{A}$. Since the center subspace $E_{h}^{c}(z)$ is three-dimensional and the flow direction $\operatorname{Span}\left\{\mathcal{X}_{h}\right\}$ corresponds to a zero exponent, we conclude that $\lambda_{2}\left(z, h^{t}\right)$ and $\lambda_{3}\left(z, h^{t}\right)$ correspond to vectors in $E_{h}^{c}(z) \backslash \operatorname{Span}\left\{\mathcal{X}_{h}\right\}$. Thus the flow $h^{t}$ has positive central Lyapunov exponents.

By Theorem 2.4, we obtain that $h^{t}$ has positive central exponents at almost every point in $\mathcal{U}$, and $h^{t} \mid \mathcal{U}$ is an ergodic flow.

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## A. Appendix. The differential and metric structures of the suspension manifold

We specify the differential and metric structure of the suspension manifold $\mathcal{N}$ and the fivedimensional manifold $\mathcal{M}$ in $\S 3$. Associated to the Anosov automorphism $A$ of $X=\mathbb{T}^{2}$, one can find smooth local charts ( $U_{x}, \phi_{x}$ ) around each $x \in X$ such that

$$
\phi_{x}: U_{x} \rightarrow\left(-u_{0}(x), u_{0}(x)\right) \times\left(-s_{0}(x), s_{0}(x)\right)
$$

satisfies $\phi_{x}(x)=(0,0)$ and

$$
\phi_{A x} \circ A \circ \phi_{x}^{-1}(u, s)=\left(\eta u, \eta^{-1} s\right),
$$

where $u_{0}(x), s_{0}(x)>0$ are sizes of charts depending on $x$, and $\eta>1$ is the expanding rate along the unstable direction. In fact, $\partial / \partial u$ and $\partial / \partial s$ are the unstable and stable directions of $A$ respectively.

Recall that the suspension manifold $\mathcal{N}$ is the quotient space $X \times \mathbb{R} / \sim$ with the equivalence relation $(x, \tau+1) \sim(A x, \tau)$. Let $\pi: X \times \mathbb{R} \rightarrow \mathcal{N}$ be the natural projection. Following [12] there is a natural differential structure on $\mathcal{N}$ with atlas $\left(U_{(x, \tau)}^{1}, \phi_{(x, \tau)}^{1}\right)$ for $\tau \in(-1 / 4,3 / 4)$ and $\left(U_{(x, \tau)}^{2}, \phi_{(x, \tau)}^{2}\right)$ for $\tau \in(1 / 4,5 / 4)$, where

$$
\begin{array}{cc}
U_{(x, \tau)}^{1}=\pi\left(U_{x} \times(-1 / 4,3 / 4)\right), & \phi_{(x, \tau)}^{1}\left(\pi\left(\phi_{x}^{-1}(u, s), \tau\right)\right)=(u, s, \tau), \\
U_{(x, \tau)}^{2}=\pi\left(U_{x} \times(1 / 4,5 / 4)\right), & \phi_{(x, \tau)}^{2}\left(\pi\left(\phi_{x}^{-1}(u, s), \tau\right)\right)=(u, s, \tau) .
\end{array}
$$

It is easy to verify that

$$
\begin{gathered}
\phi_{\left(x^{\prime}, \tau^{\prime}\right)}^{i} \circ \phi_{(x, \tau)}^{i}\left(\pi\left(\phi_{x}^{-1}(u, s), \tau\right)\right)=\left(\phi_{x^{\prime}} \circ \phi_{x}^{-1}(u, s), \tau\right), \quad i=1,2, \\
\phi_{\left(x^{\prime}, \tau^{\prime}\right)}^{1} \circ \phi_{(x, \tau)}^{2}\left(\pi\left(\phi_{x}^{-1}(u, s), \tau\right)\right)=\left(\phi_{A x^{\prime}} \circ A \circ \phi_{x}^{-1}(u, s), \tau-1\right) .
\end{gathered}
$$

In particular,

$$
\begin{equation*}
\phi_{\left(x, \tau^{\prime}\right)}^{1} \circ \phi_{(x, \tau)}^{2}\left(\pi\left(\phi_{x}^{-1}(u, s), \tau\right)\right)=\left(\eta u, \eta^{-1} s, \tau-1\right) . \tag{A.1}
\end{equation*}
$$

