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# FAMILIES OF INVARIANT MANIFOLDS CORRESPONDING TO NONZERO CHARACTERISTIC EXPONENTS 

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Ja. B. PESIN

> Abstract. A theorem on conditional stability is proved for a family of mappings of class $C^{1+\epsilon}$, satisfying a condition more general than Ljapunov regularity. Using this theorem, families of invariant manifolds are constructed for a diffeomorphism of a smooth manifold onto a set where at least one Liapunov characteristic exponent is nonzero. The property of absolute continuity is proved for these families.

> Bibliography: 10 titles.

## Introduction

0.1. In this paper we consider a cascade (that is, a dynamical system with discrete time) generated by a diffeomorphism $f$ of a smooth closed $n$-dimensional manifold $M$ and preserving a finite measure $\nu$ compatible with smoothness (that is, equivalent to the measure induced by some Riemannian metric). We fix in $M$ an auxiliary Riemannian metric: the corresponding scalar product and norm on the tangent space $T_{x} M$ will be denoted $\langle,\rangle_{x}$ and $\|\cdot\|_{x}$ (sometimes the index $x$ will be omitted). The required class of smoothness of $f$ will be specified in the statements of theorems, but the smoothness of $M$ and the Riemannian metric can, without loss of generality, be taken to be class $C^{\infty}$. There is defined on the tangent bundle $T M$ a measurable function (see [7])

$$
\begin{equation*}
\chi^{+}(x, v)=\varlimsup_{n \rightarrow \infty} \frac{1}{n} \ln \left\|d f^{n} v\right\|, \quad v \in T_{x} M, \tag{0.1}
\end{equation*}
$$

called the Ljapunov characteristic exponent (the number $\chi^{+}(x, v)$ is called the characteristic exponent of the vector $v$ at the point $x$ ).

Our basic assumption is that the measurable invariant (relative to $f$ ) set

$$
\begin{equation*}
\Lambda=\left\{x \in M: \chi^{+}(x, v) \neq 0 \quad \text { for some } \quad v \in T_{x} M\right\} \tag{0.2}
\end{equation*}
$$

has positive measure.
If we interpret the existence of a vector $v \in T_{x} M$ with a nonzero exponent as a certain "partial hyperbolicity" property at $x$, then we may say that on $\Lambda$ our cascade is "nonuniformly partially hyperbolic". "Partially" here means that the existence of vectors with zero exponent is not excluded, and "nonuniformly" means that the inequality, expressing for vec-
tors $v \in T_{x} M$ with negative exponents the change of $\left\|d f^{n} v\right\|_{f^{n}(x)}$ with increasing $n$, is not uniform in $x$. For comparison we note that the case when hyperbolicity is not partial but is, as we say, "complete" (no vectors with zero exponents) and is "uniform", is the well known case of U-systems [1], [2].

First we consider an individual trajectory $\left\{f^{n}\left(x_{0}\right)\right\}$, where $x_{0} \in \Lambda$, and the mapping along it

$$
\begin{equation*}
d f^{n}: T_{x_{0}} M \rightarrow T_{f^{n}\left(x_{0}\right)} M \tag{0.3}
\end{equation*}
$$

(in the case of continuous time there are corresponding variational equations along the trajectory $x(t)$ ). Since in a linear approximation the behavior (relative to $\left\{f^{n}\left(x_{0}\right)\right\}$ ) of neighboring trajectories is described by these mappings, there then arises the question: does some "stability condition" of the trajectory $\left\{f^{n}\left(x_{0}\right)\right\}$ follow from our weak variant of "hyperbolicity", namely, the existence of a smooth submanifold $V\left(x_{0}\right) \subset M$ such that the trajectories $\left\{f^{n}(x)\right\}$ with initial value $x \in V\left(x_{0}\right)$ become arbitrarily near to $\left\{f^{n}\left(x_{0}\right)\right\}$. It is known that the answer, in general, is negative.

The similar circumstances for the problem of stability of the trivial solution $x(t) \equiv 0$ of a system of differential equations

$$
\begin{equation*}
\frac{d x}{d t}=A(t) x+f(t, x) \tag{0.4}
\end{equation*}
$$

where $x$ and $f$ are vectors, $A$ is a matrix, uniformly bounded and continuous, and $\|f\|=$ $o(\|x\|)$ (or even $O\left(\|x\|^{2}\right)$ ) uniformly in $t$, are more widely known. Namely, it is known that negativity of all Ljapunov characteristic exponents of the "linear approximation" $\dot{x}=A(t) x$ is not sufficient for Ljapunov stability of the solution $x(t)=0$ of the system ( 0.4 ) (see, for example, [4]). There exist different sufficient conditions providing an affirmative answer to the question raised.

The problem of interest to us is concerned with the case of a stability condition. In the analytic case, Ljapunov proved such a theorem for a system of differential equations satisfying a regularity condition. In this paper we prove a similar theorem for cascades of class $C^{1+\epsilon}$ and those points $x_{0} \in \Lambda$ at which the linear approximation (0.3) satisfies a regularity condition or a certain general condition (see Theorems 2.1.1 and 2.1.2). In addition we are interested in the invariant manifolds $V\left(x_{0}\right)$ not for a single point, but for a whole set of such points having positive measure. Accordingly we obtain some information on the dependence of $V\left(x_{0}\right)$ on $x_{0}$.

This result is intended for use in the metric theory, playing for the smooth cascades considered here the same role as the well-known Hadamard-Perron theorem [6], [8] (more precisely, the version uniform relative to initial data [1]) plays for U-systems.

In metric theory it is appropriate to deal not with one individual trajectory but with a set of trajectories of positive measure; therefore it is necessary to consider not the separate invariant manifolds but families of such manifolds. For U-systems, Anosov showed that these families have a very important special property which one calls absolute continuity (see [1]). (In essence, it is just this property that allows the use of invariant manifolds in the metric theory.) In our case we will prove an analogous result (see §3).

The uniform hyperbolicity of U-cascades gives them rich metric properties: they are ergodic, have positive entropy and even have the $K$-property. In our case the hyperbolicity is weaker, but some analogue of these results is retained [9] (at least in the case when in $\Lambda$ all the exponents are different from zero; however, we will give proofs in the general case because, firstly, they do not follow from [9], and, secondly, the assertions proved are used to get a new result-a formula for the entropy of a dynamical system; see [9]). These proofs are based on the above-mentioned properties of invariant manifolds.

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0.2. In this paper we will constantly use the following notation in addition to that introduced above:
$\langle$,$\rangle and \|\cdot\|$ are the standard scalar product and corresponding norm in Euclidean space $\mathbf{R}^{n}$.
$I$ is the identity mapping of the corresponding space (which one will be clear from the context).
$\mathbf{Z}^{+}$and $\mathbf{R}^{+}$are the sets of nonnegative integers or reals.
$\rho$ and $d$ are the distances in $M$ and $T M$ respectively, induced by the Riemannian metric.
$B(x, r)$ is the open ball in the manifold $M$ with center at $x$ and radius $r$.
$T^{*} M$ is the cotangent bundle; $d^{\prime} f=\left((d f)^{*}\right)^{-1}: T^{*} M \rightarrow T^{*} M$.
0.3. In conclusion we recall some information about the Ljapunov characteristic exponents (see [4] and [7]).

The upper index "plus" in the notation of exponents means that they are obtained as "time" $n \rightarrow+\infty$ (see (0.1)). Similarly we define the exponent $\chi^{-}$as $n \rightarrow-\infty$. Correspondingly we should speak of "forward" and "backward" exponents, although, as a rule, working with the first we will omit the "forward".

For each $x \in M$ the restriction of $\chi^{+}$to the subspace $T_{x} M$ takes not more than $n$ values (distinct and different from $-\infty$ ). We denote these values in order by

$$
\begin{equation*}
\chi_{1}(x)<\chi_{2}(x)<\ldots<\chi_{s(x)}(x), \quad s(x) \leqslant n . \tag{0.5}
\end{equation*}
$$

Put $L_{i}(x)=\left\{v \in T_{x} M: \chi(x, v) \leqslant \chi_{i}(x)\right\}$. The subspaces $L_{i}(x)$ form a filtration of $T_{x} M$; that is,

$$
\begin{equation*}
0=L_{0}(x) \subset L_{1}(x) \subset \ldots \subset L_{s(x)}(x)=T_{x} M \tag{0.6}
\end{equation*}
$$

Put $\operatorname{dim} L_{i}(x)=k_{i}(x)$. The integer-valued functions $s(x), k_{1}(x), \ldots, k_{s(x)}(x)$ and the family of subspaces $L_{i}(x), i=1, \ldots, s(x)$, depend measurably on $x$. The characteristic exponent $\chi^{+}$is invariant; that is, for each $x \in M$

$$
\chi_{i}(x)=\chi_{i}(f(x)), \quad k_{i}(x)=k_{i}(f(x)), \quad L_{i}(x)=L_{i}(f(x))
$$

If we replace $d f$ by $d^{\prime} f$ on the right-hand side of ( 0.1 ), we obtain a characteristic exponent $\chi^{\prime+}$ on the cotangent bundle, called an adjoint exponent.

Let $x \in M, v \in T_{x} M$ and $\phi \in T_{x}^{*} M$. Then

$$
\begin{equation*}
a^{\prime} f(\varphi(d f v))=\varphi(v)=1 . \tag{0.7}
\end{equation*}
$$

From this follows the so-called adjoint condition: $\chi^{+}(x, v)^{+}+\chi^{\prime}(x, \phi) \geqslant 0$. A normalized basis $e(x)=\left\{e_{i}(x)\right\} \in T_{x} M$ is called normal if the first $k_{1}(x)$ vectors of the basis lie in $L_{1}(x)$, the following $k_{2}(x)-k_{1}(x)$ lie in $L_{2}(x) \backslash L_{1}(x)$, etc. If $e(x)=\left\{e_{i}(x)\right\}$ is a normal basis at the point $x$, then

$$
e^{n}(x)=\left\{\frac{d f^{n} e_{i}(x)}{\left\|d f^{n} e_{i}(x)\right\|}\right\}
$$

is a normal basis at $f^{n}(x)$.
Let $e^{\prime}(x)=\left\{e_{i}^{\prime}(x)\right\} \in T_{x}^{*} M$ be the dual basis. The defect of the pair of bases is the function

$$
\gamma\left(x, e(x), e^{\prime}(x)\right)=\max _{i}\left\{\chi^{+}\left(x, e_{i}(x)\right)+\chi^{\prime+}\left(x, e_{i}^{\prime}(x)\right)\right\}
$$

By virtue of the adjoint condition, $\gamma\left(x, e(x), e^{\prime}(x)\right) \geqslant 0$. The irregularity coefficient of the exponent $\chi^{+}$is the function

$$
\gamma(x)=\min \gamma\left(x, e(x), e^{\prime}(x)\right),
$$

where the minimum is taken over all possible pairs of dual bases $\left(e(x), e^{\prime}(x)\right)$. The exponent $\chi^{+}$(or the pair of exponents $\chi^{+}$and $\chi^{\prime+}$ ) is called regular at $x$ if $\gamma(x)=0$, and points at which this condition is satisfied are called forward regular.

Theorem 0.1 (see [4]). Let $x$ be a forward regular point. Then the following assertions are true:

1) $\chi_{i}(x)=-\chi_{i}^{\prime}(x)$.
2) The filtration connected with $\chi^{\prime+}$ consists of the subspaces $L_{i}^{\perp}(x)$ (the annihilator of the subspace $L_{i}(x)$ ).
3) Each basis dual to a normal basis is normal.
4) For each $n \in Z$ the point $f^{n}(x)$ is forward regular.

A point $x$ is called backward regular if it is forward regular for the exponent $\chi^{-}$. A point $x$ is called regular if it is both forward and backward regular. We have detailed the notion of regularity, in contrast to [4] and [7] (their notion of regularity corresponds to our "forward regularity").

If $\phi$ is a measurable invariant function on the manifold $M$, then its Ljapunov characteristic exponent at $x$ is

$$
\begin{equation*}
\chi^{+}(\varphi(x))=\varlimsup_{n \rightarrow+\infty} \frac{1}{n} \ln \left|\varphi\left(f^{n}(x)\right)\right| . \tag{0.8}
\end{equation*}
$$

We similarly define the exponent $\chi^{-}(\phi(x))$ for $n \rightarrow-\infty$.
Theorem 0.2 (see [7], Theorem 4). If $x$ is regular, then there exist subspaces $E_{i}(x)$, $i=1, \ldots, s(x)$, satisfying the following conditions:

1) $L_{i}(x)=\bigoplus_{j=1}^{k_{i}(x)} E_{j}(x), i=1, \ldots, s(x)$.
2) Uniformly in $v \in E_{j}(x)$

$$
\lim _{n= \pm \infty} \ln \frac{1}{n}\left\|d f^{n} v\right\|= \pm \chi_{j}(x)
$$

3) $\chi^{-}\left(\Gamma_{j}(x)\right)=\chi^{+}\left(\Gamma_{j}(x)\right)=\left(k_{j}(x)-k_{j-1}(x)\right) \chi_{j}(x)$, where $\Gamma_{j}(x)$ is the volume of the parallelepiped in the space $E_{j}(x)$.
4) For each $n \in \mathbf{Z}$ the point $f^{n}(x)$ is regular, and df ${ }^{n} E_{j}(x)=E_{j}\left(f^{n}(x)\right)$.
5) There is a decomposition $T_{x}^{*} M=\bigoplus_{j=1}^{s(x)} E_{j}^{\prime}(x)$ with similar properties relative to the exponent $\chi^{\prime+}$; moreover, if $e(x)=\left\{e_{i}(x)\right\}$ is a normal basis for which $e_{i}(x) \in E_{j}(x), k_{j-1}(x)$ $<i \leqslant k_{j}(x)$, and if $e^{\prime}(x)=\left\{e_{i}^{\prime}(x)\right\}$ is the dual basis, then $e_{j}^{\prime}(x) \in E_{j}^{\prime}(x), k_{j-1}(x)<i \leqslant k_{j}(x)$.

Theorem 0.3 (see [7], Theorem 1). Let $f$ be a dynamical system on the smooth manifold $M$, preserving a finite Borel measure. Then relative to this measure almost any point $x \in M$ is regular. The function $s(x)$ and the subspaces $E_{1}(x), \ldots, E_{s(x)}(x)$ depend measurably on $x$.

## §1. Invariant distribution and the Ljapunov metric

1.1. Let $f$ be a diffeomorphism of $M$ of class $C^{r}, r \geqslant 2$, preserving a measure $\nu$, and let the measurable invariant set $\Lambda$ defined by (0.2) have positive measure. We denote by $\widetilde{\Lambda}$ the set of regular points in $\Lambda$. From Theorem 0.3 it follows that $\nu(\widetilde{\Lambda})=\nu(\Lambda)$.

Let $x \in \widetilde{\Lambda}$. We consider the filtration (0.6) and the subspaces $E_{j}(x), j=1, \ldots, s(x)$, at the point $x$ (see Theorem 0.2). We denote by $k(x)$ the largest natural number such that for each $v \in L_{k(x)}(x)$

$$
\chi^{+}(x, v)<0 .
$$

According to the definition ( 0.2 ), $1 \leqslant k(x)<s(x)$ and $k(f(x))=k(x)$.
Let $l(x)$ be an arbitrary measurable invariant function satisfying the condition $l \leqslant l(x)$ $\leqslant k(x)$. Put

$$
\begin{array}{cc}
E_{1 x}=\bigoplus_{j=1}^{l(x)} E_{j}(x), & E_{2 x}=\bigoplus_{j=l(x)+1}^{s(x)} E_{j}(x),  \tag{1.1.1}\\
\lambda(x)=e^{\chi_{l(x)}(x)}, & \mu(x)=e^{x_{l(x)+1}(x)} .
\end{array}
$$

From Theorems 0.3 and 0.2 it follows that the functions $\lambda(x)$ and $\mu(x)$ and the subspaces $E_{1 x}$ and $E_{2 x}$ depend measurably on $x$ and satisfy the following conditions:

$$
\begin{gather*}
0<\lambda(x)<\mu(x), \quad \lambda(x)<1, \quad \lambda(f(x))=\lambda(x), \quad \mu(f(x))=\mu(x)  \tag{1.1.2}\\
T_{x} M=E_{1 x} \oplus E_{2 x}, \quad d f E_{i x}=E_{i f(x)}, \quad i=1,2 \tag{1.1.3}
\end{gather*}
$$

Theorem 1.1.1. There exist measurable functions $C(x, \epsilon)$ and $K(x, \epsilon), \epsilon>0, x \in \widetilde{\Lambda}$, satisfying the following conditions:

1) For every $m \in \mathbf{Z}$

$$
\begin{align*}
& C\left(f^{m}(x), \varepsilon\right) \leqslant C(x, \varepsilon) e^{\varepsilon \varepsilon|m|}  \tag{1.1.4}\\
& K\left(f^{m}(x), \varepsilon\right) \geqslant K(x, \varepsilon) e^{-\varepsilon|m|}
\end{align*}
$$

2) For any $n \in \mathbf{Z}^{+}$

$$
\begin{array}{ll}
\left\|d f^{n} v\right\| \leqslant C(x, \varepsilon) \lambda^{n}(x) e^{\varepsilon n}\|v\|, & \left(v \in E_{1 x}\right),  \tag{1.1.5}\\
\left\|d f^{-n} v\right\| \geqslant C^{-1}(x, \varepsilon) \lambda^{-n}(x) e^{-\varepsilon n}\|v\| & \\
\left\|d f^{n} v\right\| \geqslant C^{-1}(x, \varepsilon) \mu^{n}(x) e^{-\varepsilon n}\|v\|, \quad\left(v \in E_{2 x}\right) . \\
\left\|d f^{-n} v\right\| \leqslant C(x, \varepsilon) \mu^{-n}(x) e^{\varepsilon n}\|v\| &
\end{array}
$$

3. Let $\gamma(x)$ be the angle between the subspaces $E_{1 x}$ and $E_{2 x}$. Then $K(x, \epsilon) \leqslant \gamma(x)$.

Proof. We begin by proving
Lemma 1.1.1. Let $X \subset M$ be a measurable set invariant relative to $f$, and let $A(x, \epsilon)$ be a measurable function on $X$ which for some $\epsilon>0$ and all $m \in \mathbf{Z}$ satisfies the condition

$$
\begin{equation*}
M_{1}(x, \varepsilon) e^{-s \varepsilon|m|} \leqslant A\left(f^{m}(x), \varepsilon\right) \leqslant M_{2}(x, \varepsilon) e^{3 \varepsilon|m|} \tag{1.1.6}
\end{equation*}
$$

where $M_{1}(x, \epsilon)$ and $M_{2}(x, \epsilon)$ are measurable functions. Then there are measurable functions $B_{1}(x, \epsilon)$ and $B_{2}(x, \epsilon)$ such that for any $m \in \mathbf{Z}$

$$
\begin{align*}
& B_{1}(x, \varepsilon) \leqslant A(x, \varepsilon) \leqslant B_{2}(x, \varepsilon)  \tag{1.1.7}\\
& B_{1}(x, \varepsilon) e^{-4 \varepsilon|m|} \leqslant B_{1}\left(f^{m}(x), \varepsilon\right)  \tag{1.1.8}\\
& B_{2}(x, \varepsilon) e^{\varepsilon|m|} \geqslant B_{2}\left(f^{m}(x), \varepsilon\right)
\end{align*}
$$

Proof. (1.1.6) implies the existence of a natural number $m(x, \epsilon)$ such that for $m \in \mathbf{Z}$, $|m| \geqslant m(x, \epsilon)$,

$$
-4 \varepsilon \leqslant \frac{1}{|m|} \ln A\left(f^{m}(x), \varepsilon\right) \leqslant 4 \varepsilon
$$

An

$$
\begin{aligned}
& B_{1}(x, \varepsilon)=\min _{-m(x, \varepsilon) \leqslant i \leqslant m(x, \varepsilon)}\left\{1, A\left(f^{i}(x), \varepsilon\right) e^{\iota \varepsilon[i]}\right\} ; \\
& B_{2}(x, \varepsilon)=\max _{-m(x, \varepsilon) \leqslant i \leqslant m(x, \varepsilon)}\left\{1, A\left(f^{i}(x), \varepsilon\right) e^{-\varepsilon \varepsilon[i]}\right\}
\end{aligned}
$$

The functions $B_{j}(x, \epsilon), j=1,2$, are measurable. In addition, for each $n \in \mathbf{Z}$

$$
\begin{equation*}
B_{1}(x, \varepsilon) e^{-\varepsilon \varepsilon|n|} \leqslant A\left(f^{n}(x), \varepsilon\right) \leqslant B_{2}(x, \varepsilon) e^{\alpha[n]} \tag{1.1.9}
\end{equation*}
$$

If the numbers $b_{1} \leqslant 1 \leqslant b_{2}$ have the property that for each $n \in \mathbf{Z}$

$$
\begin{gather*}
b_{1} e^{-\varepsilon \varepsilon|n|} \leqslant A\left(f^{n}(x), \varepsilon\right)  \tag{1.1.10}\\
b_{2} e^{\langle\varepsilon| n]} \geqslant A\left(f^{n}(x), \varepsilon\right) \tag{1.1.11}
\end{gather*}
$$

then $b_{1} \leqslant B_{1}(x, \epsilon)$ and $b_{2} \geqslant B_{2}(x, \epsilon)$. Thus

$$
\begin{align*}
& B_{1}(x, \varepsilon)=\sup \{b \leqslant 1: \text { for each } n(1.1 .10) \text { is satisfied }\},  \tag{1.1.12}\\
& B_{\mathbf{2}}(x, \varepsilon)=\inf \{b \geqslant 1: \text { for each } n(1.1 .11) \text { is satisfied }\} .
\end{align*}
$$

The inequality (1.1.9) implies (1.1.7), and also the inequalities

$$
\begin{aligned}
& A\left(f^{n+m}(x), \varepsilon\right) \leqslant B_{2}(x, \varepsilon) e^{4 \varepsilon|n+m|} \leqslant B_{2}(x, \varepsilon) e^{4 \varepsilon|n|+4 \varepsilon|m|} \\
& A\left(f^{n+m}(x), \varepsilon\right) \geqslant B_{1}(x, \varepsilon) e^{-4 \varepsilon|n+m|} \geqslant B_{1}(x, \varepsilon) e^{-4 \varepsilon|n|-4 \varepsilon|m|}
\end{aligned}
$$

Comparing these two inequalities with (1.1.9), written at the point $f^{m}(x)$, and taking account of (1.1.12), we obtain (1.1.8). The lemma is proved.