There are three subbundles $E^{u}, E^{s}$ and $E^{\tau}$ on $\mathcal{N}$ generated by independent vector fields $d \pi(\partial / \partial u), d \pi(\partial / \partial s)$ and $d \pi(\partial / \partial \tau)$ respectively. By [3], we can choose the Riemannian metric on $\mathcal{N}$ which has the local representation $\eta^{2 \tau} d u^{2}+\eta^{-2 \tau} d s^{2}+d \tau^{2}$. Under this metric, the suspension flow $S^{t}: \mathcal{N} \rightarrow \mathcal{N}$ satisfies

$$
\begin{gathered}
\left\|d S^{t} v\right\|=\eta^{t}\|v\|, \quad v \in E^{u}, \\
\left\|d S^{t} v\right\|=\eta^{-t}\|v\|, \quad v \in E^{s}, \\
\left\|d S^{t} v\right\|=\|v\|, \quad v \in E^{\tau} .
\end{gathered}
$$

For the 2-torus $Y=\mathbb{T}^{2}$, we choose the local coordinate $(a, b)$ centered at each $y \in Y$. Given $z=(x, \tau, y) \in \mathcal{M}=\mathcal{N} \times Y$, we can hence choose a local coordinate system $(u, s, \tau, a, b)$ endowed with the product Riemannian metric $\eta^{2 \tau} d u^{2}+\eta^{-2 \tau} d s^{2}+d \tau^{2}+$ $d a^{2}+d b^{2}$. In particular, the metric on the cross-section $X \times\{0\} \times Y$ is given by the flat metric

$$
\begin{equation*}
d u^{2}+d s^{2}+d a^{2}+d b^{2} . \tag{A.2}
\end{equation*}
$$

## B. Appendix. Partial hyperbolicity of the limit flow

Let $\mathcal{M}$ be a compact smooth Riemannian manifold and $\mathcal{S} \subset \mathcal{M}$ an open subset. Let also $H^{t}$ be the flow on $\mathcal{M}$ that is pointwise partially hyperbolic on $\mathcal{S}$. Further, let $\mathcal{U}_{n} \subset \mathcal{S}$, $n \geq 1$ be a sequence of open subsets such that the following hold.
(1) $\mathcal{U}_{n} \subset \overline{\mathcal{U}}_{n} \subset \mathcal{U}_{n+1}$ and $\bigcup \mathcal{U}_{n}=\mathcal{S}$.
(2) Each $\mathcal{U}_{n}$ is $H^{t}$-invariant.
(3) $H^{t} \mid \overline{\mathcal{U}_{n}}$ is uniformly partially hyperbolic.

The goal of this appendix is to prove the following statement.
THEOREM B.1. There exists a sequence of positive numbers $\varepsilon_{n}$ such that if smooth vector fields $\mathcal{X}_{n}$ on $\mathcal{M}$ satisfy

$$
\mathcal{X}_{0}=\mathcal{X}_{H}, \quad \mathcal{X}_{n}=\mathcal{X}_{n-1} \quad \text { on } \mathcal{M} \backslash \overline{\mathcal{U}_{n}}
$$

and

$$
\left\|\mathcal{X}_{n}-\mathcal{X}_{n-1}\right\|_{C^{1}} \leq \varepsilon_{n},
$$

then for every $n \geq 1$ the corresponding flow $h_{n}^{t}$ is uniformly partially hyperbolic on the invariant set $\overline{\mathcal{U}_{n}}$ and hence pointwise partially hyperbolic on $\mathcal{S}$. Moreover, the limit vector field $\mathcal{X}=\lim _{n \rightarrow \infty} \mathcal{X}_{n}$ is of class $C^{1}$ and generates a pointwise partially hyperbolic flow $h^{t}$ on $\mathcal{S}$.

We need the following technical statements.
Lemma B.2. Given a sequence of positive numbers $\left\{a_{n}\right\}_{n \geq 1}$ satisfying $\sum_{n=1}^{\infty} a_{n} \leq \frac{1}{4}$, we have

$$
\prod_{n=1}^{\infty}\left(1+a_{n}\right) \leq 1+2 \sum_{n=1}^{\infty} a_{n}, \quad \prod_{n=1}^{\infty}\left(1-a_{n}\right) \geq 1-2 \sum_{n=1}^{\infty} a_{n} .
$$

Lemma B.3. (Gronwall's inequality) Let $\eta(t)$ be a non-negative $C^{1}$ function on $[0, T]$ satisfying

$$
\eta^{\prime}(t) \leq \phi(t) \eta(t)+\psi(t)
$$

where $\phi(t)$ and $\psi(t)$ are non-negative integrable functions. Then, for all $0 \leq t \leq T$,

$$
\eta(t) \leq e^{\int_{0}^{t} \phi(s) d s}\left[\eta(0)+\int_{0}^{t} \psi(s) d s\right] .
$$

Lemma B.4. Set $K:=2\left\|\mathcal{X}_{H}\right\|_{C^{1}}$. If $\varepsilon_{n}<K / 2^{n+1}$, then $\left\|\mathcal{X}_{n}\right\|_{C^{1}} \leq K$ for all $n \geq 0$. Moreover, given a flow $F^{t}$ with $\left\|\mathcal{X}_{F}\right\|_{C^{1}} \leq K$, we have for any $x \in \mathcal{M}$ and $t \in \mathbb{R}^{+}$

$$
e^{-t K} \leq m\left(d_{x} F^{t}\right) \leq\left\|d_{x} F^{t}\right\| \leq e^{t K}
$$