Remark. If a measurable function $\psi(x)$ has exact exponents (positive and negative), equal to zero, then the function $A(x, \epsilon)=\psi(x)$ satisfies the condition of Lemma 1.1.1 for any $\epsilon>0$.

Fix $\epsilon>0$ and consider the angle $\gamma(x)$ between the subspaces $E_{1 x}$ and $E_{2 x}$. By Theorem 0.2 , for each $x \in \widetilde{\Lambda}$

$$
\begin{equation*}
\chi^{+}(\gamma(x))=\chi^{-}(\gamma(x))=0 ; \tag{1.1.13}
\end{equation*}
$$

moreover, the exponents are exact. Applying Lemma 1.1.1, we construct a function $K(x, \epsilon)$ satisfying the second inequality in (1.1.4) and assertion 3 of the theorem.

Let $x \in \widetilde{\Lambda}$. We denote $d f_{j x}=d f_{x} l_{E_{j}(x)}, j=1, \ldots, s(x)$.
Lemma 1.1.2. For any $\epsilon>0$ there is a function $D(x, \epsilon), x \in \widetilde{\Lambda}$, satisfying for any $m \in \mathbf{Z}$ and $1 \leqslant j \leqslant l(x)$ the condition

$$
\begin{equation*}
D\left(f^{m}(x), \varepsilon\right) \leqslant D^{2}(x, \varepsilon) e^{2 \varepsilon|m|} \tag{1.1.14}
\end{equation*}
$$

such that for each $n \in \mathbf{Z}^{+}$

$$
\begin{gather*}
\left\|d f_{j x}^{n}\right\| \leqslant D(x, \varepsilon) e^{\left(x_{j}+\varepsilon\right) n},  \tag{1.1.15}\\
\left\|d f_{j x}^{-n}\right\| \geqslant D^{-1}(x, \varepsilon) e^{-\left(x_{j}+\varepsilon\right) n},
\end{gather*}
$$

where $\chi_{j}=\chi_{j}(x)($ see (0.5)).
Proof. By Theorem 0.2 there is a number $n(x, \epsilon)>0$ such that for each $n \geqslant n(x, \epsilon)$

$$
\begin{aligned}
& \chi_{j}-\varepsilon \leqslant \frac{1}{n} \ln \left\|d f_{j x}^{n}\right\| \leqslant \chi_{i}+\varepsilon, \quad-\chi_{j}-\varepsilon \leqslant \frac{1}{n} \ln \left\|d f_{i x}^{-n}\right\| \leqslant-\chi_{j}+\varepsilon, \\
& -\chi_{j}-\varepsilon \leqslant \frac{1}{n} \ln \left\|d^{\prime} f_{j x}^{n}\right\| \leqslant-\chi_{j}+\varepsilon, \quad \chi_{j}-\varepsilon \leqslant \frac{1}{n} \ln \left\|d^{\prime} f_{j x}^{-n}\right\| \leqslant \chi_{j}+\varepsilon,
\end{aligned}
$$

where we have put $d^{\prime} f_{i x}=\left.d^{\prime} f\right|_{E_{j}^{\prime}(x)}\left(E_{j}^{\prime}(x)\right.$ are the subspaces of $T_{x}^{* M}$ occurring in Theorem 0.2 ).

Define

$$
\begin{gathered}
D_{1}(x, \varepsilon)=\max _{\substack{-n(x, \varepsilon \leqslant i \leqslant n(x, \varepsilon) \\
1 \leqslant j \leqslant l(x)}}\left\{1,\left\|d f_{j x}^{i}\right\| e^{-\left((\operatorname{sign} i) x_{j}+\varepsilon\right)|i|},\left\|d^{\prime} f_{j x}^{l}\right\| e^{-\left(-(\operatorname{sign} i) x_{j}+\varepsilon\right)|d|}\right\}, \\
D_{2}(x, \varepsilon)=\min _{\substack{n(x, \varepsilon) \leqslant i \leqslant n(x, \varepsilon) \\
1 \leqslant j \leqslant l(x)}}\left\{1,\left\|d f_{j x}^{i}\right\| e^{-\left((\operatorname{sign} i) x_{j}-\varepsilon\right) \mid(i)},\left\|d^{\prime} f_{j x}^{i}\right\| e^{-\left(-(\operatorname{sign} i) x_{j}-\varepsilon\right)|i|}\right\}, \\
D(x, \varepsilon)=\max \left\{D_{1}(x, \varepsilon), D_{2}^{-1}(x, \varepsilon)\right\} .
\end{gathered}
$$

The function $D(x, \varepsilon)$ is measurable, and for any $n \in \mathbf{Z}^{+}$and $l \leqslant j \leqslant l(x)$

$$
\begin{gather*}
D^{-1}(x, \varepsilon) e^{\left(x_{j}-\varepsilon\right) n} \leqslant\left\|d f_{j x}^{n}\right\| \leqslant D(x, \varepsilon) e^{\left(x_{j}+\varepsilon\right) n}, \\
D^{-1}(x, \varepsilon) e^{\left(-x_{j}-\varepsilon\right) n} \leqslant\left\|d f_{j x}^{n}\right\| \leqslant D(x, \varepsilon) e^{\left(-x_{j} f \varepsilon\right) n}, \\
D^{-1}(x, \varepsilon) e^{\left(-x_{j}-\varepsilon\right) n} \leqslant\left\|d^{\prime} f_{j x}^{n}\right\| \leqslant D(x, \varepsilon) e^{\left(-x_{j}+\varepsilon\right) n},  \tag{1.1.16}\\
D^{-1}(x, \varepsilon) e^{\left(x_{j}-\varepsilon\right) n} \leqslant\left\|d^{\prime} f_{j x}^{-n}\right\| \leqslant D(x, \varepsilon) e^{\left(x_{j}+\varepsilon\right) n}
\end{gather*}
$$

If a number $d \geqslant 1$ has the property that for any $n \in \mathbf{Z}^{+}$(1.1.16) are satisfied on replacing $D(x, \epsilon)$ by $d$, then $d \geqslant D(x, \epsilon)$, Thus

$$
\begin{align*}
D(x, \epsilon)=\inf \{d \geqslant 1: & \text { for each } n>0(1.1 .16) \text { is satisfied }  \tag{1.1.17}\\
& \text { on replacing } D(x, \epsilon) \text { by } d\} .
\end{align*}
$$

We will compare the values of $D(x, \epsilon)$ at the points $x$ and $f^{m} x, m \in \mathbf{Z}^{+}$. We make some preliminary remarks. We define a mapping $\tau_{x}: T_{x}^{*} M \rightarrow T_{x} M$, putting $\tau_{x}(\phi)=v_{\phi}$ for $\phi \in T_{x}^{*} M$, where the vector $v_{\phi}$ satisfies

$$
\left\langle v_{\varphi}, v\right\rangle=\varphi(v)
$$

for each $v \in T_{x} M$. Let $\left\{e_{i}^{n}\right\}$ be an orthonormal basis in $E_{j}\left(f^{n}(x)\right),\left\{e_{i}^{\prime n}\right\}$ the dual basis in $E_{j}^{\prime}\left(f^{n}(x)\right)$. Then $v_{e_{i}^{\prime} n}=e_{i}^{n}$. Let the mappings $d f_{j f}^{n} m_{(x)}$ and $\tau_{f^{n+m}(x)} d^{\prime} f_{j f}^{n}{ }^{m}(x) \tau_{f^{m}(x)}^{-1}$ relative to the basis $\left\{e_{i}^{e_{i}}\right\}$ be given by matrices $A_{n, m}^{j}$ and $B_{n, m}^{j}$ respectively. Condition (0.7) written in our notation gives

$$
A_{0, m}^{j}\left(B_{0, m}^{j}\right)^{\top}=I
$$

where ${ }^{\top}$ denotes the transposed matrix. Therefore for each $n>0$ the matrix of $d f_{f^{m}(x)}^{n}$ has the form

$$
A_{0, n+m}^{j}\left(A_{0, m}^{j}\right)^{-1}=A_{0, n+m}^{j}\left(B_{0, m}^{j}\right)^{\top}
$$

We will use the inequalities (1.1.16) in order to estimate the norm of the operator $d f_{j f_{(x)}}^{n}$ for different $n$. We consider the following cases:

1. $n>0$.

$$
\begin{aligned}
&\left\|d f_{j f^{m}(x)}^{n}\right\| \leqslant D^{*}(x, \varepsilon) e^{\left(x_{j}+\varepsilon\right)(n+m)+\left(-x_{j}+\varepsilon\right) m}=D^{2}(x, \varepsilon) e^{2 \varepsilon m} e^{\left(x_{j}+\varepsilon\right) n}, \\
& \| d f^{n} \\
& \text { 2. } n>0, m-n>0
\end{aligned}
$$

$$
\begin{gathered}
\left\|d f_{j f^{m}(x)}^{n}\right\| \leqslant D^{2}(x, \varepsilon) e^{\left(x_{j}+\varepsilon\right)(m-n)+\left(-x_{i}+\varepsilon\right) m} \leqslant D^{2}(x, \varepsilon) e^{2 \varepsilon m} e^{\left(-x_{j}+\varepsilon\right) n} \\
\left\|d f_{j f^{m}(x)}^{-n}\right\| \geqslant D^{-2}(x, \varepsilon) e^{\left(x_{j}-\varepsilon\right)(m-n)+\left(-x_{j}-\varepsilon\right) m} \geqslant D^{-2}(x, \varepsilon) e^{-2 \varepsilon m^{\left(-x_{j}-\varepsilon\right) n}} e^{(2)}
\end{gathered}
$$

3. $n>0, n-m>0$.

$$
\begin{aligned}
& \left\|d f_{\| f f^{m}(x)}^{n}\right\| \leqslant D^{2}(x, \varepsilon) e^{\left(-x_{j}+\varepsilon\right)(n-m)+\left(-x_{j}+\varepsilon\right) m}=D^{2}(x, \varepsilon) e^{2 \varepsilon m} e^{\left(-x_{j}+\varepsilon\right) n} \\
& \left\|d f_{j f^{m}(x)}^{n}\right\| \geqslant D^{-2}(x, \varepsilon) e^{\left(-x_{j}-\varepsilon\right)(n-m)+\left(-x_{j}-\varepsilon\right) m}=D^{-2}(x, \varepsilon) e^{-2 \varepsilon m} e^{\left(-x_{j}-\varepsilon\right) n}
\end{aligned}
$$

Similar inequalities hold for $d^{\prime} f_{i f}^{n}{ }_{(x)}, n, m \in \mathbf{Z}^{+}$. Comparing the relations obtained and (1.1.16), written out for $f^{m}(x)$, and also taking account of (1.1.17), we get that (1.1.14) holds for any $m \in \mathbf{Z}$. Arguing similarly, it is not difficult to show that this inequality also holds for each integer $m \leqslant 0$. The lemma is proved.

If we replace $m$ by $-m$ and $x$ by $f^{m}(x)$ in (1.1.14), we obtain

$$
\begin{equation*}
D\left(f^{m}(x), \varepsilon\right) \geqslant \sqrt{D(x, \varepsilon)} e^{-2 \varepsilon|m|} \tag{1.1.18}
\end{equation*}
$$

Let $x \in \tilde{\Lambda}$. We consider two disjoint subsets $\sigma_{1}$ and $\sigma_{2}$ of the set of natural numbers from 1 to $l(x)$. Put

$$
L(x)=\bigoplus_{i \in \sigma_{1}} E_{i}(x), \quad M(x)=\bigoplus_{i \in \sigma_{1}} E_{i}(x)
$$

and consider the angle $\gamma_{\sigma_{1} \sigma_{2}}(x)$ between the subspaces $L(x)$ and $M(x)$. By Theorem 0.1, (1.1.13) is satisfied for each $x \in \widetilde{\Lambda}$. Therefore, by Lemma 1.1.1 the function $\gamma_{\sigma_{1} \sigma_{2}}(x)$ satisfies the estimate

$$
\gamma_{\sigma_{1} \sigma_{1}}(x) \geqslant K_{\sigma_{1} \sigma_{2}}(x, \varepsilon)
$$

with some function $K_{\sigma_{1} \sigma_{2}}(x, \epsilon)$ satisfying the second inequality in (1.1.4). Put

$$
T(x, \varepsilon)=\min K_{\sigma_{1} \sigma_{2}}(x, \varepsilon)
$$

where the minimum is taken over all possible pairs of disjoint subsets $\sigma_{1}$ and $\sigma_{2}$ of the set of natural numbers from 1 to $l(x)$. The function $T(x, \epsilon)$ also satisfies the second inequality in (1.1.4).

Let $v \in E_{1 x}, v=\Sigma_{1}^{l(x)} v_{j}$, where $v_{j} \in E_{j}(x)$. From the above follows the existence of a constant $L>1$ such that

$$
\begin{equation*}
\sum_{j=1}^{l(x)}\left\|v_{j}\right\| \leqslant L T^{-1}(x, \varepsilon)\|v\| . \tag{1.1.19}
\end{equation*}
$$

Put

$$
C_{1}^{\prime}(x, \varepsilon)=L D(x, \varepsilon) T^{-1}(x, \varepsilon)
$$

From (1.1.4) and (1.1.18) it follows that $C_{1}^{\prime}(x, \epsilon)$ satisfies (1.1.6) of Lemma 1.1.1, where

$$
M_{1}(x, \varepsilon)=\frac{2}{\pi} L \sqrt{D(x, \varepsilon)}, \quad M_{2}(x, \varepsilon)=L D^{2}(x, \varepsilon) T^{-1}(x, \varepsilon)
$$

Therefore there exists a function $C_{1}(x, \epsilon) \geqslant C_{1}^{\prime}(x, \epsilon)$ satisfying the first inequality in (1.1.4) for all $m \in \mathbf{Z}$. In addition, by virtue of (1.1.15) and (1.1.19), the first two inequalities in (1.1.5) are satisfied with the function $C_{1}(x, \epsilon)$ for each $v \in E_{1 x}$ and $n \in \mathbf{Z}^{+}$.

Our construction is symmetric relative to the passage to the inverse mapping. Here the subspaces $E_{1 x}$ and $E_{2 x}$ exchange roles. Thus, repeating the preceding argument with the mapping $f^{-1}$ and the subspace $E_{2 x}$, we construct a measurable function $C_{2}(x, \epsilon)$, satisfying the first inequality in (1.1.4). In addition, for any $v \in E_{2 x}$ and $n \in \mathbf{Z}^{+}$the first two inequalities in (1.1.5) are satisfied with the function $C_{2}(x, \epsilon)$.

Put

$$
C(x, \varepsilon)=\max \left\{C_{1}(x, \varepsilon), C_{2}(x, \varepsilon)\right\}
$$

It is not difficult to verify that the measurable function $C(x, \epsilon)$ satisfies the first inequality in (1.1.4), and that the inequalities (1.1.5) are satisfied with $C(x, \epsilon)$. The theorem is proved.
1.2. It follows from Theorem 1.1.1 that it is possible to introduce a nonuniform partially hyperbolic structure in $\tilde{\Lambda}$. We will briefly describe the corresponding construction.

As above, let $k(x)$ denote the largest natural number such that $\chi^{+}(x, v)<0$ for each $v \in L_{k(x)}(x)$. Put

$$
\begin{gather*}
E_{1 x}=\bigoplus_{j=1}^{k(x)} E_{j}(x), \quad E_{2 x}=\underset{j=k(x)+2}{\stackrel{s(x)}{\bigoplus}} E_{j}(x)  \tag{1.2.1}\\
E_{0 x}=E_{k(x)+1}(x), \quad \lambda(x)=e^{x_{k(x)}(x)}, \quad \mu(x)=e^{\chi_{k(x)+2}(x)}
\end{gather*}
$$

It is not difficult to see that

$$
\begin{gather*}
T_{x} M=E_{1 x} \oplus E_{0 x} \oplus E_{2 x}, \quad d f E_{i x}=E_{i f(x)}, \quad i=1,2,0, \\
0<\lambda(x)<1<\mu(x), \quad \lambda(f(x))=\lambda(x), \quad \mu(f(x))=\mu(x) . \tag{1.2.2}
\end{gather*}
$$

In addition, $\chi_{k(x)+1}(x)=0$ and $\chi^{+}(x, v)>0$ for each $v \in E_{2 x}$.
The following is proved in the same way as Theorem 1.1.1.
Theorem 1.2.1. There exist measurable functions $C(x, \epsilon)$ and $K(x, \epsilon), \epsilon>0, x \in \widetilde{\Lambda}$, satisfying (1.1.4) and the following conditions:

1. For each $v \in E_{1 x}$ (or $v \in E_{2 x}$ ) the inequalities (1.1.5) hold.
2. For any $n \in \mathbf{Z}$ and $v \in E_{0 x}$

$$
C^{-1}(x, \varepsilon) e^{-\varepsilon|n|}\|v\| \leqslant\left\|d f^{n} v\right\| \leqslant C(x, \varepsilon) e^{\varepsilon, n \mid} \downarrow v \|
$$

3. Let $\gamma_{i j}(x)$ be the angle between the subspaces $E_{i x}$ and $E_{j x}, j=0,1,2, i \neq j$. Then $K(x, \epsilon) \leqslant \gamma(x)$.
1.3. For integers $s>r \geqslant 1$ consider the set
$\tilde{\Lambda}_{s, r}=\left\{x \in \tilde{\Lambda}: \frac{r-1}{s}<\lambda(x) \leqslant \frac{r}{s}<\frac{r+2}{s} \leqslant \mu(x)\right.$, where $s$ is the
smallest number satisfying these inequalities for some $r$ \}.
It is obvious that $\tilde{\Lambda}_{s, r}$ is measurable and invariant relative to $f$. Moreover, $\bigcup_{s>r \geqslant 1} \tilde{\Lambda}_{s, r}=$ $\tilde{\Lambda}$, and if $s_{1} \neq s_{2}$ or $r_{1} \neq r_{2}$, then $\tilde{\Lambda}_{s_{1}, r_{1}} \cap \tilde{\Lambda}_{s_{2}, r_{2}}=\varnothing$.

Consider the measurable invariant function defined on $\tilde{\Lambda}$ by the equality

$$
\begin{equation*}
\varepsilon(x)=\varepsilon_{s}=\frac{1}{100} \ln \left(1+\frac{2}{s}\right), \quad x \in \tilde{\Lambda}_{s, r} \tag{1.3.1}
\end{equation*}
$$

It is easy to verify that for each $x \in \widetilde{\Lambda}$

$$
\begin{equation*}
\lambda(x) e^{500(x)}<1, \quad \frac{r}{s} e^{50 \varepsilon(x)}<\frac{r+2}{s} e^{-\delta 0 \varepsilon(x)} \tag{1.3.2}
\end{equation*}
$$

In addition, the functions $C(x)=C(x, \epsilon(x))$ and $K(x)=K(x, \epsilon(x))$ are measurable on $\widetilde{\Lambda}$.
For each integer $l \geqslant 1$, consider the set

$$
\widetilde{\Lambda}_{s, r}^{l}=\left\{x \in \widetilde{\Lambda}_{s, r}: C(x) \leqslant l, K^{-1}(x) \leqslant l\right\}
$$

It is obvious that $\widetilde{\Lambda}_{s, r}^{l}$ is measurable, $\bigcup \widetilde{\Lambda}_{s, r}^{l}=\widetilde{\Lambda}_{s, r}$ and $\widetilde{\Lambda}_{s, r}^{l} \subset \widetilde{\Lambda}_{s, r}^{l+1}$. On the basis of (1.1.4) and (1.1.5), for each $x \in \widetilde{\Lambda}_{s, r}^{l}$, any $m \in \mathbf{Z}$ and any $n \in \mathbf{Z}^{+}$we obtain

$$
\begin{align*}
& \left\|d f_{f(x)}^{n} v\right\| \leqslant l\left(\frac{r}{s}\right)^{n} e^{\varepsilon_{s^{n+4}} \varepsilon_{s}|m|}\|v\| \\
& \left\|d f_{f^{m}(x)}^{-n} v\right\| \geqslant l^{-1}\left(\frac{r}{s}\right)^{-n} e^{-\varepsilon_{s^{n-4}} \varepsilon_{s}|m|}\|v\| \\
& \left\|d f_{f^{m}(x)}^{n} v\right\| \geqslant l^{-1}\left(\frac{r+2}{s}\right)^{n} e^{-\varepsilon_{s} n^{n-4 \varepsilon_{s}|m|}}\|v\|,  \tag{1.3.3}\\
& \quad\left(v \in E_{1 f^{m}(x)}\right), \\
& \left\|d f_{f^{m}(x)}^{n} v\right\| \leqslant l\left(\frac{r+2}{s}\right)^{-n} e^{\varepsilon_{s} n+4 \varepsilon_{s}|m|}\|v\| \\
& \quad\left(v \in E_{2 f^{m}(x)}\right), \\
& l^{-1} e^{-e_{s}|m|} \leqslant \gamma\left(f^{m}(x)\right) . \tag{1.3.4}
\end{align*}
$$

We denote by $\Lambda_{s, r}^{l}$ the set of points $x \in M$ satisfying the following conditions:
(1.3.5) There exist subspaces $E_{1 x}$ and $E_{2 x}$ for which $T_{x} M=E_{1 x} \oplus E_{2 x}$.
(1.3.6) For vectors v lying in $d f^{m}\left(E_{i x}\right), i=1,2$, the estimates (1.3.3) hold.
(1.3.7) The angle $\gamma\left(f^{m}(x)\right)$ between the subspaces $d f^{m}\left(E_{1 x}\right)$ and df ${ }^{m}\left(E_{2 x}\right)$ satisfies (1.3.4).