In particular,

$$
e^{-K} \leq \min _{0 \leq t \leq 1} m\left(d_{x} F^{t}\right) \leq \max _{0 \leq t \leq 1}\left\|d_{x} F^{t}\right\| \leq e^{K}
$$

Proof. Since $\left\|\mathcal{X}_{H}\right\|_{C^{1}}=K / 2$, we find that

$$
\left\|\mathcal{X}_{n}-\mathcal{X}_{H}\right\|_{C^{1}} \leq \sum_{k=1}^{n}\left\|\mathcal{X}_{k}-\mathcal{X}_{k-1}\right\|_{C^{1}} \leq \sum_{k=1}^{n} \varepsilon_{k}<K / 2
$$

This implies that $\left\|\mathcal{X}_{n}\right\|_{C^{1}} \leq K$.
Let $F^{t}$ be a flow and $\mathcal{X}_{F}$ the corresponding vector field. Consider the variational differential equation

$$
\frac{d}{d t} d_{x} F^{t}=D \mathcal{X}_{F}\left(F^{t}(x)\right) d_{x} F^{t}
$$

for any $x \in \mathcal{M}$ and $t \in \mathbb{R}^{+}$. Then

$$
\frac{d}{d t}\left\|d_{x} F^{t}\right\| \leq\left\|\frac{d}{d t} d_{x} F^{t}\right\| \leq\left\|D \mathcal{X}_{F}\right\|\left\|d_{x} F^{t}\right\| \leq K\left\|d_{x} F^{t}\right\| .
$$

Since $\left\|d_{x} F^{0}\right\|=1$, by Lemma B.3, we obtain that $\left\|d_{x} F^{t}\right\| \leq e^{t K}$. Noting that $m\left(d_{x} F^{t}\right)=$ $\left\|d_{x} F^{-t}\right\|^{-1}$ and the flow $F^{-t}$ corresponds to the vector field $-\mathcal{X}_{F}$, we get that $m\left(d_{x} F^{t}\right) \geq$ $e^{-t K}$.

Let $F^{t}$ and $G^{t}$ be flows on $\mathcal{M}$, and $\mathcal{X}_{F}$ and $\mathcal{X}_{G}$ the corresponding vector fields. Assume that $\left\|\mathcal{X}_{F}\right\|_{C^{1}},\left\|\mathcal{X}_{G}\right\|_{C^{1}} \leq K$, where $K$ is given in Lemma B.4.
Lemma B.5. Set

$$
M:=e^{3 K}\left\|\mathcal{X}_{F}\right\|_{C^{2}}+e^{2 K}, \quad \varepsilon_{F, G}:=\left\|\mathcal{X}_{G}-\mathcal{X}_{F}\right\|_{C^{1}}
$$

Then for $t \in[0,1]$,

$$
\begin{equation*}
\rho_{C^{1}}\left(G^{t}, F^{t}\right) \leq 2 t M \varepsilon_{F, G}, \tag{B.1}
\end{equation*}
$$

where $\rho_{C^{1}}\left(G^{t}, F^{t}\right)=\max _{x \in \mathcal{M}}\left(\operatorname{dist}\left(G^{t}(x), F^{t}(x)\right)+\left\|d_{x} F^{t}-d_{x} G^{t}\right\|\right)$ is the distance between the flows in the $C^{1}$ topology.

Proof. Consider the family of flows $F^{t}(\tau)$ generated by the family of vector fields $(1-\tau) \mathcal{X}_{F}+\tau \mathcal{X}_{G}$ with $\tau \in[0,1]$. Given $x \in \mathcal{M}$ and $t \in[0,1]$, the curve $c_{t}=c_{t}^{x}: \tau \mapsto$ $F^{t}(\tau)(x)$ is of length

$$
L\left(c_{t}\right)=\int_{0}^{1}\left\|\frac{\partial c_{t}}{\partial \tau}\right\| d \tau=\int_{0}^{1}\left\|\frac{\partial}{\partial \tau} F^{t}(\tau)(x)\right\| d \tau
$$

and hence

$$
\begin{aligned}
\frac{d}{d t} L\left(c_{t}\right) & \leq \int_{0}^{1}\left\|\frac{\partial}{\partial \tau} \frac{d}{d t} F^{t}(\tau)(x)\right\| d \tau \\
& \leq\left\|\mathcal{X}_{G}-\mathcal{X}_{F}\right\|+\left[(1-\tau)\left\|D \mathcal{X}_{F}\right\|+\tau\left\|D \mathcal{X}_{G}\right\|\right] \int_{0}^{1}\left\|\frac{\partial c_{t}}{\partial \tau}\right\| d \tau \\
& \leq \varepsilon_{F, G}+K L\left(c_{t}\right)
\end{aligned}
$$

Recall that $L\left(c_{0}\right)=0, c_{t}(0)=F^{t} x$ and $c_{t}(1)=G^{t} x$. By Lemma B.3, we obtain

$$
\begin{equation*}
\operatorname{dist}\left(F^{t} x, G^{t} x\right) \leq L\left(c_{t}\right) \leq t e^{t K} \varepsilon_{F, G} \leq t M \varepsilon_{F, G} \tag{B.2}
\end{equation*}
$$

On the other hand,

$$
\begin{aligned}
\frac{d}{d t}\left\|d_{x} F^{t}-d_{x} G^{t}\right\| \leq & \left\|D \mathcal{X}_{F}\left(F^{t} x\right) d_{x} F^{t}-D \mathcal{X}_{G}\left(G^{t} x\right) d_{x} G^{t}\right\| \\
\leq & \left\|D \mathcal{X}_{F}\left(F^{t} x\right) d_{x} F^{t}-D \mathcal{X}_{F}\left(F^{t} x\right) d_{x} G^{t}\right\| \\
& +\left\|D \mathcal{X}_{F}\left(F^{t} x\right) d_{x} G^{t}-D \mathcal{X}_{F}\left(G^{t} x\right) d_{x} G^{t}\right\| \\
& +\left\|D \mathcal{X}_{F}\left(G^{t} x\right) d_{x} G^{t}-D \mathcal{X}_{G}\left(G^{t} x\right) d_{x} G^{t}\right\| \\
\leq & \left\|D \mathcal{X}_{F}\right\|\left\|d_{x} F^{t}-d_{x} G^{t}\right\| \\
& +\left\|D^{2} \mathcal{X}_{F}\right\| \operatorname{dist}\left(F^{t} x, G^{t} x\right)\left\|d_{x} G^{t}\right\| \\
& +\left\|D \mathcal{X}_{F}-D \mathcal{X}_{G}\right\|\left\|d_{x} G^{t}\right\| \\
\leq & K\left\|d_{x} F^{t}-d_{x} G^{t}\right\|+M e^{-K} \varepsilon_{F, G},
\end{aligned}
$$

where in the last inequality we use Lemma B. 4 and the inequalities (B.2). By Lemma B.3, we obtain

$$
\begin{equation*}
\left\|d_{x} F^{t}-d_{x} G^{t}\right\| \leq t M e^{(t-1) K} \varepsilon_{F, G} \leq t M \varepsilon_{F, G} . \tag{B.3}
\end{equation*}
$$

Now (B.1) follows by combining (B.2) and (B.3).
Given flows $F^{t}$ and $G^{t}$ and invariant distributions $E_{F}$ and $E_{G}$ on $\mathcal{S}$ respectively, let

$$
\begin{align*}
\Delta_{F^{t}, G^{t}, E_{F}, E_{G}}(x) & =\max \left\{\left|\frac{\left\|d_{x} G^{t} \mid E_{G}(x)\right\|}{\left\|d_{x} F^{t} \mid E_{F}(x)\right\|}-1\right|,\left|\frac{m\left(d_{x} G^{t} \mid E_{G}(x)\right)}{m\left(d_{x} F^{t} \mid E_{F}(x)\right)}-1\right|\right\},  \tag{B.4}\\
\delta_{F^{t}, G^{t}} & =\left\|G^{t}-F^{t}\right\|_{C^{1}}, \quad \theta_{E_{F}, E_{G}}(x)=\angle\left(E_{F}(x), E_{G}(x)\right) .
\end{align*}
$$

Lemma B.6. Assume that $\left\|\mathcal{X}_{F}\right\|_{C^{1}} \leq K$. Then

$$
\Delta_{F^{t}, G^{t}, E_{F}, E_{G}}(x) \leq e^{K}\left[\delta_{F^{t}, G^{t}}+C e^{K} \theta_{E_{F}, E_{G}}(x)\right]
$$

for any $x \in \mathcal{S}$ and $t \in[0,1]$, where $C>0$ is a constant which depends only on the Riemannian metric of $\mathcal{M}$.