Theorem 1.3.1. 1. $\widetilde{\Lambda}_{s, r}^{l} \subset \Lambda_{s, r}^{l} \subset \Lambda$ and $\Lambda_{s, r}^{l} \subset \Lambda_{s, r}^{l+1}$.
2. The set $\Lambda_{s, r}^{I}$ is closed.
3. The subspaces $E_{1 x}$ and $E_{2 x}$ depend continuously on $x$ in the set $\Lambda_{s, r}^{l}$.
4. For each integer $q$ and $l \geqslant 1$ there is an $\alpha=\alpha(q, l, s) \in \mathbf{Z}^{+}$such that

$$
f^{q}\left(\Lambda_{s, r}^{l}\right) \subset \Lambda_{s, r}^{a} .
$$

5. The set $\Lambda_{s, r}=\bigcup_{l \geqslant 1} \Lambda_{s, r}^{l}$ is invariant relative to $f$.

Proof. 1. The inclusion $\widetilde{\Lambda}_{s, r}^{l} \subset \Lambda_{s, r}^{l}$ is obvious. If $x \in \Lambda_{s, r}^{l}$, then for any $v \in E_{1 x}$

$$
\chi^{+}(x, v)<\ln \frac{r}{s}+\varepsilon_{s}<0
$$

Thus $x \in \Lambda$. Since (1.3.3) and (1.3.4) remain valid on replacing $l$ by $l+1$, it follows that $\Lambda_{s, r}^{l} \subset \Lambda_{s, r}^{l+1}$; moreover, the subspaces $E_{1 x}$ and $E_{2 x}$ do not depend on which of the sets $\Lambda_{s, r}^{l}$ contain $x$.

2 and 3. It is easy to see that a decomposition satisfying (1.3.5)-(1.3.7) is unique. Let $x \in M$, and let $x_{i} \in \Lambda_{s, r}^{l}$ be a sequence of points converging to $x$. Passing to a subsequence, we may suppose that for each $i$

$$
\operatorname{dim} E_{1 x_{i}}=k, \quad \operatorname{dim} E_{2 x_{i}}=n-k
$$

and the subspaces $E_{1 x_{i}}$ and $E_{2 x_{i}}$ converge, as $i \rightarrow \infty$, to subspaces $E_{1 x}, E_{2 x} \subset T_{x} M$. We will prove that these subspaces satisfy (1.3.5)-(1.3.7). Let $v \in E_{1 x}$ and $v_{i} \in E_{1 x_{i}}, v_{i} \rightarrow v$. Since the inequalities (1.3.3) hold for $v_{i}$, on putting $m=0$ in these and passing to the limit as $i \rightarrow \infty$ we obtain that (1.3.3), with $m=0$, is valid for $v$. A similar assertion holds for vectors $v \in E_{2 x}$. Hence it follows that $E_{1 x} \cap E_{2 x}=0$, so that $T_{x} M=E_{1 x} \oplus E_{2 x}$. Fix an integer $m$. Since $x_{i} \rightarrow x$, it follows that $d f_{x_{i}}^{m} \rightarrow d f_{x}^{m}$, and consequently the subspaces $d f^{m}\left(E_{j x_{i}}\right)$ converge to the subspace $d f^{m}\left(E_{j k}\right), j=1,2$. Let the sequence of vectors $v_{i} \in$ $d f^{m}\left(E_{j x_{i}}\right)$ converge to a vector $v \in d f^{m}\left(E_{j x}\right), j=1,2$. Since (1.3.3), with a given $m$, is valid for the $v_{i}$, the same will be true for $v$. Arguing similarly, we obtain that (1.3.4) holds for each $x \in \Lambda_{s, r}^{l}$ and any $m$.
4. Let $q$ and $l>0$ be integers. Choose an integer $\alpha=\alpha(q, l, s)$ such that

$$
\begin{equation*}
l e^{4 \varepsilon_{s}|q|} \leqslant \alpha \tag{1.3.8}
\end{equation*}
$$

If $x \in f^{q}\left(\Lambda_{s, r}^{l}\right)$, we put $E_{j x}=d f^{q}\left(E_{\left.j r-q_{(x)}\right)}\right), j=1,2$. For $v \in E_{1 x}$, by (1.3.8) we have

$$
\begin{aligned}
& \left\|d f_{f(x)}^{n} v\right\|=\left\|d f_{\left.f m+q_{(f-}-q_{(x)}\right)}^{\prime} v\right\| \leqslant l\left(\frac{r}{s}\right)^{n} e^{\varepsilon_{s} n+4 \varepsilon_{s}|m+q|}\|v\| \\
& \leqslant l e^{4 \varepsilon_{s}|q|}\left(\frac{r}{s}\right)^{n} e^{\varepsilon_{s} n+4 \varepsilon_{s}|m|}\|v\| \leqslant \alpha\left(\frac{r}{s}\right)^{n} e^{\varepsilon_{s} n+4 \varepsilon_{s}|m|}
\end{aligned}
$$

The remaining estimates are deduced similarly. Thus $x \in \Lambda_{s, r}^{\alpha}$.
Assertion 5 follows immediately from 4. The theorem is proved.
Put

$$
\begin{aligned}
& \hat{\Lambda}=\bigcup_{s, r} \Lambda_{s, r}, \quad A_{k}=\left\{x \in \hat{\Lambda}: \operatorname{dim} E_{1 x}=k\right\} \\
& \Lambda_{k, s, r}=A_{k} \cap \Lambda_{s, r}, \quad \Lambda_{k, s, r}^{l}=A_{k} \cap \Lambda_{s, r}^{l} \\
& \tilde{\Lambda}_{k . s, r}=A_{k} \cap \widetilde{\Lambda}_{s, r}, \quad \widetilde{\Lambda}_{k, s, r}^{l}=A_{k} \cap \widetilde{\Lambda}_{s, r}^{l}
\end{aligned}
$$

It is obvious that the sets $\widetilde{\Lambda}_{k, s, r}^{l}$ and $\Lambda_{k, s, r}^{l}$ are measurable and the sets $\widetilde{\Lambda}_{k, s, r}$ and $\Lambda_{k, s, r}$ are measurable and invariant. We also put $l(x)=\operatorname{dim} E_{1 x}$.
1.4. The starting point for further constructions lies in the nonuniform (partial) hyperbolicity condition on the set $\widetilde{\Lambda}$ given by Theorem 1.1.1 (see also Theorem 1.2.1), where in (1.1.4) and (1.1.5) we must put $\epsilon=\epsilon(x)$ (see (1.3.1)). In fact this condition holds on a set $\hat{\Lambda} \supset \tilde{\Lambda}$. Therefore the following construction can be carried out on $\hat{\Lambda}$. We note however, that, by Theorem $0.3, \Lambda=\hat{\Lambda}=\widetilde{\Lambda}(\bmod 0)$. We stress that the hyperbolicity conditions on the sets $\Lambda_{k, s, r}^{l}$ are uniform, although the sets themselves are not invariant.
1.5. Now we construct a special (in general measurable) Riemannian metric, a systematic use of which will significantly simplify the arguments.

Consider the restriction $T \hat{\Lambda}$ of the tangent bundle $T M$ to the set $\hat{\Lambda}$. Since $\hat{\Lambda}$ is measurable (and also Borel), $T \hat{\Lambda}$ can be regarded as a measurable linear bundle in the sense of [7]. A trivialization of the measurable bundle $T \hat{\Lambda}$ is a measurable family of isomorphisms $\tau_{x}$ :
$T_{x} M \rightarrow \mathbf{R}^{n}$. A measurable Riemannian metric on $T \hat{\Lambda}$ is a measurable family of positive bilinear forms (scalar products) in the spaces $T_{x} M$. Let $\psi$ be a measurable Riemannian metric. A trivialization $\tau=\left\{\tau_{x}\right\}$ is called a Riemannian trivialization (relative to $\psi$ ) if for any $x \in$ $\hat{\Lambda}$ and $v_{1}, v_{2} \in T_{x} M$

$$
\begin{equation*}
\left\langle\tau_{x} v_{1}, \tau_{x} v_{2}\right\rangle=\psi_{x}\left(v_{1}, v_{2}\right) \tag{1.5.1}
\end{equation*}
$$

It is not difficult to show that a Riemannian trivialization exists for any measurable Riemannian metric.

Theorem 1.5.1. On $\hat{\Lambda}$ there exist a measurable Riemannian metric $\langle,\rangle_{x}^{\prime}$, measurable functions $\lambda^{\prime}(x)$ and $\mu^{\prime}(x)$ invariant relative to $f$, and a measurable function $A(x)$ such that for any $x \in \hat{\Lambda}$ and $m \in \mathbf{Z}$

$$
\begin{gather*}
0<\lambda(x)<\lambda^{\prime}(x)<\mu^{\prime}(x)<\mu(x), \quad \lambda^{\prime}(x)<1 ;  \tag{1.5.2}\\
A\left(f^{m}(x)\right) \leqslant A(x) e^{5 \varepsilon(x)|m|} ;  \tag{1.5.3}\\
\sup _{x \in \Lambda_{s, r}^{l}} A(x)=A_{s, r}^{l}<\infty ;  \tag{1.5.4}\\
\left\|\left.d f\right|_{E_{1 x}}\right\|^{\prime}<\lambda^{\prime}(x), \quad\left(\left\|\left.d f^{-1}\right|_{E_{2 x}}\right\|^{\prime}\right)^{-1} \geqslant \mu^{\prime}(x), \quad \gamma^{\prime}(x)=\frac{\pi}{2}, \tag{1.5.5}
\end{gather*}
$$

where $\gamma^{\prime}(x)$ denotes the angle between the subspaces $E_{1 x}$ and $E_{2 x}$ in the metric $\langle,\rangle_{x}^{\prime}$;

$$
\begin{equation*}
\frac{\sqrt{2}}{2}\|\cdot\|_{x} \leqslant\|\cdot\|_{x}^{\prime} \leqslant A(x)\|\cdot\|_{x} \tag{1.5.6}
\end{equation*}
$$

Proof. For $x \in \hat{\Lambda}$ put

$$
\begin{equation*}
\lambda^{\prime}(x)=\lambda(x) e^{3 \varepsilon(x)}, \quad \mu^{\prime}(x)=\mu(x) e^{-3 \varepsilon(x)} \tag{1.5.7}
\end{equation*}
$$

The inequality (1.5.2) follows from (1.3.2). Further, put

$$
\begin{gather*}
\left\langle v_{1}, v_{2}\right\rangle_{x}^{\prime}=\sum_{k=0}^{\infty}\left(\lambda^{\prime}(x)\right)^{-2 k}\left\langle d f^{k} v_{1}, d f^{k} v_{2}\right\rangle_{f^{k}(x)}, \quad v_{1}, v_{2} \in E_{1 x},  \tag{1.5.8}\\
\left\langle v_{1}, v_{2}\right\rangle_{x}^{\prime}=\sum_{k=0}^{\infty}\left(\mu^{\prime}(x)\right)^{2 k}\left\langle d f^{-k} v_{1}, d f^{-k} v_{2}\right\rangle_{f}-_{(x)}, \quad v_{1}, \quad v_{2} \in E_{2 x},
\end{gather*}
$$

and if $v_{i}=v_{i}^{1}+v_{i}^{2}$, where $v_{i}^{j} \in E_{j x}, i, j=1,2$, then

$$
\begin{equation*}
\left\langle v_{1}, v_{2}\right\rangle_{x}^{\prime}=\left\langle v_{1}^{1}, v_{2}^{1}\right\rangle_{x}^{\prime}+\left\langle v_{1}^{2}, v_{2}^{2}\right\rangle_{x}^{\prime} . \tag{1.5.9}
\end{equation*}
$$

We note that the series on the right-hand side of (1.5.8) converges. In addition, by (1.1.5), (1.5.2) and the Cauchy-Schwarz inequality,

$$
\begin{gather*}
\left|\left\langle v_{1}, v_{2}\right\rangle_{x}^{\prime}\right| \leqslant C^{2}(x)\left[1-\left(\frac{\lambda(x)}{\lambda^{\prime}(x)} e^{\varepsilon(x)^{\prime}}\right)^{2}\right]^{-1}\left\|v_{1}\right\|_{x}\left\|v_{2}\right\|_{x}, \\
v_{1}, v_{2} \in E_{1 x}  \tag{1.5.10}\\
\left|\left\langle v_{1}, v_{2}\right\rangle_{x}^{\prime}\right| \leqslant C^{2}(x)\left[1-\left(\frac{\mu^{\prime}(x)}{\mu(x)} e^{\varepsilon(x)}\right)^{2}\right]^{-1}\left\|v_{1}\right\|_{x}\left\|v_{2}\right\|_{x} . \\
v_{1}, v_{2} \in E_{2 x}
\end{gather*}
$$

Using (1.5.8), we obtain that for any $x \in \hat{\Lambda}$ and $v \in E_{1 x}$

$$
\begin{aligned}
& \|d f \cup\|^{\prime}=\left[\sum_{k=0}^{\infty}\left(\lambda^{\prime}(x)\right)^{-2 k}\left\|d f^{k} d f v\right\|_{f^{k+1}(x)}^{2}\right]^{1 / 2} \\
& =\left[\left(\lambda^{\prime}(x)\|v\|_{x}^{\prime}\right)^{2}-\left(\lambda^{\prime}(x)\|v\|_{x}\right)^{2}\right]^{1 / 2} \leqslant \lambda^{\prime}(x)\|v\|_{x}^{\prime} .
\end{aligned}
$$

Arguing similarly, it can be shown that for $v \in E_{2 x}$

$$
\|d f v\|^{\prime} \geqslant \mu^{\prime}(x)\|v\| \dot{x} .
$$

Finally, it follows from (1.5.9) that $\gamma(x)=\pi / 2$.
We denote by $\tau_{x}$ any Riemannian trivialization corresponding to the Riemannian metric $\langle,\rangle_{x}$. Let $\kappa^{i}(x)$ be orthogonal (in the metric $\langle,\rangle_{x}$ ) mappings:

$$
x^{1}(x): E_{1 x} \rightarrow \tau_{x}^{-1}\left(\mathbf{R}^{l(x)}\right), \quad x^{2}(x): E_{2 x} \rightarrow \tau_{x}^{-1}\left(\mathbf{R}^{n-l(x)}\right),
$$

and $\kappa(x): T_{x} M \rightarrow T_{x} M$ the mapping defined in the following way:

$$
x(x)\left(v_{1}+v_{2}\right)=x^{1}(x) v_{1}+x^{2}(x) v_{2},
$$

where $v_{i} \in E_{i x}, i=1,2$. It is easily seen that there is a constant $M>0$ such that for any $x \in \hat{\Lambda}$

$$
M_{\gamma}(x) \leqslant\|x(x)\|_{x} \leqslant \sqrt{2} .
$$

From this and (1.5.10) it follows that for each $x \in \hat{\Lambda}$ and any $v \in T_{x} M$ the inequality (1.5.6) is satisfied, where

$$
\begin{equation*}
A(x)=C(x) M^{-1} K^{-1}(x) \max \left\{\left[1-\left(\frac{\lambda(x)}{\lambda^{\prime}(x)} e^{\varepsilon(x)}\right)^{2}\right]^{-\frac{1}{2}},\left[1-\left(\frac{\mu^{\prime}(x)}{\mu(x)} e^{\varepsilon(x)}\right)^{2}\right]^{-\frac{1}{2}}\right\} \tag{1.5.11}
\end{equation*}
$$

The inequality (1.5.3) follows from (1.1.4), and (1.5.4) follows from the definition of $\Lambda_{s, r}^{l}$. The theorem is proved.

The measurable Riemannian metric $\langle,\rangle_{x}^{\prime}$ will be called a measurable:Ljapunov metric.
1.6. In this subsection we will introduce special coordinates in a neighborhood of each point $x \in \hat{\Lambda}$ with the help of a certain Riemannian trivialization, corresponding to the measurable Ljapunov metric, and also give the representation of $f$ in these coordinates. This representation will be used in $\S 2$ for the construction of families of local stable manifolds and in $\S 3$ for the proof of their absolute continuity.

For each $x \in \hat{\Lambda}$ consider the mapping

$$
\begin{equation*}
\tilde{f}_{x}=\exp _{f(x)}^{-1} \circ f \circ \exp _{x} \tag{1.6.1}
\end{equation*}
$$

defined in some neighborhood $U_{x} \subset T_{x} M$. The neighborhoods $U_{x}$ can be chosen so that for each $x \in \hat{\Lambda}$

$$
\begin{equation*}
\tau_{x}\left(U_{x}\right)=U \tag{1.6.2}
\end{equation*}
$$

Here $U$ is some fixed neighborhood of zero in $\mathbf{R}^{n}$ and $\tau_{x}$ is a Riemannian trivialization corresponding to the measurable Ljapunov metric. This trivialization can be chosen so that for each $x \in \hat{\Lambda}$

$$
\begin{equation*}
\tau_{x}^{\prime}\left(E_{1 x}\right)=\mathbf{R}^{l(x)}, \quad \tau_{x}^{\prime}\left(E_{2 x}\right)=\mathbf{R}^{n-l(x)} \tag{1.6.3}
\end{equation*}
$$

Consider the mapping

$$
\begin{equation*}
f_{x}^{\prime}=\tau_{f(x)}^{\prime} \circ f_{x} \circ\left(\tau_{x}^{\prime}\right)^{-1}: U \rightarrow \mathbf{R}^{n} \tag{1.6.4}
\end{equation*}
$$

which, on the basis of (1.6.3), can be written in the form

$$
\begin{equation*}
f_{x}^{\prime}(u, v)=\left(A_{x} u+g_{x}(u, v), B_{x} v+h_{x}(u, v)\right) . \tag{1.6.5}
\end{equation*}
$$

Неге $u \in \mathbf{R}^{l(x)}, v \in \mathbf{R}^{n-l(x)}, A_{x}: \mathbf{R}^{l(x)} \rightarrow \mathbf{R}^{l(x)}, B_{x}: \mathbf{R}^{n-l(x)} \rightarrow \mathbf{R}^{n-l(x)}, g_{x}: U \rightarrow \mathbf{R}^{l(x)}$, $g_{x} \in C^{r}, h_{x}: U \rightarrow \mathbf{R}^{n-l(x)}$ and $h_{x} \in C^{r}$.

Theorem 1.6.1.

$$
\begin{gather*}
\left\|A_{x}\right\| \leqslant \lambda^{\prime}(x), \quad\left\|B_{x}^{-1}\right\|^{-1} \geqslant \mu^{\prime}(x)  \tag{1.6.6}\\
g_{x}(0)=0, \quad h_{x}(0)=0, \quad d g_{x}(0)=0, \quad d h_{x}(0)=0 . \tag{1.6.7}
\end{gather*}
$$

There exists $M>0$ such that for any $x \in \hat{\Lambda}, z_{1}, z_{2} \in U$ and $l=1,2, \ldots, r-1$

$$
\begin{equation*}
\left\|d^{l} t_{x}\left(z_{1}\right)-d^{l} t_{x}\left(z_{2}\right)\right\| \leqslant M A(x)\left\|z_{1}-z_{2}\right\| \tag{1.6.8}
\end{equation*}
$$

where $t_{x}: T_{x} M \rightarrow T_{f(x)} M, t_{x}=\left(g_{x}, h_{x}\right)$.
Proof. Since $\tau_{x}^{\prime}$ is an isometry (see (1.5.1)), on the basis of (1.5.5), for any $u \in \mathbf{R}^{l(x)}$ we obtain

$$
\left\|A_{x} u\right\|=\left\|\tau_{f(x)}^{\prime} d f_{x}^{\prime}\left(\tau_{x}^{\prime}\right)^{-1} u\right\|=\left\|d f_{x}^{\prime}\left(\tau_{x}^{\prime}\right)^{-1} u\right\|_{f(x)}^{\prime} \leqslant \lambda^{\prime}(x)\left\|\left(\tau_{x}^{\prime}\right)^{-1} u\right\|_{x}^{\prime}=\lambda^{\prime}(x)\|u\|^{\prime}
$$

and if $v \in \mathbf{R}^{n-l(x)}$ similar inequalities give

$$
\left\|B_{x}^{-1} v\right\|^{-1} \geqslant \mu^{\prime}(x)\|v\|
$$

Since $f \in C^{r}, r \geqslant 2$, from the definition of $f_{x}$ (see (1.6.1)) and the compactness of $M$ follows the existence of a constant $C>0$ such that for any $x \in M, w_{1}, w_{2} \in T_{x} M$ and $l=1,2, \ldots$, $r-1$

$$
\left\|d^{2} f_{x}\left(w_{1}\right)-d^{l} f_{x}\left(w_{2}\right)\right\|_{f(x)} \leqslant C\left\|w_{1}-w_{2}\right\|_{x} .
$$

From this and (1.5.6) it follows that for any $z_{1}, z_{2} \in U$ and $l=1,2, \ldots, r-1$ ( $w_{i}=$ $\left.\left(\tau_{x}^{\prime}\right)^{-1} z_{i}, i=1,2\right)$

$$
\begin{gathered}
\left\|d^{l} t_{x}\left(z_{1}\right)-d^{l} t_{x}\left(z_{2}\right)\right\|=\left\|d^{l}\left(f_{x}^{\prime}\left(z_{1}\right)-d f_{x}^{\prime}(0) z_{1}\right)-d^{l}\left(f_{x}^{\prime}\left(z_{2}\right)-d f_{x}^{\prime}(0) z_{2}\right)\right\| \\
=\| \tau_{f(x)}^{\prime}\left(d^{l} f_{x}\left(w_{1}\right)-d^{l} f_{x}\left(w_{2}\right)\|=\| d^{l} f_{x}\left(w_{1}\right)-d^{l} f_{x}\left(w_{2}\right) \| f(x)\right. \\
\leqslant A(x)\left\|d^{l} f_{x}\left(w_{1}\right)-d^{l} f_{x}\left(w_{2}\right)\right\|_{f(x)} \leqslant A(x) C\left\|w_{1}-w_{2}\right\|_{x} \\
\leqslant A(x) C \sqrt{2}\left\|w_{1}-w_{2}\right\|_{x}^{\prime}=A(x) \sqrt{2} C\left\|z_{1}-z_{2}\right\|
\end{gathered}
$$

The theorem is proved.

## §2. The construction of local stable manifolds.

2.1. Let $F$ be Euclidean space with norm $\|\cdot\|$. Suppose that $F$ can be decomposed as a direct sum of subspaces $E_{1}$ and $E_{2}, F=E_{1} \oplus E_{2}$. Choose open convex neighborhoods of zero $U_{1}$ and $U_{2}$ in $E_{1}$ and $E_{2}$ respectively. Let $F_{m}$ be a countable family of Euclidean spaces with norms $\|\cdot\|_{m}, m \in \mathbf{Z}^{+}$, isomorphic to $F$ by isometries (a trivialization) $\tau_{m}: F_{m} \rightarrow$ F. Put $E_{i m}=\tau_{m}^{-1} E_{i}$ and $U_{i m}=\tau_{m}^{-1}\left(U_{i}\right), i=1,2$, and consider the family of mappings $f_{m}$,

$$
\begin{equation*}
\tau_{m+1} \circ f_{m} \circ \tau_{m}^{-1}(u, v)=\left(A_{m} u+g_{1 m}(u, v), B_{m} v+g_{2 m}(u, v)\right), \tag{2.1.1}
\end{equation*}
$$

where $u \in U_{1}$ and $v \in U_{2}$.
For brevity we denote the composition $f_{m}{ }^{\circ} f_{m-1}{ }^{\circ} \ldots{ }^{\circ} f_{0}=\Pi_{0}^{m} f_{i}$.
Theorem 2.1.1. Suppose that the following conditions are satisfied:

1. There exist numbers $\lambda$ and $\mu$ such that

$$
\begin{equation*}
0<\lambda<1, \quad \lambda<\mu \tag{2.1.2}
\end{equation*}
$$

and for any $m \in \mathbf{Z}^{+}$

$$
\begin{equation*}
\left\|A_{m}\right\| \leqslant \lambda, \quad\left\|B_{m}^{-1}\right\|^{-1} \geqslant \mu \tag{2.1.3}
\end{equation*}
$$

2. The function $g_{m}=\left(g_{1 m}, g_{2 m}\right) \in C^{1}$, and for any $m \in \mathbf{Z}^{+}$

$$
\begin{equation*}
g_{m}(0)=0, \quad d g_{m}(0)=0 \tag{2.1.4}
\end{equation*}
$$

3. There exist constants $K, \alpha$ and $\nu$ such that

$$
\begin{equation*}
\lambda^{\alpha}<v<1, \quad 0<\alpha \leqslant 1, \quad K>0 \tag{2.1.5}
\end{equation*}
$$

and for any $z_{1}, z_{2} \in U_{1} \times U_{2}$ and $m \in \mathbf{Z}^{+}$

$$
\begin{equation*}
\left\|d g_{m}\left(z_{1}\right)-d g_{m}\left(z_{2}\right)\right\| \leqslant K v^{-m}\left\|z_{1}-z_{2}\right\|^{\alpha} \tag{2.1.6}
\end{equation*}
$$

Let $\kappa$ be any number sat: ying

$$
\begin{equation*}
\lambda<x<\min \left(\mu, v^{\frac{1}{a}}\right) . \tag{2.1.7}
\end{equation*}
$$

Then there exist positive constants $C$ and $r_{0}$ and a mapping $\phi: S \rightarrow E_{20}$, where $S$ is the ball in $E_{10}$ with center at zero and radius $r_{0}$, satisfying the following conditions:

1) $\phi \in C^{1}$,

$$
\begin{equation*}
\varphi(0)=0, \quad d \varphi(0)=0 \tag{2.1.8}
\end{equation*}
$$

and for any $u_{1}, u_{2} \in S$

$$
\begin{equation*}
\left\|d \varphi\left(u_{1}\right)-d \varphi\left(u_{2}\right)\right\|_{0} \leqslant C\left\|u_{1}-u_{2}\right\|_{0}^{\alpha} . \tag{2.1.9}
\end{equation*}
$$

2) For each $m \in \mathbf{Z}^{+}$and any $u \in S$

$$
\begin{gather*}
\left(\prod_{t=0}^{m} f_{i}\right)(u, \varphi(u)) \in U_{1 m} \times U_{2 m} \\
\left\|\left(\prod_{i=0}^{m} f_{i}\right)(u, \varphi(u))\right\|_{m} \leqslant 200 x^{m}\|(u, \varphi(u))\|_{0} . \tag{2.1.10}
\end{gather*}
$$

3) Let $(u, v) \in F_{0}, u \in F$, and let there be a number $C>0$ such that for any $m \in \mathbf{Z}^{+}$

$$
\left(\prod_{i=0}^{m} f_{i}\right)(u, v) \in U_{1 m} \times U_{2 m}, \quad\left\|\left(\prod_{i=0}^{m} f_{i}\right)(u, v)\right\|_{m} \leqslant C x^{m} .
$$

Then $v=\phi(u)$.
4) There exist continuous functions $\psi_{i}=\psi_{i}(\lambda, \mu, \nu, \kappa, \alpha), i=1,2$, defined on the open set in $\mathbf{R}^{4}$ given by the inequalities (2.1.2), (2.1.5) and (2.1.7), such that

$$
\begin{equation*}
r_{0}=K^{-\frac{1}{\alpha}} \psi_{1}(\lambda, \mu, v, x, \alpha), \quad C=K^{-\frac{1}{\alpha}} \psi_{2}(\lambda, \mu, v, x, \alpha) \tag{2.1.11}
\end{equation*}
$$

Proof. Consider the linear space $\Gamma_{\kappa}$ of sequences of vectors $\{z(m)\}, z(m) \in F_{m}$, satisfying

$$
\|z\|_{x}=\sup _{m \in Z^{+}} x^{-m}\|z(m)\|_{m}<\infty
$$

The norm $\|\cdot\|_{\kappa}$ makes $\Gamma_{\kappa}$ into a Banach space. Consider the open set

$$
W=\left\{z \in \Gamma_{\kappa}: z(m) \in U_{1 m} \times U_{2 m}\right\}
$$

and the mapping $\Phi_{\kappa}: U_{1} \times W \rightarrow \Gamma_{\kappa}$ :

$$
\begin{gather*}
\tau_{0}\left(\Phi_{x}(y, z)(0)=\left(y,-\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k} B_{s}^{-1}\right) g_{2 k}\left(\tau_{k} z(k)\right)\right)\right. \\
\tau_{m} \Phi_{x}(y, z)(m)=\left(\left(\prod_{s=0}^{m-1} A_{s}\right) y+\sum_{k=0}^{m-1}\left(\prod_{s=k+1}^{m-1} A_{s}\right) g_{1 k}\left(\tau_{k} z(k)\right)\right.  \tag{2.1.12}\\
\left.-\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k} B_{s+m}^{-1}\right) g_{2 k+m}\left(\tau_{k+m} z(k+m)\right)\right)-z(m)
\end{gather*}
$$

Here for uniformity of notation we regard $\prod_{s=m}^{m-1} A_{s}=I$. First we will show that the mapping $\Phi_{\kappa}$ is well defined. For this we note that by the mean value theorem, (2.1.4) and (2.1.6), for any $z \in U_{1 m} \times U_{2 m}, m \in \mathbf{Z}^{+}$,

$$
\begin{gather*}
\left\|g_{m}\left(\tau_{m} z\right)\right\|=\left\|g_{m}\left(\tau_{m} z\right)-g_{m}(0)\right\| \leqslant\left\|d g_{m}(\xi)\right\|\left\|\tau_{m} z\right\|=\left\|d g_{m}(\xi)-d g_{m}(0)\right\|  \tag{2.1.13}\\
\times\left\|\tau_{m} z\right\| \leqslant K v^{-m}\|\xi\|^{\alpha}\left\|\tau_{m} z\right\| \leqslant K v^{-m}\left\|\tau_{m} z\right\|^{1+\alpha} .
\end{gather*}
$$

Here $\xi=\xi(m)$ lies on the segment joining 0 and $\tau_{m} z$. Using (2.1.3), (2.1.13) and the definition (2.1.12) of $\Phi_{\kappa}$, we get

$$
\begin{gathered}
\left\|\Phi_{x}(y, z)\right\|_{\kappa}=\sup _{m \geqslant 0} x^{-m}\left\|\tau_{m} \Phi_{x}(y, z)(m)\right\| \\
\leqslant \sup _{m \geqslant 0}\left\{x ^ { - m } \left[\prod_{s=0}^{m-1}\left\|A_{s}\right\|\|y\|+\sum_{k=0}^{m-1}\left(\prod_{s=k+1}^{m-1}\left\|A_{k}\right\|\right) K v^{-k}\|z(k)\|_{k}^{1+\alpha}\right.\right. \\
\left.\left.+\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k}\left\|B_{s+m}^{-1}\right\|\right) K v^{-(m+k)}\|z(m+k)\|_{m+k}^{1+\alpha}\right]\right\}+\sup _{m \geqslant 0} x^{-m}\|z(m)\|_{m} \\
\leqslant \sup _{m \geqslant 0} x^{-m} \lambda^{m}\|y\|+\sup _{m \geqslant 0} x^{-m} K\|z\|_{x}\left[\sum_{k=0}^{m-1} \lambda^{m-k-1} v^{-k} \chi^{(1+\alpha) k}\right. \\
\left.+\sum_{k=0}^{\infty} \mu^{-(k+1)} v^{-(m+k)} x^{(1+\alpha)(m+k)}\right]+\|z\|_{x} .
\end{gathered}
$$

(2.1.7) implies the estimates

$$
\begin{gather*}
\sup _{m \geqslant 0} \varkappa^{-m \lambda^{m}}=1,  \tag{2.1.14}\\
\sup _{m \geqslant 0} x^{-m \lambda^{m-1}} \sum_{k=0}^{m-1}\left(\lambda^{-1} v^{-1} x^{1+\alpha}\right)^{k} \leqslant-\lambda^{-1} e^{-1}\left(\ln \max \left\{\frac{x^{\alpha}}{v}, \frac{\lambda}{x}\right\}\right)^{-1},  \tag{2.1.15}\\
\sup _{m \geqslant 0} x^{-m} v^{-m} \mathcal{K}^{(1+\alpha) m} \mu^{-1} \sum_{k=0}^{\infty}\left(\mu^{-1} v^{-1} x^{1+\alpha}\right)^{k}=\frac{\mu^{-1}}{1-\mu^{-1} v^{-1} \mu^{1+\alpha}} \tag{2.1.16}
\end{gather*}
$$

Put

$$
\begin{equation*}
M=\max \left\{-\frac{1}{\lambda e \ln \max \left\{\frac{x^{\alpha}}{v}, \frac{\lambda}{\mu}\right\}}, \frac{\mu^{-1}}{1-\mu^{-1} v^{-1} x^{1+\alpha}}\right\} \tag{2.1.17}
\end{equation*}
$$

From (2.1.14)-(2.1.17) it follows that

$$
\begin{equation*}
\left\|\Phi_{x}(y, z)\right\|_{\varkappa} \leqslant\|y\|+2 K M\|z\|_{x}^{1+\alpha}+\|z\|_{x} . \tag{2.1.18}
\end{equation*}
$$

We have thus proved that $\Phi_{\kappa}$ is well defined. From (2.1.4) it follows that

$$
\begin{equation*}
\Phi_{\varkappa}(0,0)=(0,0) \tag{2.1.19}
\end{equation*}
$$

We will show that $\Phi_{\kappa} \in C^{1}$. For this it is sufficient to prove the existence and continuity of the partial derivatives of $\Phi_{\kappa}$ with respect to $y$ and $z$. Let $y \in U_{1}, h \in E_{1}$ and $y+h \in$ $U_{1}$. It follows from (2.1.12) that for any $z \in W$ and $m \in \mathbf{Z}^{+}$

$$
\tau_{m}\left[\Phi_{x}(y+h, z)(m)-\Phi_{\kappa}(y, z)(m)\right]=\left(\left(\prod_{s=0}^{m-1} A_{s}\right) h, 0\right)
$$

Therefore

$$
\begin{equation*}
\tau_{m} d_{y} \Phi_{x}(y, z)(m)=\left(\prod_{s=0}^{m-1} A_{s}, 0\right) \tag{2.1.20}
\end{equation*}
$$

We will prove the continuous differentiability of the mapping $\Phi_{\kappa}$ relative to $z$.
Consider the difference $\Phi_{\kappa}(y, z+h)-\Phi_{\kappa}(y, z)$, where $z \in W, h \in \Gamma_{\kappa}$ and $z+h \in W$. Using (2.1.12) and condition 2 of the theorem, we obtain

$$
\Phi_{x}(y, z+h)-\Phi_{x}(y, z)=\left(\gamma_{x}(z)-I\right) h+o(y, z, h)
$$

where $I$ is the identity mapping of $\Gamma_{\kappa}$ and

$$
\begin{gathered}
\gamma_{k}(z) h(m)=\tau_{m}^{-1}\left(\sum_{k=0}^{m-1}\left(\prod_{s=k+1}^{m-1} A_{s}\right) d g_{1 k}\left(\tau_{k} z(k)\right) \tau_{k} h(k)\right. \\
\left.-\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k} B_{s+m}^{-1}\right) d g_{2 m+k}\left(\tau_{m+k} z(m+k)\right) \tau_{m+k} h(m+k)\right) \\
o(y, z, h)(m)=\tau_{m}^{-1}\left(\sum_{k=0}^{m-1}\left(\prod_{s=k+1}^{m-1} A_{s}\right) o_{1}(z, h)(k)\right. \\
\left.-\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k} B_{s+m}^{-1}\right) o_{2}(z, h)(m+k)\right)
\end{gathered}
$$

Here $o_{i}(z, h)(m), i=1,2$, is defined by

$$
\begin{equation*}
o_{i}(z, h)(m)=g_{i m}\left(\tau_{m}(z+h)(m)\right)-g_{i m}\left(\tau_{m} z(m)\right)-d g_{i m}\left(\tau_{m} z(m)\right) \tau_{m} h(m) \tag{2.1.21}
\end{equation*}
$$

Inequality (2.1.6) and the mean value theorem imply

$$
\begin{gather*}
\left\|o_{i}^{\prime}(z, h)(m)\right\| \leqslant \| d g_{i m}(\xi(m)) \tau_{m} h(m)  \tag{2.1.22}\\
-d g_{i m}\left(\tau_{m} z(m)\right) \tau_{m} h(m)\left\|\leqslant K v^{-m}\right\| h(m) \|_{m}^{1+\alpha}
\end{gather*}
$$

where $\xi(m)$ lies on the segment joining $\tau_{m} z(m)$ and $\tau_{m}(z+h)(m)$.
It follows from (2.1.22) that for any $z_{1}, z_{2} \in W$ and $h \in \Gamma_{\kappa}$

$$
\begin{gathered}
\left\|\left(\gamma_{\varkappa}\left(z_{1}\right)-\gamma_{\varkappa}\left(z_{2}\right)\right) h\right\|_{x} \\
\leqslant \sup _{m \geqslant 0}\left\{x ^ { - m } \left[\sum_{k=0}^{m-1}\left(\prod_{s=k+1}^{m-1}\left\|A_{s}\right\|\right) K v^{-k}\left\|z_{1}(k)-z_{2}(k)\right\|_{k}^{\alpha}\left\|h_{1}(k)\right\|_{k}\right.\right. \\
\left.\left.+\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k}\left\|B_{s+m}^{-1}\right\|\right) K v^{-(m+k)}\left\|z_{1}(m+k)-z_{2}(m+k)\right\|_{m+k}^{\alpha}\|h(m+k)\|_{m+k}\right]\right\}
\end{gathered}
$$

The relations (2.1.14)-(2.1.17) imply

$$
\begin{equation*}
\left\|\left(\gamma_{x}\left(z_{1}\right)-\gamma_{x}\left(z_{2}\right)\right) h\right\|_{x} \leqslant 2 K M\left\|z_{1}-z_{2}\right\|_{x}\|h\|_{x} . \tag{2.1.23}
\end{equation*}
$$

In a similar way it is proved that

$$
\begin{equation*}
\|o(y, z, h)\|_{x} \leqslant 2 K M\|h\|_{x}^{1+\alpha} . \tag{2.1.24}
\end{equation*}
$$

It follows from (2.1.4) that $\gamma_{\kappa}(0)=0$. Therefore by (2.1.23)

$$
\begin{equation*}
\left\|\gamma_{x}(z)\right\|_{x} \leqslant 2 K M\|z\|_{\kappa}^{\alpha}, \quad z \in W \tag{2.1.25}
\end{equation*}
$$

This inequality together with (2.1.24) shows that $\Phi_{\kappa}$ is continuously differentiable with respect to $z$ and

$$
\begin{equation*}
d_{z} \Phi_{x}(y, z)=\gamma_{x}(z)-I . \tag{2.1.26}
\end{equation*}
$$

In addition, for any $y \in U_{1}$

$$
\begin{equation*}
d_{z} \Phi_{\varkappa}(y, 0)=-I \tag{2.1.27}
\end{equation*}
$$

From the above and the equalities (2.1.19) and (2.1.27) it follows that $\Phi_{\kappa}$ satisfies all the conditions of the implicit function theorem. In what follows we will need a version of this theorem, which we quote below.

Lemma 2.1.1. Let $E, F$ and $G$ be three Banach spaces and $f$ a continuously differentiable mapping of the product $A=A_{1} \times A_{2}$ into $G$, where $A_{1} \subset E, A_{2} \subset F$, and $A_{i}$ is a ball about zero with radius $r_{i}, i=1,2$. Suppose that df satisfies a Hölder condition in $A$ with constant a and exponent $\alpha, f(0,0)=0$ and the partial derivative $T_{0}=D_{2} f(0,0)$ is a linear isomorphism of $F$ onto $G$. Let $S$ be the ball in $E$ with center at zero and radius $r_{0}$,

$$
\begin{equation*}
r_{0}=\min \left\{r_{1}, r_{2}, \frac{r_{2}}{2 c b}, \frac{1}{(1+2 c b)(2 c a)^{1 / a}}\right\}, \tag{2.1.28}
\end{equation*}
$$

where $b=\max _{(x, 0) \in A}\left\|d_{x} f(x, 0)\right\|$ and $C=\left\|T_{0}^{-1}\right\|$. Then there exists a unique mapping $u: S \rightarrow A_{2}$ of class $C^{1}$ satisfying the following conditions:

$$
\begin{gather*}
f(x, u(x))=0, \quad u(0)=0 \\
\left\|\frac{d u}{d x}\left(x_{1}\right)-\frac{d u}{d x}\left(x_{2}\right)\right\| \leqslant 8 a c(1+2 b c)^{2}\left\|x_{1}-x_{2}\right\|^{a}  \tag{2.1.29}\\
\left\|\frac{d u}{d x}(x)\right\| \leqslant 1+2 b c
\end{gather*}
$$

Proof. We use the method of proof of Theorem 10.21 in [5]. Denote

$$
g(x, y)=y-T_{0}^{-1} f(x, y) .
$$

Let $\left(x, y_{1}\right),\left(x, y_{2}\right) \in A$. Then

$$
g\left(x, y_{1}\right)-g\left(x, y_{2}\right)=T_{0}^{-1}\left(T_{0}\left(y_{1}-y_{2}\right)-\left(f\left(x, y_{1}\right)-f\left(x, y_{2}\right)\right)\right)
$$

Put $r_{3}=2 c b r_{0}$. By (2.1.28), $r_{0} \leqslant r_{1}, r_{3} \leqslant r_{2}$. Let $S \subset E$ and $Q \subset F$ be balls about zero with radii $r_{0}$ and $r_{3}$ respectively. If $x \in S$ and $y_{1}, y_{2} \in Q$, then by the choice of $r_{0}$ and $r_{3}$ and the mean value theorem we have

$$
\left\|f\left(x, y_{1}\right)-f\left(x, y_{2}\right)-T_{0}\left(y_{1}-y_{2}\right)\right\| \leqslant a\left(r_{0}+r_{3}\right)^{\alpha}\left\|y_{1}-y_{2}\right\| \leqslant \frac{1}{2 c}\left\|y_{1}-y_{2}\right\|
$$

Therefore for any $x \in S$ and $y_{1}, y_{2} \in Q$

$$
\left\|g\left(x, y_{1}\right)-g\left(x, y_{2}\right)\right\| \leqslant \frac{1}{2}\left\|y_{1}-y_{2}\right\|
$$

On the other hand, for any $x \in S$

$$
\|g(x, 0)\|=\left\|-T_{0}^{-1} f(x, 0)\right\| \leqslant c b r_{0}=\frac{1}{2} r_{3} .
$$

Thus the mapping $g: S \times Q \rightarrow Q$ satisfies the conditions of a fixed point theorem of [5] (Theorem 10.14), from which follows the existence and uniqueness of a continuous mapping $u: S \rightarrow Q$ such that $u(0)=0$ and $f(x, u(x))=0$ for each $x \in S$. The proof that $u$ is continuously differentiable in $S$ is a verbatim repeat of the corresponding argument in [5] (see p. 267). From the inequality

$$
\|d f(x, y)-d f(0,0)\| \leqslant a\left(r_{0}+r_{3}\right)^{\alpha} \leqslant \frac{1}{2 c}
$$

valid for all $(x, y) \in S \times Q$, it follows that

$$
\begin{equation*}
\left\|d_{x} f(x, y)\right\| \leqslant b+\frac{1}{2 c} \tag{2.1.30}
\end{equation*}
$$

in addition,

$$
\left\|d_{y f} f(x, y)\right\| \geqslant \frac{1}{c}-\frac{1}{2 c}=\frac{1}{2 c} .
$$

Therefore the mapping $d_{y} f(x, y)$ is invertible in $S \times Q$, and it is not difficult to show that

$$
\begin{equation*}
\left\|d_{y}^{-1} f(x, y)\right\| \leqslant 2 c \tag{2.1.31}
\end{equation*}
$$

Differentiating the equality $f(x, u(x))=0$ with respect to $x$, we get

$$
d_{x} f(x, u(x))+d_{y f} f(x, u(x)) \frac{d u(x)}{d x}=0 .
$$

Inequalities (2.1.30) and (2.1.31) imply that for $x \in S$

$$
\begin{equation*}
\left\|\frac{d u(x)}{d x}\right\| \leqslant 1+2 b c . \tag{2.1.32}
\end{equation*}
$$