Proof. Similarly to the proof of Lemma B.4, we can show that if $\left\|\mathcal{X}_{F}\right\|_{C^{1}} \leq K$ then for any $x \in \mathcal{S}$ and $t \in[0,1]$,

$$
e^{-K} \leq m\left(d_{x} F^{t}\right) \leq\left\|d_{x} F^{t}\right\| \leq e^{K}
$$

We have that

$$
\begin{aligned}
& \left|\left\|d_{x} G^{t}\left|E_{G}(x)\|-\| d_{x} F^{t}\right| E_{F}(x)\right\|\right| \\
& \quad \leq\left|\left\|d_{x} G^{t}\left|E_{G}(x)\|-\| d_{x} F^{t}\right| E_{G}(x)\right\|\right|+\left|\left\|d_{x} F^{t}\left|E_{G}(x)\|-\| d_{x} F^{t}\right| E_{F}(x)\right\|\right| \\
& \quad \leq\left\|d_{x} G^{t}-d_{x} F^{t}\right\|+\left\|d_{x} F^{t}\right\| \operatorname{dist}\left(E_{G}(x), E_{F}(x)\right) \\
& \quad \leq\left\|d_{x} G^{t}-d_{x} F^{t}\right\|+C\left\|d_{x} F^{t}\right\| \angle\left(E_{G}(x), E_{F}(x)\right)
\end{aligned}
$$

for some constant $C>0$ depending only on the Riemannian metric of $\mathcal{M}$. Dividing both sides of the inequality by $\left\|d_{x} F^{t} \mid E_{F}(x)\right\|$ and noting that $\left\|d_{x} F^{t} \mid E_{F}(x)\right\| \geq m\left(d_{x} F^{t}\right)$, we obtain that

$$
\begin{aligned}
\left|\frac{\left\|d_{x} G^{t} \mid E_{G}(x)\right\|}{\left\|d_{x} F^{t} \mid E_{F}(x)\right\|}-1\right| & \leq \frac{1}{m\left(d_{x} F^{t}\right)}\left[\left\|d_{x} G^{t}-d_{x} F^{t}\right\|+C\left\|d_{x} F^{t}\right\| \angle\left(E_{G}(x), E_{F}(x)\right)\right] \\
& \leq e^{K}\left[\delta_{F^{t}, G^{t}}+C e^{K} \theta_{E_{F}, E_{G}}(x)\right]
\end{aligned}
$$

Similarly, one can show that $\left|\left(m\left(d_{x} G^{t} \mid E_{G}(x)\right) / m\left(d_{x} F^{t} \mid E_{F}(x)\right)\right)-1\right|$ admits the same upper bound.

Lemma B.7. A flow $F^{t}$ is uniformly partially hyperbolic on a compact invariant subset $\Lambda \subset \mathcal{S}$ if and only if the time-1 map $F^{1} \mid \Lambda$ is uniformly partially hyperbolic.

Proof. See [9].
Lemma B.8. Suppose that $F^{t}$ is uniformly partially hyperbolic on a compact invariant subset $\Lambda \subset \mathcal{S}$. Pick numbers $0<\lambda<\widetilde{\lambda} \leq 1 \leq \widetilde{\mu}<\mu$ such that

$$
\begin{array}{cc}
\lambda \geq \lambda\left(F^{1}, \Lambda\right)=\sup _{x \in \Lambda}\left\|d_{x}^{s} F^{1}\right\|, \quad \tilde{\lambda} \leq \tilde{\lambda}\left(F^{1}, \Lambda\right)=\inf _{x \in \Lambda} m\left(d_{x}^{c} F^{1}\right), \\
\tilde{\mu} \geq \tilde{\mu}\left(F^{1}, \Lambda\right)=\sup _{x \in \Lambda}\left\|d_{x}^{c} F^{1}\right\|, \quad \mu \leq \mu\left(F^{1}, \Lambda\right)=\inf _{x \in \Lambda} m\left(d_{x}^{u} F^{1}\right),
\end{array}
$$

where $d_{x}^{\omega} F^{t}=d_{x} F^{t} \mid E_{F}^{\omega}(x), \omega=s, c$, u. Given $\Delta>0$, there is $\varepsilon=\varepsilon(\Delta, \lambda, \tilde{\lambda}, \tilde{\mu}, \mu)$ such that if $\left\|\mathcal{X}_{G}-\mathcal{X}_{F}\right\|_{C^{1}}<\varepsilon$ and $\mathcal{X}_{G}=\mathcal{X}_{F}$ on $\mathcal{S} \backslash \Lambda$, then $G^{t} \mid \Lambda$ is also a uniformly partially hyperbolic flow and

$$
\begin{equation*}
\Delta_{F^{t}, G^{t}}^{\omega}(x):=\Delta_{F^{t}, G^{t}, E_{F}^{\omega}, E_{G}^{\omega}}(x) \leq \Delta t, \quad \omega=s, c, u, x \in \Lambda, t \in[0.5,1] \tag{B.5}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
1-\Delta \leq \frac{\lambda\left(G^{1}, \Lambda\right)}{\lambda\left(F^{1}, \Lambda\right)}, \quad \frac{\tilde{\lambda}\left(G^{1}, \Lambda\right)}{\tilde{\lambda}\left(F^{1}, \Lambda\right)}, \quad \frac{\tilde{\mu}\left(G^{1}, \Lambda\right)}{\widetilde{\mu}\left(F^{1}, \Lambda\right)}, \quad \frac{\mu\left(G^{1}, \Lambda\right)}{\mu\left(F^{1}, \Lambda\right)} \leq 1+\Delta \tag{B.6}
\end{equation*}
$$

Proof. Consider the time-1 map $F^{1}$. By [13], there is $\varepsilon<\Delta e^{-K} / 4 M$ depending on $\Delta, \lambda, \tilde{\lambda}, \tilde{\mu}, \mu$ such that if $\left\|\mathcal{X}_{G}-\mathcal{X}_{F}\right\|_{C^{1}}<\varepsilon$ and $\mathcal{X}_{G}=\mathcal{X}_{F}$ on $\mathcal{S} \backslash \Lambda$, then $G^{1} \mid \Lambda$ is uniformly partially hyperbolic on $\Lambda$ with

$$
\begin{equation*}
\sup _{x \in \Lambda} \angle\left(E_{G^{1}}^{\omega}(x), E_{F^{1}}^{\omega}(x)\right)<\frac{\Delta}{4 C e^{2 K}} . \tag{B.7}
\end{equation*}
$$

By Lemma B.7, the flow $G^{t}$ is uniformly partially hyperbolic on $\Lambda$ with the same invariant distributions as its time-1 map $G^{1}$. Moreover, it follows from Lemmas B. 5 and B. 6 that

$$
\Delta_{F^{t}, G^{t}, E_{F}^{\omega}, E_{G}^{\omega}}(x) \leq \frac{\Delta t}{2}+\frac{\Delta}{4} \leq \Delta t, \quad \omega=s, c, u, x \in \Lambda, t \in[0.5,1] .
$$

In particular,

$$
\left\|d_{x}^{s} G^{1}\right\| \leq\left\|d_{x}^{s} F^{1}\right\|(1+\Delta) \leq \lambda(1+\Delta)
$$

and hence $\lambda\left(G^{1}, \Lambda\right) / \lambda\left(F^{1}, \Lambda\right) \leq 1+\Delta$. The other inequalities in (B.6) can be shown in a similar fashion.

We will now specify how to choose the sequence of numbers $\varepsilon_{n}$ in the theorem. First choose four sequences of numbers $0<\lambda_{n}<\widetilde{\lambda}_{n} \leq 1 \leq \widetilde{\mu}_{n}<\mu_{n}$ such that the following hold.
(1) $\lambda_{n} \geq \lambda\left(H^{1}, \overline{\mathcal{U}_{n}}\right), \tilde{\lambda}_{n} \leq \tilde{\lambda}\left(H^{1}, \overline{\mathcal{U}_{n}}\right), \tilde{\mu}_{n} \geq \widetilde{\mu}\left(H^{1}, \overline{\mathcal{U}_{n}}\right)$, and $\mu_{n} \leq \mu\left(H^{1}, \overline{\mathcal{U}_{n}}\right)$.
(2) $\lambda_{n}, \widetilde{\mu}_{n}$ are strictly increasing while $\tilde{\lambda}_{n}, \mu_{n}$ are strictly decreasing.

For all $x \in \mathcal{S}$, let

$$
\gamma(x)=\min \left\{\frac{\min \left\{1, m\left(d_{x}^{c} H^{1}\right)\right\}}{\left\|d_{x}^{s} H^{1}\right\|}, \frac{m\left(d_{x}^{u} H^{1}\right)}{\max \left\{1,\left\|d_{x}^{c} H^{1}\right\|\right\}}\right\},
$$

and choose a strictly decreasing sequence of numbers $\gamma_{n}$ such that

$$
\begin{equation*}
0<\gamma_{n} \leq \inf _{x \in \overline{\mathcal{U}_{n}}} \frac{\gamma(x)-1}{8} \tag{B.8}
\end{equation*}
$$

Now choose a sequence of positive numbers $\Delta_{n}$ such that

$$
\begin{gather*}
\max \left\{\frac{\tilde{\lambda}_{n+1}}{\tilde{\lambda}_{n}}, \frac{\mu_{n+1}}{\mu_{n}}\right\} \leq 1-\Delta_{n}<1+\Delta_{n} \leq \min \left\{\frac{\lambda_{n+1}}{\lambda_{n}}, \frac{\tilde{\mu}_{n+1}}{\widetilde{\mu}_{n}}\right\},  \tag{B.9}\\
\Delta_{n}<\frac{1}{2^{n+2}}, \quad \sum_{k=n}^{\infty} \Delta_{k}<\gamma_{n} . \tag{B.10}
\end{gather*}
$$