In addition, for any $x_{1}, x_{2} \in S$

$$
\begin{gathered}
d_{x} f\left(x_{1}, u\left(x_{1}\right)\right)-d_{x} f\left(x_{2}, u\left(x_{2}\right)\right)+\left(d_{y} f\left(x_{1}, u\left(x_{1}\right)\right)-d_{y} f\left(x_{2}, u\left(x_{2}\right)\right)\right) \frac{d u\left(x_{1}\right)}{d x} \\
+d_{y} f\left(x_{2}, u\left(x_{2}\right)\right)\left(\frac{d u\left(x_{1}\right)}{d x}-\frac{d u\left(x_{2}\right)}{d x}\right)=0 .
\end{gathered}
$$

Inequalities (2.1.30), (2.1.31) and (2.1.32) imply

$$
\begin{gathered}
\left\|\frac{d u\left(x_{1}\right)}{d x}-\frac{d u\left(x_{2}\right)}{d x}\right\| \leqslant 4 a c(1+b c)^{2}\left\|x_{1}-x_{2}\right\|^{\alpha} \\
+4 a c(1+2 b c)^{2}\left\|x_{1}-x_{2}\right\|^{\alpha} \leqslant 8 a c(1+2 b c)^{2}\left\|x_{1}-x_{2}\right\|^{\alpha} .
\end{gathered}
$$

The lemma is proved.
It follows from (2.1.23), (2.1.20) and (2.1.27) that $\Phi_{\kappa}$ satisfies the conditions of the lemma with

$$
\begin{equation*}
c=1, \quad b=1, \quad a=2 K M \tag{2.1.33}
\end{equation*}
$$

According to the lemma there exists a ball $S$ in $E_{10}$ of radius $r_{0}, 0 \in S \subset U$, and a mapping $\psi: S \rightarrow W, \psi \in C^{1}$, such that for any $y \in S$

$$
\begin{equation*}
\psi(0)=0, \quad \Phi_{x}(y, \psi(y))=0 \tag{2.1.34}
\end{equation*}
$$

By virtue of (2.1.17), (2.1.28) and (2.1.33) the number $r_{0}$ has the form (2.1.11).
Differentiating the second equality in (2.1.34) with respect to $y$, we get

$$
\begin{equation*}
d \psi(y)=-\left[d_{\mathbf{2}} \Phi_{x}(y, \psi(y))\right]^{-1} d_{y} \Phi_{x}(y, \psi(y)) \tag{2.1.35}
\end{equation*}
$$

Putting $y=0$ in this equality and taking account of (2.1.20), (2.1.27) and (2.1.34), we find that

$$
\begin{equation*}
d \psi(0)(m)=\left(\prod_{s=0}^{m-1} A_{s}, 0\right) \tag{2.1.36}
\end{equation*}
$$

We represent $\psi(y)(m)$ in the form $\psi(y)(m)=\left(\psi_{1}(y)(m), \psi_{2}(y)(m)\right)$, where $\psi_{i}(y)(m) \in E_{i m}, i=1,2$. It follows from (2.1.12) and (2.1.34) that

$$
\begin{gather*}
\tau_{0} \psi_{1}(y)(0)=y, \\
\tau_{m} \psi_{1}(y)(m)=\left(\prod_{s=0}^{m-1} A_{5}^{-}\right) y+\sum_{k=0}^{m-1}\left(\prod_{s=k+1}^{m-1} A_{s}\right) g_{1 k}\left(\tau_{k} \psi(y)(k)\right), \quad m \geqslant 1,  \tag{2.1.37}\\
\tau_{m} \psi_{2}(y)(m)=-\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k} B_{s+m}^{-1}\right) g_{2 m+k}\left(\tau_{m+k} \psi(y)(m+k)\right), \quad m \geqslant 0 . \tag{2.1.38}
\end{gather*}
$$

It is easy to see that these equalities imply the relations

$$
\begin{aligned}
& \tau_{m+1} \psi_{1}(y)(m+1)=A_{m} \tau_{m} \psi_{1}(y)(m)+g_{1 m}\left(\tau_{m} \psi(y)(m)\right), \\
& \tau_{m+1} \psi_{2}(y)(m+1)=B_{m} \tau_{m} \psi_{2}(y)(m)+g_{2 m}\left(\tau_{m} \psi(y)(m)\right),
\end{aligned}
$$

which mean that for all $m \in \mathbf{Z}^{+}$and $u \in S$

$$
\begin{equation*}
f_{m}(\psi(y)(m))=\psi(y)(m+1) \tag{2.1.39}
\end{equation*}
$$

Define $\phi: S \rightarrow E_{20}$ from the conditions

$$
\begin{equation*}
\varphi(u)=\psi_{2}(y)(0), \quad u=\psi_{1}(y)(0), \quad y \in S \tag{2.1.40}
\end{equation*}
$$

We will show that $\phi(u)$ satisfies assertions $1-3$ of the theorem. The differentiability of $\phi$ follows from the differentiability of $\psi$, and (2.1.8) follows from (2.1.34) and (2.1.36). Using (2.1.29), (2.1.33), (2.1.37) and (2.1.40), we obtain that for any $u_{1}, u_{2} \in S$

$$
\begin{gathered}
\left\|d \varphi\left(u_{1}\right)-d \varphi\left(u_{2}\right)\right\|_{0} \leqslant\left\|\left(d \psi\left(y_{1}\right)-d \psi\left(y_{2}\right)\right)(0)\right\|_{0} \\
\leqslant\left\|d \psi\left(y_{1}\right)-d \psi\left(y_{2}\right)\right\|_{x} \leqslant 200 K M\left\|y_{1}-y_{2}\right\|^{\alpha}=200 K M\left\|u_{1}-u_{2}\right\|^{\alpha}
\end{gathered}
$$

This proves (2.1.9). Assertion 1 is proved. In addition, $C=200 \mathrm{KM}$ has the form (2.1.11).
Since $\psi(y) \in W$ for any $y \in S$, using (2.1.29), (2.1.33), (2.1.37) and (2.1.40), we find that for any $u \in S$ and $m \in \mathbf{Z}^{+}$

$$
\begin{gathered}
\left\|\left(\prod_{i=1}^{m} f_{i}\right)(u, \varphi(u))\right\|_{m}=\left\|\left(\prod_{i=1}^{m} f_{i}\right)(\psi(y)(0))\right\|_{m}=\|\psi(y)(m)\|_{m} \leqslant x^{m}\|\psi(y)\|_{s} \\
\leqslant x^{m}\|d \psi(y)\|\left\|y \leqslant 200 x^{m}\right\| y\left\|=200 x^{m}\right\| u\left\|_{0} \leqslant 200 x^{m}\right\|(u, \varphi(u)) \|_{0} .
\end{gathered}
$$

In addition,

$$
\left(\prod_{i=1}^{m} f_{i}\right)(u, \varphi(u))=\psi(y)(m) \in U_{1 m} \times U_{2 m}
$$

Assertion 2 is proved.
Let the point $(u, v)$ and the number $C>0$ be chosen in accordance with the conditions of assertion 3. Consider the sequence $\psi(l)=\left(\Pi_{i=1}^{l} f_{i}\right)(u, v)$. Since

$$
\sup _{l \geqslant 0} x^{-l}\left\|\left(\prod_{i=1}^{l} f_{i}\right)\left(u, i^{\prime}\right)\right\|_{l \geq 0} \leqslant C \sup _{D} x^{-l} x^{l}=C,
$$

we see that $\psi \in \Gamma_{\kappa}$. From the definition of $\psi$ it follows that $f_{l} \psi(l)=\psi(l+1)$. Therefore

$$
\begin{gathered}
\tau_{l+1} \psi_{1}(l+1)=A_{l} \tau_{l} \psi_{1}(l)+g_{1 l}\left(\tau_{l} \psi(l)\right), \\
\tau_{l+1} \psi_{2}(l+1)=B_{l} \tau_{l} \psi_{2}(l)+g_{2 l}\left(\tau_{l} \psi(l)\right) .
\end{gathered}
$$

Hence it follows that for any $n \geqslant l$

$$
\begin{gather*}
\tau_{l} \psi_{1}(l)=\left(\prod_{s=0}^{l-1} A_{s}\right) y+\sum_{k=0}^{i-1}\left(\prod_{s=k+1}^{l-1} A_{s}\right) g_{1 k}\left(\tau_{k} \psi(k)\right), \\
\tau_{l} \psi_{2}(l)=\left(\prod_{s=0}^{n-l-1} B_{s+l}^{-1}\right) \tau_{n} \psi_{2}(n)-\sum_{k=0}^{n-l-1}\left(\prod_{s=0}^{k} B_{s+l}^{-1}\right) g_{2 k+l}\left(\tau_{k+l} \psi(k+l)\right) . \tag{2.1.41}
\end{gather*}
$$

Since

$$
\left\|\left(\prod_{s=0}^{n-l-1} B_{s+l}^{-1}\right) \tau_{n} \psi_{2}(n)\right\| \leqslant \frac{1}{\mu^{n-l}}\left\|\left(\prod_{i=0}^{n-l-1} f_{i+l}\right) \psi(l)\right\|_{n} \leqslant C\left(\frac{\gamma}{\mu}\right)^{n-l}
$$

on passing to the limit as $n \rightarrow \infty$ in the second equality of (2.1.41) we obtain

$$
\tau_{l} \psi_{2}(l)=-\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k} B_{s+l}^{-1}\right) g_{2 k+l}\left(\tau_{k+l} \psi(k+l)\right) .
$$

Hence from (2.1.38) and (2.1.41) it follows that $\Phi_{\kappa}(y, \psi)=0$. Since $y \in S$, by the uniqueness statement in Lemma 2.1.1 we have $\psi=\psi(y)$. Therefore

$$
\psi_{1}(y)=\psi_{1}(0)=u, \quad \varphi(u)=\psi_{2}(y)(0)=\psi(0)=v .
$$

Assertion 3, and together with it the theorem, is proved.
Under additional smoothness assumptions on $f_{m}$ it can be proved that the mapping $\phi$ constructed in Theorem 2.1.1 also has a higher degree of smoothness.

Theorem 2.1.2. Assume that, in the conditions of Theorem 2.1.1, for each $m \in \mathbf{Z}^{+}$ the function $g_{m}(u, v) \in C^{r},(u, v) \in U_{1} \times U_{2}, r \geqslant 2$, and moreover there exist positive constants $K$ and $K_{l}, l=1, \ldots, r$, such that

$$
\begin{gather*}
\sup _{z \in U_{1} \times U_{z}}\left\|d^{l} g_{m}(z)\right\| \leqslant K_{l} v^{-m},  \tag{2.1.42}\\
\left\|d^{r} g_{m}\left(z_{1}\right)-d^{r} g_{m}\left(z_{2}\right)\right\| \leqslant K v^{-m}\left\|z_{1}-z_{2}\right\|^{a}, \tag{2.1.43}
\end{gather*}
$$

where $z_{1}, z_{2} \in U_{1} \times U_{2}$. Let $\phi(u)$ be the mapping constructed in Theorem 2.1.1. Then there exist positive numbers $r_{0}, N$ and $N_{l}, l=1, \ldots, r$, depending only on $K, K_{l}, \lambda, \mu, \nu$, $\kappa$ and $\alpha$, such that for any $u, u_{1}, u_{2} \in S$ and $\phi(u) \in C^{r}$

$$
\begin{gather*}
\sup _{u \in S}\left\|d^{l} \varphi(u)\right\| \leqslant N_{l}  \tag{2.1.44}\\
\left\|d^{r} \varphi\left(u_{1}\right)-d^{r} \varphi\left(u_{2}\right)\right\| \leqslant N\left\|u_{1}-u_{2}\right\|^{\alpha} . \tag{2.1.45}
\end{gather*}
$$

Proof. We use the notation and constructions of Theorem 2.1.1 and show that $\Phi_{\kappa} \in$ $C^{r}$. It follows from (2.1.20) that for any $y \in U_{1}$ and $z \in W$

$$
d_{y}^{l} \Phi_{x}(y, z)=(0,0), \quad 2 \leqslant l \leqslant r
$$

We will show that $\Phi_{\kappa}$ is $r$ times continuously differentiable with respect to $z$. Formally differentiating (2.1.12) $l$ times with respect to $z$ gives

$$
\begin{gathered}
\tau_{m} d_{z}^{l} \Phi_{\mu}(y, z)(m)=\left(\sum_{k=0}^{m-1}\left(\prod_{s=k+1}^{m-1} A_{s}\right) d^{l} g_{1 k}\left(\tau_{k} z(k)\right)\right. \\
\left.-\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k} B_{s+m}^{-1}\right) d^{l} g_{2 m+k}\left(\tau_{m+k} z(m+k)\right)\right)
\end{gathered}
$$

The inequality (2.1.42) means that the multilinear form $d^{l} g_{m}(z)$ is uniformly bounded; that is, for any $z \in U_{1} \times U_{2}$ and $h_{1}, \ldots, h_{l} \in F, l \geqslant 2$, we have

$$
\left\|d^{l} g_{m}(z)\left(h_{1}, \ldots, h_{l}\right)\right\| \leqslant K_{l} v^{-m} \prod_{i=1}^{i}\left\|h_{i}\right\| .
$$

Let $h_{i} \in \Gamma_{\kappa}$. From the preceding inequality and (2.1.14)-(2.1.17) it follows that

$$
\begin{aligned}
& \left\|d_{z}^{l} \Phi_{x}(y, z)\left(h_{1}, \ldots, h_{l}\right)\right\|_{x} \leqslant \sup _{m \geqslant 0} x^{-m}\left[\sum_{k=0}^{m-1} \lambda^{m-k-1} K_{l} v^{-k} x^{l k} \prod_{i=1}^{l}\left\|h_{i}\right\|_{x}\right. \\
& \left.\quad+\sum_{k=0}^{\infty} \mu^{-(k+1)} K_{l} v^{-(m+k)} x^{(m+k) l} \prod_{i=1}^{l}\left\|h_{i}\right\|_{x}\right] \leqslant 2 K_{l} M \prod_{i=1}^{l}\left\|h_{i}\right\|_{x} .
\end{aligned}
$$

Thus the multilinear form $d_{z}^{l} \Phi_{\kappa}(y, z)$ is bounded. Obviously it is continuous relative to $y$ and $z$. Now we show by induction that $d_{z}^{l} \Phi_{k}(y, z)$ is in fact the $l$ th derivative with respect to $z$ of $\Phi_{\kappa}(y, z)$. For $l=1$ this was proved in Theorem 2.1.1. If $h \in \Gamma_{\kappa}$ is a small increment, then

$$
d_{z}^{l-1} \Phi_{x}(y, z+h)-d_{z}^{l-1} \Phi_{x}(y, z)=d_{z}^{l} \Phi_{x}(y, z) h+o_{l}(y, z, h)
$$

where $o_{l}(y, z, h)$ is defined by

$$
\begin{gathered}
\boldsymbol{\tau}_{m} o_{l}(y, z, h)(m)=\left(\sum_{k=0}^{m-1}\left(\prod_{s=k+1}^{m-1} A_{s}\right) o_{1 l}(y, z, h)(k),\right. \\
\left.-\sum_{k=0}^{\infty}\left(\prod_{s=0}^{k} B_{s+m}^{-1}\right) o_{2 l}(y, z, h)(m+k)\right)
\end{gathered}
$$

moreover $o_{i l}(y, z, h)$ is defined by (2.1.22), in which $g_{i}$ is replaced by its $l$-derivative. Using (2.1.42) and (2.1.43), we obtain that

$$
\begin{aligned}
& \left\|o_{i l}(y, z, h)(m)\right\|=\| d^{l} g_{i m}(\xi(m)) \tau_{m} h(m) \\
& -d^{l} g_{i m}\left(\tau_{m} z(m)\right) \tau_{m} h(m)\left\|\leqslant K_{i+1} v^{-m}\right\| h(m) \|_{m}^{2} \quad \text { for } 2 \leqslant l<r
\end{aligned}
$$

and

$$
\left\|o_{i r}(y, z, h)(m)\right\| \leqslant K v^{-m}\|h(m)\|_{m}^{1+\alpha} \quad \text { for } l=r .
$$

Hence, arguing as above and using (2.1.14)-(2.1.17), we obtain

$$
\left\|o_{l}(y, z, h)\right\|_{x} \leqslant 2 K_{l+1} M\|h\|_{x}^{2} \quad \text { for } 2 \leqslant l<r
$$

and

$$
\left\|o_{r}(y, z, h)\right\|_{x} \leqslant 2 K M\|h\|_{\varkappa}^{l+a} \text { for } l=r
$$

This proves that $\Phi_{\kappa} \in C^{r}$. From the implicit function theorem we conclude that $\psi \in C^{r}$. Furthermore, there exist constants $N_{1}>0$, depending only on $\lambda, \mu, \nu, \kappa$ and $\alpha$, such that

$$
\sup _{y \in S}\left\|d^{l} \psi(y)\right\|_{x} \leqslant N_{l}, \quad 2 \leqslant l \leqslant r
$$

We will use a version of the implicit function theorem whose proof is similar to that of Lemma 2.1.1.

Lemma 2.1.2. In the conditions of Lemma 2.1.1 assume that $f \in C^{r}$ and $d^{r} f$ satisfies a Hölder condition in $A$ with constant $a$ and exponent $\alpha$. Then $u \in C^{r}$, and there is a constant $N$ such that for any $x_{1}, x_{2} \in S$

$$
\begin{equation*}
\left\|\frac{d^{r} u\left(x_{1}\right)}{d x^{r}}-\frac{d^{r} u\left(x_{2}\right)}{d x^{r}}\right\| \leqslant N\left\|x_{1}-x_{2}\right\|^{\alpha} . \tag{2.1.46}
\end{equation*}
$$

From the definition of $\phi(u)$ (see (2.1.40)) and Lemma 2.1.2 (see (2.1.46)) it follows that for any $u_{1}, u_{2} \in S$

$$
\left\|d^{r} \varphi\left(u_{1}\right)-d^{r} \varphi\left(u_{2}\right)\right\|_{0} \leqslant\left\|d^{r} \psi\left(y_{1}\right)-d^{r} \psi\left(y_{2}\right)\right\|_{x} \leqslant N\left\|u_{1}-u_{2}\right\|_{0}^{\alpha}
$$

The theorem is proved.
2.2. Let $f$ be a diffeomorphism of $M$ of class $C^{r}, r \geqslant 2$, and $\Lambda$ the invariant set of positive measure defined by (0.2). We use the notation of 1.5 and consider the bundle $T \hat{\Lambda}$ with the Riemannian trivialization $\tau^{\prime}$, corresponding to the Ljapunov measurable metric $\langle,\rangle_{x}^{\prime}$, constructed in Theorem 1.5.1.

Proposition 2.2.1. Let $x \in \Lambda_{s, r}$. The family of trivializations

$$
\tau_{m}=\tau_{f(m)}^{\prime}: T_{f m_{(x)}} \rightarrow \mathbf{R}^{n}
$$

and mappings $f_{m}=f_{f m_{(x)}}$ (see (1.6.6)), defined in the neighborhoods $U_{m}=U_{f m_{(x)}}$ (see (1.6.2)), satisfies the conditions of Theorems 2.1.1 and 2.1.2, where

$$
\begin{gather*}
\lambda=\lambda_{s, r}=\frac{r}{s} e^{3 \varepsilon_{s}}, \quad \mu=\mu_{s, r}=\frac{r+2}{s} e^{-3 \varepsilon_{s}}, \\
\nu=i_{s, r}=e^{-j \varepsilon_{s}}, \quad x=x_{s, r}=\lambda_{s, r} e^{\varepsilon_{s}}  \tag{2.2.1}\\
\alpha=1, \quad \varepsilon_{s}=\frac{1}{100} \ln \left(1+\frac{2}{s}\right), \quad K=M A(x)
\end{gather*}
$$

Here the function $A(x)$ constructed in Theorem 1.5 .1 and the constant $M$ are the same as in (1.6.8).