Finally, choose

$$
\varepsilon_{n}<\frac{1}{2} \min \left\{\frac{K}{2^{n+1}}, \varepsilon\left(\Delta_{n}, \lambda_{n}, \tilde{\lambda}_{n}, \tilde{\mu}_{n}, \mu_{n}\right)\right\},
$$

where $\varepsilon(\Delta, \lambda, \tilde{\lambda}, \tilde{\mu}, \mu)$ is given by Lemma B.8.
Proof of Theorem B.1. First we shall show that for every $n>0$ the map $h_{n}^{t}$ is uniformly partially hyperbolic on $\overline{\mathcal{U}}_{n}$. This is clearly true for $h_{0}^{t}$ and we will use induction assuming that $h_{k}^{t} \mid \overline{\mathcal{U}_{k}}$ for $k=1, \ldots, n-1$ are uniformly partially hyperbolic. By Lemma B.6, we obtain that

$$
1-\Delta_{k} \leq \frac{\lambda\left(h_{k}^{1}, \overline{\mathcal{U}_{k}}\right)}{\lambda\left(h_{k-1}^{1}, \overline{\mathcal{U}_{k}}\right)}, \quad \frac{\tilde{\lambda}\left(h_{k}^{1}, \overline{\mathcal{U}_{k}}\right)}{\tilde{\lambda}\left(h_{k-1}^{1}, \overline{\mathcal{U}_{k}}\right)}, \quad \frac{\tilde{\mu}\left(h_{k}^{1}, \overline{\mathcal{U}_{k}}\right)}{\widetilde{\mu}\left(h_{k-1}^{1}, \overline{\mathcal{U}_{k}}\right)}, \quad \frac{\mu\left(h_{k}^{1}, \overline{\mathcal{U}_{k}}\right)}{\mu\left(h_{k-1}^{1}, \overline{\mathcal{U}_{k}}\right)} \leq 1+\Delta_{k} .
$$

Note that

$$
\begin{aligned}
\lambda\left(h_{k}^{1}, \overline{\mathcal{U}_{k+1}}\right) & \leq \max \left\{\lambda\left(H^{1}, \overline{\mathcal{U}_{k+1}}\right), \lambda\left(h_{k}^{1}, \overline{\mathcal{U}_{k}}\right)\right\} \\
& \leq \max \left\{\lambda_{k+1}, \lambda\left(h_{k}^{1}, \overline{\mathcal{U}_{k}}\right)\right\} \\
& \leq \max \left\{\lambda_{k+1}, \lambda\left(h_{k-1}^{1}, \overline{\mathcal{U}_{k}}\right)\left(1+\Delta_{k}\right)\right\} .
\end{aligned}
$$

The fact that $\lambda\left(h_{0}^{1}, \overline{\mathcal{U}_{1}}\right) \leq \lambda_{1}$ and the choice of $\Delta_{n}$ in (B.9) guarantee that

$$
\lambda_{n}^{\prime}:=\lambda\left(h_{n-1}^{1}, \overline{\mathcal{U}_{n}}\right) \leq \lambda_{n} .
$$

Similarly, we have

$$
\begin{gathered}
\tilde{\lambda}_{n}^{\prime}:=\tilde{\lambda}\left(h_{n-1}^{1}, \overline{\mathcal{U}_{n}}\right) \geq \tilde{\lambda}_{n}, \quad \widetilde{\mu}_{n}^{\prime}:=\widetilde{\mu}\left(h_{n-1}^{1}, \overline{\mathcal{U}_{n}}\right) \leq \tilde{\mu}_{n}, \\
\mu_{n}^{\prime}:=\mu\left(h_{n-1}^{1}, \overline{\mathcal{U}_{n}}\right) \geq \mu_{n} .
\end{gathered}
$$

It follows that

$$
\varepsilon_{n} \leq \varepsilon\left(\Delta_{n}, \lambda_{n}, \tilde{\lambda}_{n}, \tilde{\mu}_{n}, \mu_{n}\right) \leq \varepsilon\left(\Delta_{n}, \lambda_{n}^{\prime}, \tilde{\lambda}_{n}^{\prime}, \tilde{\mu}_{n}^{\prime}, \mu_{n}^{\prime}\right) .
$$

Since $\left\|\mathcal{X}_{n}-\mathcal{X}_{n-1}\right\|_{C^{1}} \leq \varepsilon_{n}$, by Lemma B.8, we obtain that $h_{n}^{t} \mid \overline{\mathcal{U}_{n}}$ is uniformly partially hyperbolic.

Next we will show that $\mathcal{X}=\lim _{n \rightarrow \infty} \mathcal{X}_{n}$ exists and is smooth. In fact, $\left\{\mathcal{X}_{n}\right\}$ is a Cauchy sequence in the $C^{1}$ topology since, for any $n, m \in \mathbb{N}$,

$$
\left\|\mathcal{X}_{n+m}-\mathcal{X}_{n}\right\|_{C^{1}} \leq \sum_{l=1}^{m}\left\|\mathcal{X}_{n+l}-\mathcal{X}_{n+l-1}\right\|_{C^{1}} \leq \sum_{l=1}^{m} \varepsilon_{n+l} \leq \frac{K}{2^{n+1}}
$$

Hence $\mathcal{X}=\lim _{n \rightarrow \infty} \mathcal{X}_{n}$ exists and is $C^{1}$.
It remains to show that the flow $h^{t}$ generated by $\mathcal{X}$ is pointwise partially hyperbolic on $\mathcal{S}$. First we construct invariant distributions for $h^{t}$. Given $x \in \mathcal{S}$, we have

$$
\angle\left(E_{h_{n}}^{\omega}(x), E_{h_{n-1}}^{\omega}(x)\right) \leq \frac{\Delta_{n}}{4 C e^{2 K}}<\frac{1}{2^{n+4} C e^{2 K}}, \quad \omega=s, c, u
$$

Hence the sequence of subspaces $E_{h_{n}}^{\omega}(x)$ is Cauchy and converges to

$$
E_{h}^{\omega}(x)=\lim _{n \rightarrow \infty} E_{h_{n}}^{\omega}(x)
$$

which is clearly $d h^{t}$-invariant for all $t \in \mathbb{R}^{+}$.
Now we would like to estimate $\Delta_{h^{1}, H^{1}}^{\omega}(x)$. Fixing $x \in \mathcal{U}_{n} \backslash \overline{\mathcal{U}_{n-1}}$, we have

$$
\Delta_{h_{k}, h_{k-1}^{1}}^{\omega}(x) \begin{cases}=0, & k<n \\ \leq \Delta_{k}, & k \geq n\end{cases}
$$

Noting that

$$
\frac{\left\|d_{x}^{\omega} h_{l}^{1}\right\|}{\left\|d_{x}^{\omega} H^{1}\right\|}=\prod_{k=1}^{l} \frac{\left\|d_{x}^{\omega} h_{k}^{1}\right\|}{\left\|d_{x}^{\omega} h_{k-1}^{1}\right\|}, \quad \frac{m\left(d_{x}^{\omega} h_{l}^{1}\right)}{m\left(d_{x}^{\omega} H^{1}\right)}=\prod_{k=1}^{l} \frac{m\left(d_{x}^{\omega} h_{k}^{1}\right)}{m\left(d_{x}^{\omega} h_{k-1}^{1}\right)},
$$

and $\sum \Delta_{k}<1 / 4$, we obtain, by Lemma B.2,

$$
\Delta_{h_{l}^{1}, H^{1}}^{\omega}(x) \leq \prod_{k=1}^{l}\left(1+\Delta_{h_{k}^{1}, h_{k-1}^{1}}^{\omega}(x)\right)-1 \leq \prod_{k=n}^{\infty}\left(1+\Delta_{k}\right)-1 \leq 2 \sum_{k=n}^{\infty} \Delta_{k}
$$

Letting $l \rightarrow \infty$, we have

$$
\Delta_{h^{1}, H^{1}}^{\omega}(x) \leq 2 \sum_{k=n}^{\infty} \Delta_{k}, \quad \omega=s, c, u, x \in \mathcal{U}_{n} \backslash \overline{\mathcal{U}_{n-1}}
$$

Therefore,

$$
\begin{aligned}
\frac{\left\|d_{x}^{s} h^{1}\right\|}{\min \left\{1, m\left(d_{x}^{c} h^{1}\right)\right\}} & \leq \frac{1+2 \sum_{k=n}^{\infty} \Delta_{k}}{1-2 \sum_{k=n}^{\infty} \Delta_{k}} \frac{\left\|d_{x}^{s} H^{1}\right\|}{\min \left\{1, m\left(d_{x}^{c} H^{1}\right)\right\}} \\
& <\left(1+8 \gamma_{n}\right) \frac{\left\|d_{x}^{s} H^{1}\right\|}{\min \left\{1, m\left(d_{x}^{c} H^{1}\right)\right\}} \\
& \leq \gamma(x) \frac{\left\|d_{x}^{s} H^{1}\right\|}{\min \left\{1, m\left(d_{x}^{c} H^{1}\right)\right\}}<1 .
\end{aligned}
$$

Similarly, one can show that $m\left(d_{x}^{u} h^{1}\right)>\max \left\{1,\left\|d_{x}^{c} h^{1}\right\|\right\}$. It follows that $h^{1}$ is pointwise partially hyperbolic on $\mathcal{S}$, and so is the flow $h^{t}$ by definition.

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