Proof. The representation (2.1.1) follows from (1.6.3)-(1.6.5). Inequalities (2.i.2), (2.1.5) and (2.1.7) follow from (1.3.2); inequality (2.1.3), from (1.6.6) and (1.5.7). Conditions (2.1.4) are corollaries of (1.6.7). The inequalities (1.6.8) and (1.5.3) imply that for any $m \in \mathbf{Z}^{+}, z_{1}, z_{2} \in U($ see $(1.6 .2))$ and $l=1, \ldots, r-1$

$$
\left\|d^{l} t_{m}\left(z_{1}\right)-d^{l} t_{m}\left(z_{2}\right)\right\| \leqslant M A(x) e^{b \varepsilon_{s^{m}}}\left\|z_{1}-z_{2}\right\|
$$

Finally, (2.1.43) is an immediate corollary of (1.6.8). The proposition is proved.
Put

$$
\begin{equation*}
x(x)=k_{s, r}, \quad x \in \Lambda_{s, r} \tag{2.2.2}
\end{equation*}
$$

If $\delta(x)$ is a positive measurable function on $\hat{\Lambda}$ we denote

$$
\begin{gather*}
B^{i}(\delta(x))=\left\{u \in E_{i x}:\|u\|_{x} \leqslant \delta(x)\right\}, \quad i=1,2 \\
B(\delta(x))=B^{1}(\delta(x)) \times B^{2}(\delta(x))  \tag{2.2.3}\\
U(x, \delta(x))=\exp _{x} B(\delta(x)), \quad U_{x}^{\prime}=\exp _{x}\left(\tau_{x}^{\prime}\right)^{-1} Q
\end{gather*}
$$

Theorem 2.2.1. There exist a measurable function $\delta(x), x \in \hat{\Lambda}$, a family of mappings $\phi(x): B^{1}(\delta(x)) \rightarrow B^{2}(\delta(x))$ of class $C^{r-1}$, depending measurably on $x \in \hat{\Lambda}$, and a constant $N>0$ satisfying the following conditions:

1. The set $V(x)=\left\{\exp _{x}(u, \phi(x) u): u \in B^{1}(\delta(x))\right\}$ is a submanifold in $M$ of class $C^{r-1}$.
2. $x \in V(x)$.
3. $T_{x} V(x)=E_{1 x}$.
4. For $y \in V(x)$ and $n \in \mathbf{Z}^{+}$we have $f^{n}(y) \in U_{f(x)}^{\prime}$ and

$$
\begin{equation*}
\rho\left(f^{n}(x), f^{n}(y)\right) \leqslant N A(x) \mathcal{x}^{n}(x) \rho(x, y) \tag{2.2.4}
\end{equation*}
$$

5. If there are a point $y \in U(x, \delta(x))$ and a constant $C>0$ such that for each $n \in \mathbf{Z}^{+}$ we have $f^{n}(y) \in U_{f^{n}(x)}^{\prime}$ and

$$
\begin{equation*}
\rho\left(f^{n}(x), f^{n}(y)\right) \leqslant C x^{n}(x) \tag{2.2.5}
\end{equation*}
$$

then $y \in V(x)$.
6. $U(x, \delta(x)) \in U_{x}^{\prime}$, where $U_{x}^{\prime}$ is a neighborhood of $x$ (see (2.2.3)) and for any $m \in \mathrm{Z}^{+}$

$$
\begin{gather*}
\delta\left(f^{m}(x)\right) \geqslant \delta(x) e^{-108(x) m} \\
\delta_{s, r}^{l}=\inf _{x \in \Lambda_{s, r}^{l}} \delta(x)>0, \quad \delta(x)<1 \tag{2.2.6}
\end{gather*}
$$

7. $f(V(x)) \cap U(f(x), \delta(f(x))) \subseteq V(f(x))$.
8. There exists a measurable function $G(x), x \in \hat{\Lambda}$, such that

$$
G(f(x))=G(x), G_{s, r}=\sup _{x \in \Lambda_{s . r}} G(x)<\infty
$$

and if $y \in V(x)$, then

$$
\begin{equation*}
d\left(T_{y} V(x), T_{x} V(x)\right) \leqslant G(x) A^{2}(x) \rho(x, y) \tag{2.2.7}
\end{equation*}
$$

Proof. Let $x \in \hat{\Lambda}$. Choose a number $r_{0}$ and a mapping $\phi$ in accordance with Theorem 2.1.1 and Proposition 2.2.1. By Theorem 2.1.2, $\phi \in C^{r-1}$. The number $r_{0}$ and the mapping
$\phi$ depend on the choice of $x$, so $r_{0}=r_{0}(x)$ and $\phi=\phi(x)$. Put $\delta(x)=r_{0}(x) A^{-1}(x)$. The inequalities (2.2.6) are a corollary of (2.1.11), (2.2.1), (1.5.4) and (1.5.3). Assertion 1 is obvious, and 2 and 3 follow from (2.1.8). We will prove 4. It is known (see [10]) that there exist $\delta, C_{1}$ and $C_{2}$ such that for any $x, y \in M$ for which $\rho(x, y) \leqslant \delta$ we have

$$
\begin{equation*}
C_{1}\left\|\exp _{x}^{-1} y\right\|_{x} \leqslant \rho(x, y) \leqslant C_{2}\left\|\exp _{x}^{-1} y\right\|_{x} \tag{2.2.8}
\end{equation*}
$$

From (1.5.6) it follows that

$$
\begin{equation*}
A^{-1}(x) C_{1}\left\|\exp _{x}^{-1} y\right\|_{x}^{\prime} \leqslant \rho(x, y) \leqslant \sqrt{2} C_{2}\left\|\exp _{x}^{-1} y\right\|_{x}^{\prime} \tag{2.2.9}
\end{equation*}
$$

It follows from the above and (2.1.10) that

$$
\rho\left(f^{n}(x), f^{n}(y)\right) \leqslant \sqrt{2} C_{2}\left\|\left(\prod_{i=0}^{n} f_{i}\right)\left(u, \varphi_{i}(u)\right)\right\|_{f^{n}(x)}
$$

$$
\leqslant 200 \sqrt{2} C_{2} x^{n}(x)\|(u, \varphi(u))\|_{x}^{\prime} \leqslant 200 \sqrt{2 C_{2} C_{1}^{-1} A(x) x^{n}(x) \rho(x, y) . . . ~}
$$

Now 4 follows from 2 of Theorem 2.1.1. Assertion 5 is deduced from 3 of Theorem 2.1.1 in a similar way. Let $y \in f(V(x)) \cap U\left(f(x), \delta(f(x))\right.$ ); then $f^{-1}(y) \in V(x)$, and, by assertion $4, f^{n}(y)=f^{n+1}\left(f^{-1}(y)\right) \in U_{f^{n+1}(x)}^{\prime}$ and

$$
\begin{gathered}
\rho\left(f^{n}(x), f^{n}\left(f^{-1}(y)\right)=\rho\left(f^{n-1}(f(x)), f^{n-1}(y)\right)\right. \\
\leqslant N A(x) x^{n}(x) \rho\left(x, f^{-1}(y)\right)=N A(x) x(x) x^{n-1}(x) \rho\left(x, f^{-1}(y)\right) .
\end{gathered}
$$

Since $y \in U(f(x), \delta(f(x)))$, by 5 , in which we must put $C=N A(x) \kappa(x) \rho\left(x, f^{-1}(y)\right)$, we obtain $y \in V(f(x))$. Thus 7 is proved. From (2.1.9) and (2.1.8) it follows that for any $v \in S$

$$
\begin{equation*}
\|d \varphi v\|_{0}^{\prime} \leqslant C(x)\|v\|_{0}^{\prime} . \tag{2.2.10}
\end{equation*}
$$

Here $C(x)$ is the measurable function on $\hat{\Lambda}$ given by (2.1.11). Therefore $C(x)=A(x) L(x)$, where $L(x)$ is a measurable function constant on the trajectories of the diffeomorphism $f$. From (2.1.11), (2.1.9), (2.2.10), (2.2.1) and (1.5.6) it follows that for any $y \in V(x)$

$$
\begin{aligned}
& d\left(T_{y} V(x), T_{x} V(x)\right) \leqslant \sqrt{2} C_{2} \max _{\|v\|_{0} \leqslant\| \| \|_{0}}\|d \varphi(v)\|_{0}^{\prime} \\
& \leqslant \sqrt{2} C_{2} C(x) \max _{\|v\|_{0} \leqslant\|u\|_{0}}\|v\|_{0}^{0_{0}} \leqslant \sqrt{2} C_{2} C(x) A(x) \max _{\|p\|_{0} \leqslant\|x\|_{0}}\|v\|_{0} \\
& \leqslant \sqrt{2} C_{2} C(x) A(x)\left\|\exp _{x}^{-1} y\right\|_{x} \leqslant \sqrt{2} C_{2} C_{1}^{-1} L(x) A^{2}(x) \rho(x, y) .
\end{aligned}
$$

The theorem is proved.
Definition 2.2.1. The submanifold $V(x)$ is called the local stable manifold passing through $x$.
2.3. Some additional properties of local stable manifolds are described in the following result.

THEOREM 2.3.1. 1. If $x \in \Lambda_{s, r}^{l}, y \in \Lambda_{s, r}^{l_{1}}, l_{1} \geqslant l, y \in U\left(x, 1 / 4 \delta_{s, r}^{l}\right)$ and $y \in V(x)$, then

$$
V(x) \cap V(y) \cap U\left(y, \frac{1}{4} \delta_{s, r}^{l_{1}}\right)=\varnothing .
$$

2. If $x \in \Lambda_{s, r}^{l}, y \in V(x) \cap \Lambda_{s, r}^{l_{1}}$ and $l_{1} \geqslant l$, then

$$
V(y) \cap U\left(x, \delta_{s, r}^{i}\right) \subseteq V(x) .
$$

3. If $x \in \Lambda_{s, r}^{l}, x_{i} \in \Lambda_{s, r}^{l}, i=1,2, \ldots$, and $x_{i} \rightarrow x$, then $V\left(x_{i}\right) \rightarrow V(x)$ in the $C^{1}$. topology in a neighborhood $U(x, q)$, where $0<q<\delta_{s, r}^{l}$.

Proof. 1. Let $z \in V(x) \cap V(y) \cap U\left(y, 1 / 4 \delta_{s, r}^{1}\right)$. Then $z \in U\left(x, \delta_{s, r}^{l}\right)$, and for every $\epsilon>0$ there are points $z_{1}, z_{2} \in U\left(x, \delta_{s, r}^{l}\right)$ such that $p\left(z_{1}, z_{2}\right) \leqslant \epsilon, z_{1} \in V(x) \cap V(y), z_{2} \in$ $V(y), z_{2} \notin V(x)$. On the basis of Theorem 2.2.1 (see 4), for every $n \in \mathbf{Z}^{+}$we have

$$
\begin{gathered}
\rho\left(f^{n}(x), f^{n}\left(z_{1}\right)\right) \leqslant N A(x) x^{n}(x) \rho\left(x, z_{1}\right), \\
\rho\left(f^{n}\left(z_{1}\right), f^{n}\left(z_{2}\right)\right) \leqslant N A(y) x^{n}(y)\left(\rho\left(z_{1}, y\right)+\rho\left(z_{2}, y\right)\right) .
\end{gathered}
$$

Put $C=N A_{s, r}^{l_{1}}\left(\rho\left(z_{1}, y\right)+\rho\left(z_{2}, y\right)+\rho\left(z_{1}, x\right)\right)$. From the latter two inequalities, the triangle inequality and (1.5.4), it follows that for any $n \in \mathbf{Z}^{+}$

$$
\begin{equation*}
\rho\left(f^{n}(x), f^{n}\left(z_{2}\right)\right) \leqslant C\left(x_{s, r}\right)^{n} . \tag{2.3.1}
\end{equation*}
$$

Choose an integer $n_{0}$ so that for all $n \geqslant n_{0}$

$$
\begin{equation*}
C\left(\chi_{s, r}\right)^{n} \leqslant \delta_{s, r}^{l}, e^{-10 \varepsilon_{s}{ }^{n}} . \tag{2.3.2}
\end{equation*}
$$

This can be done because $\kappa_{s, r}<e^{-10 \epsilon_{s} n}<1$. Choosing $\epsilon$ sufficiently small and using (2.3.1) and (2.3.2), we conclude that $f^{n}\left(z_{2}\right) \in U_{f}^{\prime}{ }^{n(x)}$ for all $n \in \mathbf{Z}^{+}$. Therefore from Theorem 2.2.1 (see 5) it follows that $z_{2} \in V(x)$, which contradicts the choice of $z_{2}$. Assertion 2 is proved similarly.
3. By Theorem 1.3.1 (see 3), $E_{1 x_{i}} \rightarrow E_{1 x}$. Therefore from Theorem 2.2.1 it follows that for any $q, 0<q<\delta_{s, r}^{l}$, for sufficiently large $i, V\left(x_{i}\right)$ has the form

$$
V\left(x_{i}\right)=\left\{\exp _{x}\left(u, \chi_{i}(u)\right):\|u\|_{x} \leqslant q\right\},
$$

where $\chi_{i}(u)$ is a mapping of class $C^{r-1}$ of the ball of radius $q$ in $E_{1 x}$ into the space $E_{2 x}$. From (2.2.9) and 8 of Theorem 2.2 .1 follows the existence of a constant $C_{s, r}^{l}$ such that for sufficiently large $i$ and any $u_{1}, u_{2} \in E_{1 x},\left\|u_{1}\right\|_{x} \leqslant q,\left\|u_{2}\right\|_{x} \leqslant q$,

$$
\begin{gathered}
\sup _{\| \| \|_{x} \leqslant q}\left\|d \chi_{i}(u)\right\|_{x} \leqslant C_{s, r}^{l}, \\
\| d X_{i}\left(u_{1}\right)-d X_{i}\left(u_{2}\left\|_{x} \leqslant C_{s, r}^{l}\right\| u_{1}-u_{2} \|_{x} .\right.
\end{gathered}
$$

Therefore the family of functions $\chi_{i}$ is compact in the $C^{1}$-topology. Let $\psi$ be a limit point of the sequence $\chi_{i}$ and $\chi_{i_{p}} \rightarrow \psi$. Since $\chi_{i_{p}}(0) \rightarrow 0$ and $d \chi_{i_{p}}(0) \rightarrow 0$, we have $\psi(0)=0$ and $d \psi(0)=0$. Fix $m \in \mathbf{Z}^{+}$and $u \in E_{1 x},\|u\|_{x} \leqslant q$. Put

$$
z_{i_{p}}=\exp _{x}\left(u, \chi_{i_{p}}(u)\right), \quad z=\exp _{x}(u, \psi(u)) .
$$

Since $z_{i_{p}} \rightarrow z$, for every $\epsilon>0$ for sufficiently large $i_{p}$ we have

$$
\begin{gathered}
\rho\left(f^{m}\left(z_{i_{p}}\right), f^{m}(z)\right) \leqslant \varepsilon, \quad \rho\left(f^{m}\left(x_{i_{p}}\right), f^{m}(x)\right) \leqslant \varepsilon, \\
\rho\left(z_{i_{p}}, z\right) \leqslant \varepsilon, \quad \rho\left(x_{i_{p}}, x\right) \leqslant \varepsilon .
\end{gathered}
$$

Since $z_{i_{p}} \in V\left(x_{i_{p}}\right)$, (2.2.4) and (1.5.4) imply

$$
\rho\left(f^{m}\left(z_{i_{p}}\right), f^{m}\left(x_{i_{p}}\right)\right) \leqslant N\left(x_{s . r}\right)^{m} A_{s, r}^{l} \rho\left(z_{i_{p}}, x_{i_{p}}\right) .
$$

From the above and the triangle inequality it follows that

$$
\rho\left(f^{m}(z), f^{m}(x)\right) \leqslant N A_{s, r}^{l}\left(\varkappa_{s, r}\right)^{m} \rho(x, z)
$$

Since $m$ is arbitrary, from the latter inequality and Theorem 2.2.1 (see assertion 5), we obtain

$$
\exp _{x}\{(u, \psi(u))\}=V(x), \quad\|u\|_{x} \leqslant q .
$$

Thus $\chi_{i} \rightarrow \psi$ in the $C^{1}$-topology. The theorem is proved.
Remark 2.3.1. Without loss of generality, on the strength of Theorems 2.2.1 and 2.3.1 we may suppose that the number $\delta_{s, r}^{l}$ is so small that for any $x, y \in \Lambda_{k, s, r}^{l}$, $y \in U\left(x, \delta_{s, r}^{l} / 8\right)$,

$$
V(y) \cap U\left(x, \delta_{\mathrm{s} . r}^{l}\right) \supset\left\{\exp _{x}\left(u, \varphi_{y}(u)\right): u \in B\left(\frac{1}{4} \delta_{\mathrm{s}, r}^{l}\right)\right\},
$$

where $\phi_{y}: B\left(\frac{1}{4} \delta_{s, r}^{l}\right) \rightarrow E_{2 x}$ is a mapping of class $C^{r-1}$, with

$$
\max _{y \in A_{s, r}^{l} \cap U\left(x, \frac{1}{8} \delta_{s, r}^{l}\right)} \max _{u \in B\left(\frac{1}{4} \delta_{s, r}^{l}\right)}\left[\left\|\varphi_{y}(u)\right\|+\left\|d \varphi_{y}(u)\right\|\right] \leqslant 1 .
$$

Definition 2.3.1. Let $x \in \Lambda_{k, s, r}^{l}$. The collection of local stable manifolds passing through points $y \in \Lambda_{k, s, r}^{l} \cap U\left(x, \delta_{s, r}^{l} / 8\right)$ is called the family of local stable manifolds $S_{k, s, r}^{I}(x)$.

In conclusion we give one further result on local stable manifolds. Let $x \in \hat{\Lambda}$ and $y \in V(x)$. Since the trajectories of the points $x$ and $y$, under the action of $f^{n}$, $n=0,1, \ldots$, come together with exponential speed, the variational equation along the trajectory of $y$ can be considered as obtained from the variational equation along the trajectory of $x$ by a rapidly decaying (with increasing time) perturbation of the first order. From this and Theorems 15.2.1 and 17.1.1 of [4] follows

Proposition 2.3.1. 1. Let $x \in \tilde{\Lambda}$ and $y \in V(x)$. Then $y$ is fonward regular, and

$$
s(x)=s(y), \quad \chi_{i}(x)=\chi_{i}(y), \quad i=1, \ldots, s(x)
$$

where $\chi_{i}(x)$ is the value of the multiplier $\chi^{+}$at $x$ (see §0.3).
2. Let $x \in \widetilde{\Lambda}_{s, r}^{l}$ and $y \in \widetilde{\Lambda} \cap V(x)$. There exists $K=K(l, s, r)$ such that for any $n \in \mathbf{Z}^{+}$

$$
\begin{align*}
& \left\|d f^{n} v\right\|_{f n_{(y)}} \leqslant K\left(\frac{r}{s}\right)^{n} e^{5 \varepsilon_{s} \pi}\|v\|_{y}, \quad v \in E_{1 y}=T_{y} V(x),  \tag{2.3.3}\\
& \left\|d f^{n} v\right\|_{f^{n}(y)} \geqslant K^{-1}\left(\frac{r+2}{s}\right)^{n} e^{-5 \varepsilon_{s} n}\|v\|_{y}, \quad v \in E_{2 y} . \tag{2.3.4}
\end{align*}
$$

3. Let $x \in \hat{\Lambda}$ and $y \in V(x)$. Then for every $x \in T_{y} V(x)$

$$
\chi^{+}(y, v)<0 .
$$

## §3. Absolute continuity of families of local stable manifolds

3.1. Let $\Lambda$ be the set of positive measure, invariant relative to $f$, satisfying (0.2). We consider a set $\Lambda_{k, s, r}^{l}$ (see §1.3) having positive measure. For $y \in \Lambda_{k, s, r}^{l}$ we denote by $V(y)$ the local stable manifold in the neighborhood $U\left(y, \delta_{s, r}^{l}\right)$ of $y$ (see Theorem 2.2.1). Let $x$ be a point of density of $\Lambda_{k, s, r}^{l}$. We consider the family $S_{k, s, r}^{l}(x)$ of local stable manifolds (see (2.3.1)). Choose $q, 0<q \leqslant \delta_{s, r}^{l} / 8$, and put

$$
\begin{equation*}
\hat{\Lambda}_{k, s, r}^{l}(x)=\bigcup_{y \in \Delta_{k, 5, r}^{l} \cap \bar{U}(x, q)} \bar{V}(y) \cap \bar{U}(x, q) . \tag{3.1.1}
\end{equation*}
$$

It follows from Theorem 1.3.1 that $\hat{\Lambda}_{k, s, r}^{l}$ is closed in $\bar{U}(x, q)$. Consider a local open smooth submanifold $W \subset U(x, q)$ such that the set $\exp _{x}^{-1} W$ is the graph of a smooth mapping $\psi: U \rightarrow E_{1 x}$ defined in some neighborhood $U \subset E_{2 x}$ by the relation

$$
\begin{equation*}
t(\psi(u))=u, \quad u \in U \tag{3.1.2}
\end{equation*}
$$

where $t$ is the projection onto $E_{2 x}$ parallel to $E_{1 x}$. Put

$$
\begin{equation*}
|W|=\max _{u \in U}\|\psi(u)\|_{x}+\max _{u \in U}\|d \psi(u)\|_{x} \tag{3.1.3}
\end{equation*}
$$

It is easy to see that there is a constant $\epsilon_{k, s, r}^{l}>0$ such that if

$$
\begin{equation*}
|W| \leqslant \varepsilon_{k, s, r}^{l} \tag{3.1.4}
\end{equation*}
$$

then the submanifold $W$ intersects each local stable manifold $V(y), y \in \Lambda_{k, s, r}^{l} \cap U(x, q)$, in not more than one point; moreover, the intersection is transversal. A submanifold $W$ satisfying the above conditions (see (3.1.2) and (3.1.4)) will be called a transversal to the family $S_{k, s, r}^{l}(x)$.

Let $W^{1}$ and $W^{2}$ be two smooth submanifolds, transversals to the family $S_{k, s, r}^{l}(x)$. There exist open subsets $\widetilde{W}^{1} \subset W^{1}$ and $\widetilde{W}^{2} \subset W^{2}$ for which the successor mapping

$$
p: \hat{\Lambda}_{k, s, r}^{l} \cap \widetilde{W}^{1} \rightarrow \hat{\Lambda}_{k, s, r}^{l} \cap \widetilde{W}^{2}
$$

is defined. Namely if $y \in \widetilde{W}^{1} \cap V(w)$ and $w \in U(x, q) \cap \Lambda_{k, s, r}^{\prime}$, then

$$
\begin{equation*}
p(y)=\widetilde{W}^{2} \cap V(w) . \tag{3.1.5}
\end{equation*}
$$

Theorem 2.3.1 implies that $p$ is a homeomorphism.
Definition 3.1.1. The family $S_{k, s, r}^{l}(x)$ is called absolutely continuous if any successor mapping, constructed as above, is absolutely continuous.
3.2. Theorem 3.2.1. There exist constants $q_{k, s, r}^{l}$ and $J_{k, s, r}^{l}$ satisfying the following conditions:

1. The family $S_{k, s, r}^{l}(x)$ in the neighborhood $U\left(x, q_{k, s, r}^{l}\right)$ is absolutely continuous.
2. If $y$ is a point of density of $\hat{\Lambda}_{k, s, r}^{l} \cap \widetilde{W}^{1}$, then the Jacobian $J(p)(y)$ satisfies

$$
\begin{equation*}
|J(p)(y)-1| \leqslant J_{k, s, r}^{l}\left(\left|W^{1}\right|+\left|W^{2}\right|\right) \tag{3.2.1}
\end{equation*}
$$

Proof. Consider a point $w \in \Lambda_{k, s, r}^{l} \cap U(x, q), 0<q \leqslant \delta_{s, r}^{l} / 8$, a family of trivializations $\tau_{m}=\tau_{f(w)}^{\prime}: T_{f}^{m(w)} M \rightarrow \mathbf{R}^{n}$, and the mappings $f_{m}=f_{f m-1(w)}$ defined in the neighborhoods $U_{m}=U_{f m(w)} \subset T_{f^{m}(w)} M$ (see $\S 1.6$ ). We denote by $\|\cdot\|_{m}^{\prime}$ the measurable Ljapunov norm in $T_{f^{m}(w)} M$, by $\|\cdot\|_{m}$ the Riemannian norm in $T_{f(w)} M$, and by $\|\cdot\|$ the norm in $\mathbf{R}^{n}$. We define the numbers $\lambda, \mu, \nu, \kappa, \alpha$ and $\epsilon_{s}$ by (2.2.1). Let

$$
\varphi\left(f^{m}(w)\right): B^{1}\left(\delta\left(f^{m}(w)\right)\right) \rightarrow B^{2}\left(\delta\left(f^{m}(w)\right)\right)
$$

be the smooth mapping constructed in Theorem 2.2.1. We put $\phi_{m}=\phi\left(f^{m}(w)\right)$ and

$$
B_{m}^{i}=\left\{v \in E_{i f m_{(w)}}:\|v\|_{m} \leqslant q_{m}\right\}, \quad i=1,2, \quad B_{m}=B_{m}^{1} \times B_{m}^{2}
$$

where $q_{m}=q_{0} e^{-10 \epsilon_{s} m}, q_{0} \leqslant \delta(w)$. By Theorem 2.2.1, $B_{m}^{i} \subset B^{i}\left(\delta\left(f^{m}(w)\right)\right), i=1,2$.
Consider a point $P=\left(u_{0}, v_{0}\right) \in B_{0}$ on the graph of the function $\phi_{0}, v_{0}=\phi_{0}\left(u_{0}\right)$, and put, for $m \in \mathbf{Z}^{+}$,

$$
P_{m}=\left(u_{m}, v_{m}\right)=f_{m}(P), \quad u_{m} \in B_{m}^{1}, \quad v_{m} \in B_{m}^{2}
$$

Theorem 2.2.1 (see assertion 7) implies that $v_{m}=\phi_{m}\left(u_{m}\right)$.
Put

$$
\xi=\frac{r}{s} e^{\theta \varepsilon_{s}}, \quad \eta=e^{-\varepsilon \varepsilon_{s}} .
$$

It is not difficult to verify, using (1.3.1) and (2.2.1), that

$$
\begin{equation*}
x<\xi<\mu, \quad \frac{\lambda}{\mu}<\eta, \quad \xi<\eta<v, \quad \xi v^{-1}<\eta . \tag{3.2.2}
\end{equation*}
$$

Fix a $\delta_{0}>0$, and write $\delta_{m}=\delta_{0} \xi^{m}$ and

$$
B\left(v_{m}, \delta_{m}\right)=\left\{v \in E_{2 f^{m l}(m)}:\left\|v-v_{m}\right\|_{m} \leqslant \delta_{m}\right\}
$$

Lemma 3.2.1. If $\left\|v_{0}\right\|_{0} \leqslant q_{0}\left(200 \sqrt{2} A_{s, r}^{l}\right)^{-1}$ and $\delta_{0} \leqslant q_{0} / 2$, then $B\left(v_{m}, q_{m}\right) \subset B_{m}^{2}$.
Proof. If $v \in B\left(v_{m}, \delta_{m}\right)$, then from (1 4$)$, (1.5.6), (2.1.10), the conditions of the lemma and (3.2.2) it follows that

$$
\begin{gathered}
\|v\|_{m} \leqslant\left\|v_{m}\right\|_{m}+\delta_{m} \leqslant \sqrt{2}\left\|v_{m}\right\|_{m}^{\prime}+\delta_{m} \\
\leqslant 200 \sqrt{2} x^{m}\left\|v_{0}\right\|_{0}^{\prime}+\delta_{m} \leqslant 200 \sqrt{2} A_{s, r}^{l} x^{m}\left\|v_{0}\right\|_{0}+\delta_{0} \xi^{m} \leqslant q_{0} \xi^{m}
\end{gathered}
$$

The lemma is proved.

Lemma 3.2.2. Let $\psi_{0}: B\left(v_{0}, \delta_{0}\right) \rightarrow E_{10}$ be a smooth mapping, where

$$
\begin{equation*}
\max _{v \in B\left(v_{0}, \delta_{0}\right)}\left\|\psi_{0}(v)\right\|_{0}^{\prime}=C^{\prime}, \quad \max _{v \in B\left(v_{00} \delta_{0}\right)}\left\|d \psi_{0}(v)\right\|_{0}^{\prime}=C \tag{3.2.3}
\end{equation*}
$$

Then there exists a $\hat{q}=\hat{q}(l, s, r, C) \leqslant \delta_{s, r}^{l} / 8$ such that if $q_{0} \leqslant \hat{q}$ and $\delta_{0} \leqslant q_{0} / 2$, then there is a sequence of smooth mappings $\psi_{m}: B\left(v_{m}, \delta_{m}\right) \rightarrow E_{1 f m(w)}, m>0$, for which

$$
\begin{gather*}
\left\{\left(\psi_{m}(v), v\right): v \in B\left(v_{m}, \delta_{m}\right)\right\} \subset\left\{f_{m}\left(\psi_{m-1}(v), v\right): v \in B\left(v_{m-1}, \delta_{m-1}\right)\right\}  \tag{3.2.4}\\
a_{m}=\max _{v \in B\left(v_{m}, \delta_{m}\right)}\left\|\psi_{m}(v)\right\|_{m}^{\cdot} \leqslant\left(200+C A_{s, r}^{l}\right) q_{0} \xi^{m}  \tag{3.2.5}\\
b_{m}=\max _{v \in B\left(v_{m}, \delta_{m}\right)}\left\|d \psi_{m}(v)\right\|_{m}^{\cdot} \leqslant C \eta^{m} \tag{3.2.6}
\end{gather*}
$$

Proof. Denote

$$
\begin{gather*}
C_{2}=C, \quad C_{1}=200+C A_{s, r}^{l}  \tag{3.2.7}\\
\hat{q}=(C+1) M A_{s, r}^{l}\left(C_{1}+A_{s, r}^{l}\right) \min \{(\mu \eta-\lambda),(\mu-\xi)\}
\end{gather*}
$$

Here $M$ is the same as in (1.6.8). Lemma 3.2.1 implies that $\hat{q}>0$. It is easy to see that if $q_{0} \leqslant \hat{q}$, then

$$
\begin{equation*}
\mu>M A_{s, r}^{l}\left(C_{1}+A_{s, r}^{l}\right)(C+1) q_{a}+\xi \tag{3.2.8}
\end{equation*}
$$

For the proof of the lemma we proceed by induction. We suppose that we have constructed functions $\psi_{k}, k=0,1, \ldots, m$, satisfying the assertions of the lemma. Put

$$
\tilde{v}=\tau_{m} v, \quad \widetilde{B}_{m}=\tau_{m} B\left(v_{m}, \delta_{m}\right), \quad \tilde{\psi}_{m}(\tilde{v})=\tau_{m} \psi_{m}(v)
$$

It follows from (1.6.5) that for $v \in B\left(v_{m}, \delta_{m}\right)$

$$
\begin{gather*}
\tau_{m+1} f_{m+1}\left(\psi_{m}(\tilde{v}), v\right)=(\hat{u}, \hat{v})=\left(A_{m+1} \widetilde{\psi}_{m}(\tilde{v})+g_{m+1}\left(\tilde{\psi}_{m}(\tilde{v}), \tilde{v}\right),\right.  \tag{3.2.9}\\
\left.B_{m+1} \widetilde{v}+h_{m+1}\left(\tilde{\psi}_{m}(\widetilde{v}), \tilde{v}\right)\right)
\end{gather*}
$$

where

$$
A_{m+1}=A_{f^{m}(w)}, \quad B_{m+1}=B_{f m_{(w)}}, \quad g_{m+1}=g_{f m_{(w)}}, \quad h_{m+1}=h_{f m_{(w)}}
$$

Denote by $t_{m}: \widetilde{B}_{m} \rightarrow \mathbf{R}^{n-k}$ the mapping associating the vector $\widetilde{\sigma}$ to the vector $\hat{v}$ (see (3.2.9)). We will show that

$$
\begin{equation*}
\left\|t_{m+1} \tilde{v}_{1}-t_{m+1} \tilde{v}_{2}\right\| \geq \xi\left\|\tilde{v}_{1}-\tilde{v}_{2}\right\| . \tag{3.2.10}
\end{equation*}
$$

From the mean value theorem it follows that for $v_{1}, v_{2} \in B\left(v_{m}, \delta_{m}\right)$

$$
\begin{equation*}
\left\|h_{m+1}\left(\tilde{\psi}_{m}\left(\tilde{v}_{1}\right), \tilde{v}_{1}\right)-h_{m+1}\left(\tilde{\psi}_{m}\left(\tilde{v}_{2}\right), \tilde{v}\right)\right\| \leqslant \sup \left\|d h_{m+1}(z)\right\|\left(\left\|\tilde{\psi}_{m}\left(\tilde{v}_{1}\right)-\tilde{\psi}_{m}\left(\tilde{v}_{2}\right)\right\|+\left\|\tilde{v}_{1}-\tilde{v}_{2}\right\|\right) \tag{3.2.11}
\end{equation*}
$$

Here the supremum is taken over all points $z$ lying on the segment joining the points $\left(\tilde{\psi}_{m}\left(\tilde{v}_{1}\right), \tilde{v}_{1}\right)$ and $\left(\widetilde{\psi}_{m}\left(\tilde{v}_{2}\right), \tilde{v}_{2}\right)$. By (1.5.3), (1.5.4) and (1.6.8) we have

$$
\begin{equation*}
\left\|d h_{m+1}(z)\right\|=\left\|d h_{m+1}(z)-d h_{m+1}(0)\right\| \leqslant M A_{s, r}^{l} v^{-m}\|z\| . \tag{3.2.12}
\end{equation*}
$$

Since $\psi_{m}$, by the induction hypothesis, satisfies (3.2.5) and (3.2.6), it follows that (3.2.7) and the inequalities $(1.5 .3),(1.5 .4)$ and (1.5.6) imply the estimates

$$
\begin{gather*}
\left\|\widetilde{\psi}_{m}\left(\tilde{v}_{1}\right)-\widetilde{\psi}_{m}\left(\widetilde{v}_{2}\right)\right\| \leqslant C \eta^{m}\left\|\tilde{v}_{1}-\widetilde{v}_{2}\right\| \\
\|z\| \leqslant \max _{i=1,2}\left\{\left\|\widetilde{\psi}_{m}\left(\widetilde{v}_{i}\right)\right\|+\left\|\tilde{v}_{i}\right\|\right\} \leqslant C_{1} q_{0} \xi^{m}+A_{s, r}^{l} v^{-m} \delta_{0} \xi^{m} \tag{3.2.13}
\end{gather*}
$$

Therefore (3.2.11)-(3.2.13) imply that

$$
\begin{equation*}
h_{m+1}\left(\tilde{\psi}_{m}\left(\tilde{v_{1}}\right), \tilde{v}_{1}\right)-h_{m+1}\left(\tilde{\psi}_{m}\left(\tilde{v_{2}}\right), \tilde{v_{2}}\right) \| \leqslant M A_{s, r}^{l}\left(C_{1}+A_{s, r}^{l}\right)(C+1) q_{0}\left(\eta \xi v^{-2}\right)^{m} \tag{3.2.14}
\end{equation*}
$$

On the other hand, by (1.6.6)

$$
\begin{equation*}
\left\|B_{m+1}\left(\tilde{v}_{1}-\tilde{v}_{2}\right)\right\| \geqslant \mu\left\|\tilde{v}_{1}-\tilde{v}_{2}\right\| \tag{3.2.15}
\end{equation*}
$$

Therefore (3.2.2) and (3.2.8) imply (3.2.10). Since $\delta_{m+1}=\xi \delta_{m}$, we have

$$
\begin{equation*}
f_{m}\left(B_{m}\left(v_{m}, \delta_{m}\right)\right) \supseteq B\left(v_{m+1}, \delta_{m+1}\right) \tag{3.2.16}
\end{equation*}
$$

Hence, in turn, it follows that the mapping $\psi_{m+1}(v)=\tau_{m+1}^{-1} t_{m+1} \tau_{m}(v)$ satisfies (3.2.4).
On the basis of (2.1.10), we obtain

$$
\begin{equation*}
a_{m+1} \leqslant\left\|\tilde{\psi}_{m+1}\left(\tilde{v}_{m+1}\right)\right\|+\max _{v \in \tilde{B}_{m}}\left\|d \tilde{\psi}_{m+1}(v)\right\| \tilde{v} \| \leqslant 200 q_{0} x^{m+1}+b_{m+1} g_{0} \xi^{m+1} A_{s, r}^{l} \nu^{-(m+1)} \tag{3.2.17}
\end{equation*}
$$

Choose $\hat{\tau}$, near to zero, such that $\hat{v}+\hat{\tau} \in \widetilde{B}_{m+1}$. There exists a unique $\tilde{\tau} \in \mathbf{R}^{n-k}$ such that

$$
\hat{v}+\hat{\tau}=B_{m+1}(\tilde{v}+\tilde{\tau})+h_{m+1}\left(\tilde{\psi} \psi_{m}(\tilde{v}+\tilde{\tau}), \tilde{v}+\tilde{\tau}\right)
$$

Using (3.2.9), we obtain that

$$
\tilde{\tau}=B_{m+1} \tilde{\tau}+h_{m+1}\left(\tilde{\psi_{m}}(\tilde{v}+\tilde{\tau}), \tilde{v}+\tilde{\tau}\right)-h_{m+1}\left(\tilde{\psi}_{m}(\tilde{v}), \tilde{v}\right)
$$

From (3.2.11)-(3.2.15), in which we must put $\tilde{v}_{1}=\tilde{v}+\tilde{\tau}$ and $\tilde{v}_{2}=\tilde{v}$, it follows that

$$
\|\hat{\tau}\|>\left(\mu-M A_{s, r}^{l}\left(C_{1}+A_{\mathrm{s}, r}^{l}\right)(C+1) q_{0}\left(\eta v^{-2}\right)^{m}\right)\|\tilde{\tau}\|
$$

Again using (3.2.9), we obtain that

$$
\begin{gather*}
\left\|\tilde{\psi}_{m+1}(\tilde{v}+\tilde{\tau})-\widetilde{\psi}_{m+1}(\tilde{v})\right\| \leqslant\left\|A_{m+1}\right\|\left\|\widetilde{\psi}_{m}(\tilde{v}+\tilde{\tau})-\widetilde{\psi}_{m}(\widetilde{v})\right\| \\
+\left\|g_{m+1}\left(\widetilde{\psi}_{m}(\tilde{v}+\tilde{\tau}), \tilde{v}+\tilde{\tau}\right)-g_{m+1}\left(\tilde{\psi}_{m}(\tilde{v}), \tilde{v}\right)\right\| \tag{3.2.18}
\end{gather*}
$$

Since inequalities similar to (3.2.11) and (3.2.12) hold for $g_{m+1}$, on dividing (3.2.18) by $\|\hat{\tau}\|$, passing to the limit as $\hat{\tau} \rightarrow 0$ and taking the least upper bound over all $\widetilde{v} \in \widetilde{B}_{m}$, we obtain

$$
b_{m+1} \leqslant \frac{\lambda b_{m}+M A_{s, r}^{l}\left(C_{1}+A_{s, r}^{l}\right)(C+1) q_{0}\left(\eta \xi v^{-2}\right)^{m}}{\mu-M A_{\mathrm{s}, r}^{l}\left(C_{1}+A_{s, r}^{l}\right)(C+1) q_{0}\left(\eta \xi v^{-2}\right)^{m}}
$$

Since, according to the induction hypothesis, $b_{m} \leqslant C_{2} \eta^{m}=C \eta^{m}$, (3.2.7) implies that $b_{m+1} \leqslant C \eta^{m+1}$. From this (3.2.17) and (3.2.2) it follows that

$$
a_{m+1} \leqslant\left(200 q_{0}+C A_{s, r}^{l} q_{0}\right)\left(\eta \xi v^{-1}\right)^{m+1} \leqslant C_{1} q_{0} \xi^{m+1}
$$

The lemma is proved.
From Lemma 3.2.2 and (1.5.3), (1.5.4) and (1.5.6) follows

$$
\begin{gather*}
\max _{v \in B\left(v_{m}, \delta_{m}\right)}\left\|\psi_{m}(v)\right\| \leqslant \sqrt{2} C_{1} q_{0} \xi^{m},  \tag{3.2.19}\\
\max _{v \in B\left(v_{m}, \delta_{\dot{m}}\right)}\left\|d \psi_{m}(v)\right\| \leqslant \sqrt{2} C_{2} A_{s_{, r}}^{l}\left(\eta v^{-1}\right)^{m} .
\end{gather*}
$$

Taking $C=1$ in (3.2.7), we define a number $\hat{q}$ depending only on the numbers $l, k, s$ and $r$. There exists a $q_{k, s, r}^{l}$ such that

$$
\begin{equation*}
q_{k, s, r}^{l} \leqslant\left(200 \sqrt{2} A_{s, r}^{l}\right)^{-1} \hat{q} \tag{3.2.20}
\end{equation*}
$$

and for any point $w \in U\left(x, q_{k, s, r}^{l}\right)$

$$
\begin{equation*}
\exp _{w} B_{s}\left(\left(200 \sqrt{2} A_{s, r}^{l}\right)^{-1} \hat{q}\right) \supset U\left(x, q_{k, s, r}^{l}\right) \tag{3.2.21}
\end{equation*}
$$

We consider in the neighborhood $U\left(x, q_{k, s, r}^{l}\right)$ a smooth submanifold $W$, transviersal to the family $S_{k, s, r}^{l}(x)$. Let $w \in \Lambda_{k, s, r}^{l} \cap U\left(x, q_{k, s, r}^{l}\right)$. We choose $q_{k, s, r}^{l}$ so small that

$$
\left.\begin{array}{l}
\text { the set } \exp _{w}^{-1} w \text { is the graph of a smooth mapping }  \tag{3.2.22}\\
\psi: U \rightarrow E_{1 w} \text { (see (3.1.2)), given in some } \\
\text { neighborhood } U \subset E_{2 w} .
\end{array}\right\}
$$

Thus we will suppose that $q_{k, s, r}^{l}$ is chosen in accordance with the conditions (3.2.20)(3.2.22). Without loss of generality we can assume that (see (2.3.3))

$$
\max _{v \in U}\|\psi(v)\|_{0} \leqslant 1, \quad \max _{v \in U}\|d \Psi(v)\|_{0} \leqslant 1 .
$$

Let $y=V(w) \cap W$. Put

Then

$$
\begin{gathered}
\exp _{w}^{-1} y=P=\left(u_{0}, v_{0}\right), \quad f_{m}(P)=P_{m}=\left(u_{m}, v_{m}\right), \quad m \in \mathbf{Z}^{+}, \\
q(w, W)=200 \sqrt{2} A_{s, r}^{l} \max _{v \in U}\|\psi(u)\| .
\end{gathered}
$$

$$
\begin{align*}
& q(w, W) \leqslant \frac{1}{8} \delta_{s, r}^{l}, \quad q(w, W) \leqslant \hat{q},  \tag{3.2.23}\\
& \left\|v_{0}\right\|_{0} \leqslant\left(200 \sqrt{2} A_{s, r}^{l}\right)^{-1} \quad q(w, W) .
\end{align*}
$$

Further, by (3.1.3), (3.2.3), (3.2.7), (1.5.6) and (3.2.21) there is a constant $C_{3}=C_{3}(l, s, r)$ such that for any $w \in \Lambda_{k, s, r}^{l} \cap U\left(x, q_{k, s, r}^{l}\right)$

$$
\begin{equation*}
C_{1} q(w, W)+C_{2} \leqslant C_{3}|W| . \tag{3.2.24}
\end{equation*}
$$

Taking $\delta_{0}=1 / 2 q(w, W)$ and applying Lemmas 3.2.1 and 3.2.2, we construct a mapping $\psi_{m}: B\left(v_{m}, \delta_{m}\right) \rightarrow E_{1 f m(w)}$. We denote

$$
\begin{equation*}
W_{m}=W_{m}\left(w, y, \delta_{0}\right)=\left\{\exp _{f^{m}(w)}\left(\psi_{m}(v), v\right): v \in B\left(v_{m}, \delta_{m}\right)\right\} \tag{3.2.25}
\end{equation*}
$$

It is obvious that $W_{m}$ is a smooth submanifold in $M$. By (3.2.4) and (3.2.16)

$$
\begin{equation*}
f^{-1}\left(W_{m}\right) \subset W_{m-1} \tag{3.2.26}
\end{equation*}
$$

We consider in the neighborhood $U\left(x, q_{k, s, r}^{l}\right)$ the set $\hat{\Lambda}_{k, s, r}^{l}$ (see (3.1.1)), and two smooth submanifolds $W^{1}$ and $W^{2}$, transversal to the family $S_{k, s, r}^{l}(x)$ and intersecting $V\left(w_{1}\right), w_{1} \in$ $\Lambda_{k, s, r}^{l} \cap U\left(q_{k, s, r}^{l}\right)$, in the points $y^{1}$ and $y^{2}$ respectively.

Let $W_{m}^{i}=W_{m}\left(w_{1}, y^{i}, \delta_{0}\right)$ be the submanifolds defined by (3.2.25), $i=1,2$. By (3.2.26)

$$
z^{1}=f^{-m}\left(W_{m}^{1} \cap f^{m}\left(\hat{\Lambda}_{k, s, r}^{l}\right)\right) \in W_{0}^{1} \cap \hat{\Lambda}_{k, s, r}^{l}
$$

Therefore there exists a point $w_{2} \in A_{k, s, r}^{l} \cap U\left(x, q_{k, s, r}^{l}\right)$ such that $z^{1} \in V\left(w_{1}\right)$. Put

$$
\begin{gathered}
P_{0}^{i}=\exp _{w_{1}}^{-1} y^{i}=\left(u_{0}^{i}, v_{0}^{i}\right), \quad P_{k}^{i}=f_{k}\left(P^{i}\right)=\left(u_{k}^{i}, v_{k}^{i}\right), \\
z^{2}=V\left(w_{2}\right) \cap W^{2}, \quad \exp _{w_{w}}^{-1} z^{i}=\left(u_{0}, v_{0}\right), \\
\left.\exp _{p^{k}\left(w_{1}\right)}\right)^{k}\left(z^{i}\right)=\left(\vec{u}_{k}^{i}, \bar{v}_{k}^{i}\right)=f_{k}\left(\bar{u}_{0}^{i}, \bar{v}_{0}^{i}\right), \\
i=1,2, \quad k=1, \ldots, m .
\end{gathered}
$$

Lemma 3.2.3. For every $\alpha>0$ there exists $m_{1}=m_{1}(l, s, r, \alpha)$ such that for any $m \geqslant m_{1}$

$$
\begin{equation*}
\bar{v}_{m}^{2} \in B\left(v_{m}^{2}, \bar{\delta}_{m}\right) \tag{3.2.27}
\end{equation*}
$$

where $\bar{\delta}_{m}=\left(\delta_{0}+\alpha\right) \xi^{m}$.
Proof. We will estimate $\left\|\bar{v}_{m}^{2}-v_{m}^{2}\right\|_{m}$. We have

$$
\begin{equation*}
\left\|\bar{v}_{m}^{2}-v_{m}^{2}\right\|_{m} \leqslant\left\|\bar{v}_{m}^{2}-\bar{v}_{m}^{1}\right\|_{m}+\left\|\bar{v}_{m}^{1}-v_{m}^{2}\right\|_{m} . \tag{3.2.28}
\end{equation*}
$$

Since $f^{m}\left(z^{1}\right) \in W_{m}^{1},(2.1 .10),(1.5 .4)$ and (1.5.6) imply

$$
\begin{gather*}
\left\|\bar{v}_{m}^{1}-v_{m}^{2}\right\|_{m} \leqslant\left\|\bar{v}_{m}^{1}-v_{m}^{1}\right\|_{m}+\left\|v_{m}^{1}-v_{m}^{2}\right\|_{m} \\
\leqslant \delta_{m}+\sqrt{2}\left(\left\|v_{m}^{2}\right\|_{m}^{m}+\left\|v_{m}^{2}\right\|_{m}^{\prime}\right) \leqslant \delta_{m}+200 \sqrt{2} A_{s, r}^{2}\left(x_{s, r}\right)^{m}\left(|W|+\left|W^{2}\right|\right) \tag{3.2.29}
\end{gather*}
$$

Since $z^{1}, z^{2} \in V\left(w_{2}\right)$, (2.2.8), (2.2.4) and (1.5.3) imply

$$
\begin{gather*}
\left\|\bar{v}_{m}^{2}-\bar{v}_{m}^{1}\right\|_{m} \leqslant C_{1}^{-1} \rho\left(f^{m}\left(z^{1}\right), f^{m}\left(z^{2}\right)\right) \\
\leqslant C_{1}^{-1} N\left(x_{s, r}\right)^{m} A_{s, r}^{l} \cdot\left(z^{2}, z^{2}\right) \leqslant C_{1}^{-1} N A_{s, r}^{1}\left(x_{s, r}^{-r}\right)^{m}\left(\left|W^{1}\right|+\left|W^{2}\right|\right) . \tag{3.2.30}
\end{gather*}
$$

Thus it follows from (3.2.28)-(3.2.30) and (3.2.2) that $\left\|\bar{v}_{m}^{2}-v_{m}^{2}\right\|_{m} \leqslant \delta_{m}$ if $m$ is so large that

$$
A_{s, r}^{l}\left(200 \sqrt{2}+C_{1}^{-1} N\right)\left(\chi_{s, r}\right)^{m}\left(\left|W^{1}\right|+\left|W^{2}\right|\right) \leqslant \alpha \xi^{m}
$$

The lemma is proved.
Lemma 3.2.4. For any $\alpha>0, \delta_{0} \leqslant 1 / 2\left(\left|W^{1}\right|+\left|W^{2}\right|\right)$ and $m \geqslant m_{1}(\alpha)$ (see Lemma 3.2.3)

$$
p\left(f^{-m}\left(W_{m}\left(w, y^{1}, \delta_{0}\right)\right) \cap \hat{\Lambda}_{k, s, r}^{l}\right) \subset f^{-m}\left(W_{m}\left(w, y^{2}, \delta_{0}+\alpha\right)\right)
$$

where $p$ is the successor mapping connected with the manifolds $W^{1}$ and $W^{2}$.
The proof follows from Lemmas 3.2.1, 3.2.3, the inclusion (3.2.26) and the inequality (3.2.23).

The last assertion shows that to compare the measures of the sets $A \subset \hat{\Lambda}_{k, s, r}^{l} \cap W^{1}$ and $p(A)$ we must learn to compare the parts of these sets lying in the inverse images of the submanifolds $W_{m}\left(w, y^{1}, \delta_{0}\right)$ and $W_{m}\left(w, y^{2}, \delta_{0}+\alpha\right)$ respectively.

Lemma 3.2.5 (see [3], Lemma 6). Let $z^{1}=f^{m}(z) \in W_{m}^{1}$. There exist constants $C_{5}=C_{5}(l, s, r)$ and $C_{6}=C_{6}(l, s, r)$ such that

$$
\begin{align*}
& \left.\left|\left|d f_{f^{m}(z)}^{-m}\right|_{T_{f^{m}(z)} W_{m}^{1}}\right| \cdot\left|d f_{f^{m}\left(y^{1}\right)}^{-m}\right|_{T_{f m_{( }\left(y^{1}\right)} W_{m}^{1}}\right|^{-1}-1 \mid \leqslant C_{5}\left(\left|W^{\mathbf{1}}\right|+\left|W^{2}\right|\right),  \tag{3.2.31}\\
& \left.\left|\left|d f_{f^{m}\left(y^{2}\right)}^{-m}\right|_{T_{m^{m}\left(y^{2}\right)} W_{m}^{2}}\right| \cdot\left|d f_{F^{m}\left(y^{1}\right)}^{-m}\right|_{T_{m_{\left(y^{1}\right)}} W_{m}^{1}}\right|^{-1}-1 \mid \leqslant C_{6}\left(\left|W^{1}\right|+\left|W^{2}\right|\right),
\end{align*}
$$

where $|B|$ denotes the coefficient of volume expansion for the mapping $B$.
Proof. Since $f$ is of class $C^{2}$, there exists a $C_{7}>0$ such that

$$
\left\|d f_{z_{2}}^{-1}-d f_{z_{2}}^{-1}\right\| \leqslant C_{7} \rho\left(z_{1}, z_{2}\right)
$$

where $z_{1}$ and $z_{2}$ are arbitrary points in $M$. Therefore there exists $C_{8}$ such that for any subspace $A$

$$
\begin{equation*}
\left|\left|d f_{z_{1}}^{-1}\right|_{A}\right|-\left|d f_{z_{2}}^{-1}\right|_{A}| | \leqslant C_{8} \rho\left(z_{1}, z_{2}\right) \tag{3.2.32}
\end{equation*}
$$

In addition, there exists a $C_{9}>0$ such that for any $z \in M$ and two subspaces $A_{1}, A_{2} \subset$ $T_{z} M$ of the same dimension

$$
\begin{equation*}
\left|\left|d f_{z}^{-1}\right|_{A_{1}}\right|-\left|d f_{z}^{-1}\right|_{A_{2}}| | \leqslant C_{9} d\left(A_{1}, A_{2}\right) \tag{3.2.33}
\end{equation*}
$$

Denote

$$
E_{z}^{k}=T_{f^{-k}(z)} W_{k}^{1}, \quad J(k, z)=\left|d f_{f^{-k}(z)}^{-1}\right|_{E_{z}^{k}} \mid
$$

It follows from (3.2.19), (3.2.24), (3.2.32) and (3.2.33) that

$$
\begin{align*}
& \left|J(k, z)-J\left(k, y_{1}\right)\right| \leqslant C_{8} \rho\left(f^{-k}(z), f^{-k}\left(y_{1}\right)\right)  \tag{3.2.34}\\
& +C_{9} d\left(E_{y_{1}}^{k}, E_{z}^{k}\right) \leqslant C_{10}\left(\left|W^{1}\right|+\left|W^{2}\right|\right)\left(\eta v^{-1}\right)^{k}
\end{align*}
$$

where $C_{10}=C_{10}(l, s, r)$ is a constant. Obviously there exists a $J_{0}>0$ such that

$$
\begin{equation*}
J_{0}^{-1} \leqslant J(k, z) \leqslant J_{0} \tag{3.2.35}
\end{equation*}
$$

The first inequality in (3.2.31) follows from (3.2.34) and (3.2.35). The second inequality follows similarly. The lemma is proved.

We denote by $\mu_{m}^{i}$ the Riemannian volume in $f^{m}\left(W^{i}\right)$ induced by the Riemannian metric of the manifold $M, i=1,2$.

Lemma 3.2.6. There exists $C_{11}=C_{11}(l, s, r)>0$ such that for any $m \in \mathbf{Z}^{+}, a>0$, $\delta_{0}, \delta_{1} \leqslant 1 / 2\left(\left|W^{1}\right|+\left|W^{2}\right|\right)$ and any two measurable sets $X^{1} \subset W_{m}^{1}\left(w, y^{1}, \delta_{0}\right)$ and $X^{2} \subset$ $W_{m}^{2}\left(w, y^{2}, \delta_{1}\right)$ for which $\left|\mu_{m}^{1}\left(X^{1}\right)\left(\mu_{m}^{2}\left(X^{2}\right)\right)^{-1}-1\right| \leqslant a$, we have

$$
\begin{equation*}
\left|\frac{\mu_{0}^{1}\left(f^{-m}\left(X^{1}\right)\right)}{\mu_{0}^{2}\left(f^{-m}\left(X^{2}\right)\right)}-1\right| \leqslant C_{11}\left(\left|W^{1}\right|+\left|W^{2}\right|+a\right) \tag{3.2.36}
\end{equation*}
$$

Proof. Denote

$$
A_{i}(z)=\left|d f_{z}^{-m}\right|_{T_{z} w_{m}^{i}} \mid, \quad i=1,2
$$

For $i=1,2$ we have

$$
\mu_{0}^{i}\left(f^{-m}\left(X^{i}\right)\right)=\int_{X^{i}} A_{i}(z) d \mu_{m}^{i}(z)=A_{i}\left(z_{i}\right) \mu_{m}^{i}\left(X^{i}\right)
$$

where $z_{i} \in X^{i}$. Therefore from the conditions of the lemma and (3.2.31) it follows that

$$
\begin{gathered}
\left|\frac{\mu_{0}^{1}\left(f^{-m}\left(X^{1}\right)\right)}{\mu_{0}^{2}\left(f^{-m}\left(X^{2}\right)\right)}-1\right| \leqslant\left|\frac{A_{1}\left(z_{1}\right)}{A_{2}\left(z_{2}\right)}-1\right| \frac{\mu_{m}^{1}\left(X^{1}\right)}{\mu_{m}^{2}\left(X^{2}\right)} \\
+\left|\frac{\mu_{m}^{2}\left(X_{1}\right)}{\mu_{m}^{2}\left(X^{2}\right)}-1\right| \leqslant\left|\frac{A_{1}\left(z_{1}\right)}{A_{1}\left(y^{1}\right)} \frac{A_{2}\left(y^{2}\right)}{A_{2}\left(z_{2}\right)} \frac{A_{1}\left(y^{1}\right)}{A_{2}\left(y^{2}\right)}-1\right| \\
\quad \times \frac{\mu_{m}^{1}\left(X^{1}\right)}{\mu_{m}^{2}\left(X^{2}\right)}+a \leqslant C_{\mathbf{u}}\left(\left|W^{1}\right|+\left|W^{2}\right|+a\right) .
\end{gathered}
$$

The lemma is proved.
We will now show how to cover the "nice" parts of $f^{m}\left(W^{1}\right)$, i.e. the sets $f^{m}\left(\hat{\Lambda}_{k, s, r}^{l} \cap W^{1}\right)$, by certain sets lying in the submanifolds $W_{m}\left(w_{j}, y^{i}, \delta_{0}\right)$ with centers at certain points $w_{j}$. This covering will not only have finite multiplicity not depending on $m$ (such a cover was constructed, for example, in [3]), but will be almost a decomposition; that is, the sum of the measures of the intersections of elements of the cover is sufficiently small by comparison with the value $\left|W^{1}\right|+\left|W^{2}\right|$. Here certain "superfluous" pieces of the submanifold $f^{m}\left(W^{1}\right)$ are covered and we must take care that the measure of these pieces should be sufficiently small. For this we use a method proposed in [3].

Let $B_{t}$ be a closed ball in $W^{\mathbf{1}}$ of radius $t$ and $A=\hat{\Lambda}_{k, s, r}^{l} \cap B_{\boldsymbol{t}}$. We select a number $\beta>0$ such that

$$
\begin{equation*}
\left|\frac{\mu^{1}\left(B_{t+\beta}\right)}{\mu^{1}\left(B_{i}\right)}-1\right| \leqslant \frac{1}{2}\left(\left|W^{1}\right|+\left|W^{2}\right|\right) . \tag{3.2.37}
\end{equation*}
$$

Lemma 3.2.7. There exist $\delta_{0}, 0<\delta_{0} \leqslant 1 / 2\left(\left|W^{1}\right|+\left|W^{2}\right|\right)$, and for each $m \in \mathbf{Z}^{+}$sets of points $\left\{w_{j}\right\}$ and $\left\{y_{j}^{1}\right\}, j=1, \ldots, n_{1}$, satisfying the following conditions:

1. $w_{j} \in \Lambda_{k, s, r}^{l} \cap U\left(x, q_{k, s, r}^{l}\right), y_{j}^{1} \in U\left(x, q_{k, s, r}^{l}\right) \cap V\left(w_{j}\right), j=1, \ldots, n_{1}$.
2. $f^{m}(A) \subset \widetilde{W}_{m}^{1}=\bigcup_{j=1}^{n_{1}} W_{m}^{1}\left(w_{j}, y_{j}^{1}, \frac{1}{2} \delta_{0}\right) \subset \hat{W}_{m}^{1}=\bigcup_{j=1}^{n_{1}} \hat{W}_{m}^{1}\left(w_{j}, y_{j}^{1}, \quad \delta_{0}\right)$ $\subset f^{m}\left(B_{i+\beta}\right)$.

Proof. Let $w \in \Lambda_{k, s, r}^{l} \cap U\left(x, q_{k, s, r}^{l}\right)$ and $y^{1} \in W^{1} \cap V(w)$, and let $Q\left(y^{1}, \gamma\right)$ be the ball in the submanifold $W^{1}$ with center at $y^{1}$ and radius $\gamma$. There exists $t=t(\gamma) \leqslant \gamma$, not depending on $w$ and $y$, such that

$$
\begin{equation*}
W_{0}^{1}\left(w, y^{1}, t(\gamma)\right) \subset Q\left(y^{1}, \gamma\right) \tag{3.2.38}
\end{equation*}
$$

Put

$$
\begin{equation*}
\delta_{0}=\min \left\{t\left(\frac{\beta}{4}\right), \frac{1}{2}\left(\left|W^{1}\right|+\left|W^{2}\right|\right)\right\} \tag{3.2.39}
\end{equation*}
$$

Fix $m>0$ and choose from the covering set

$$
\bigcup_{w \in \Lambda_{k, s, r}^{l}} \bar{W}^{1}\left(w, y^{1}, \frac{\delta_{0}}{2}\right)
$$

a finite subcover. The corresponding sets we denote by $W_{m}^{1}\left(w_{j}, y_{j}^{1}, 1 / 2 \delta_{0}\right), j=1, \ldots, n_{1}$. We must show that $\hat{W}_{m}^{1} \subset f^{m}\left(B_{t+\beta}\right)$. If this is not so, then there exists a point

$$
\rho \in \partial\left(f^{m}\left(B_{t+\beta}\right)\right) \cap W_{m}^{1}\left(w_{j}, y_{i}^{\mathbf{1}}, \delta_{0}\right)
$$

for some $j$. By (3.2.38), (3.2.39) and (3.2.26), $f^{-m}(p) \in Q\left(y_{j}^{1}, 1 / 4 \beta\right)$ and, by the same token, $f^{-m}(p) \notin \partial\left(B_{t+\beta}\right)$. This contradiction proves the lemma.

The following lemma on covers with bounded (independent of $m$ ) multiplicity is proved by a simple modification of the argument of [3] ( $\S 2$, Lemma 7).

Lemma 3.2.8. There exist $d_{0}>0$, an integer $L>0$ and for every $m \in \mathbf{Z}^{+} a$ set of balls $Q\left(z_{j}, d_{m}\right), j=1, \ldots, n_{2}$, in the submanifold $f^{m}\left(W^{1}\right)$ with centers at points $z_{j}$ and radii $d_{m}=d_{0} \xi^{m}$ satisfying the following conditions:

1) $\widetilde{W}_{m}^{1} \subset \bigcup_{j=1}^{n_{2}} Q\left(z_{j}, d_{m}\right) \subset \bigcup_{j=1}^{n_{2}} Q\left(z_{j}, 2 d_{m}\right) \subset \hat{W}_{m}^{1}$.
2) For each $j=1, \ldots, n_{2}$ there is an index $i, 1 \leqslant i \leqslant n_{1}$, such that

$$
Q_{\Delta}\left(z^{i}, 2 d_{m}\right) \subset W_{m}^{1}\left(w_{i}, y_{i}^{1}, \delta_{0}\right)
$$

3) The multiplicity of the cover of $\widetilde{W}_{m}^{1}$ by the sets $Q\left(z_{j}, 2 d_{m}\right)$ is equal to $L$.

Fix $m>0$ and consider any ball $Q\left(z_{j}, 2 d_{m}\right)$. For some $i, 1 \leqslant i \leqslant n_{1}$,

$$
Q\left(z_{j}, 2 d_{m}\right) \subset W_{m}^{1}\left(w_{i}, y_{i}^{1}, \delta_{0}\right)
$$

We denote

$$
\hat{Q}_{i}=\exp _{w_{i}}^{-1} Q\left(z_{j}, 2 d_{m}\right), \quad \hat{z}_{i}=\exp _{w_{i}}^{-1} z_{j}, \quad \tilde{z}_{j}=t_{w_{i}}^{m}\left(\hat{z_{j}}\right),
$$

where $t_{w_{i}}^{m}$ is the projection onto $E_{2 f m\left(w_{i}\right)}$ parallel to $E_{1 f m\left(w_{i}\right)}$.
Lemma 3.2.9. For each $\epsilon, 0<\epsilon \leqslant 1$, there exists $m_{2}=m_{2}(\epsilon, l, s, r) \in \mathbf{Z}^{+}$such that for any $m \geqslant m_{2}$

$$
\left.\rho\left(\partial\left(t_{w_{i}}^{m}\left(\hat{Q}_{j}\right)\right), \partial B \widetilde{z}_{j}, 2 d_{m}\right)\right) \leqslant \varepsilon d_{m}
$$

The proof follows from (3.2.1) of Lemma 3.2.2, (3.2.19) and (3.2.25).
Fix an arbitrary $\epsilon, 0<\epsilon \leqslant 1$, and consider a cover of the ball $B\left(\widetilde{z_{j}},(2-\epsilon) d_{m}\right)$ by closed cubes $\widetilde{D}_{i j}, i=1, \ldots, N_{j}$, of diameter $\epsilon d_{m}$ with disjoint interiors. The cubes are constructed in the obvious way relative to the Euclidean structure on $E_{2 f} m_{\left(w_{j}\right)}$ generated by the Riemannian metric. We denote by $V\left(\widetilde{D}_{i j}\right)$ the Riemannian volume of the cube $\widetilde{D}_{i j}$. It is easy to see that for $\alpha_{m}=\alpha \xi^{m}, \alpha>0$,

$$
\begin{equation*}
\left|\frac{V\left(\left(\widetilde{D}_{i j}\right)_{\alpha_{m}}\right)}{V\left(\widetilde{D}_{i j}\right)}-1\right| \leqslant C_{12} \frac{\alpha}{\varepsilon d_{0}}, \tag{3.2.40}
\end{equation*}
$$

where $C_{12}$ is a constant and $\left(\widetilde{D}_{i j}\right)_{\alpha}$ denotes the $\alpha$-blowing up of $\widetilde{D}_{i j}$. Put (see (3.2.25))

$$
\begin{equation*}
D_{i j}=\exp _{w_{i}}\left\{\left(\psi_{m}^{1}(v), v\right): v \in \widetilde{D}_{i j}\right\} \tag{3.2.41}
\end{equation*}
$$

(We recall that the graph of $\psi_{m}^{1}$ coincides with $\exp _{w_{i}}^{-1} W_{m}^{1}$.) It is easy to see that the sets $D_{i j}$ intersect only on the boundary, and if $m \geqslant m_{2}$ (see Lemma 3.2.9), then

$$
\begin{equation*}
Q\left(z_{j}, 2 d_{m}\right) \supset \bigcup_{i=1}^{N_{j}} D_{i j} \supset Q\left(z_{j}, d_{m}\right) \tag{3.2.42}
\end{equation*}
$$

Making similar constructions for each ball $Q\left(z_{i}, 2 d_{m}\right), j=1, \ldots, n_{2}$, we construct sets $D_{i j}$ which, by Lemma 3.2 .8 (see assertion 1) and (3.2.42), cover the set $\widetilde{W}_{m}^{1}$. The multiplicity of the cover, as before, does not exceed $L$. We will show how it is possible, neglecting part of the elements of this cover, to obtain a cover which at most (in measure) points has multiplicity 1 . For this, passing successively from $Q\left(z_{j}, 2 d_{m}\right)$ to $Q\left(z_{j+1}, 2 d_{m}\right)$, we neglect those sets $D_{i j+1}$ whose centers are contained in the union of the sets $D_{i k}, k \leqslant$ $j$, retained in the preceeding steps. Let $D_{i}, i=1, \ldots, N$, be the renumbered elements of the cover obtained. It is not difficult to see that

$$
\begin{equation*}
\bigcup_{i=1} D_{i}=\bigcup_{j=1}^{n_{2}} \bigcup_{i=1}^{N_{j}} D_{i j} \tag{3.2.43}
\end{equation*}
$$

Lemma 3.2.10. There exists a $C_{13}=C_{13}(l, s, r)$ such that for every $\epsilon>0$ and any $m \geqslant m_{2}$ (see Lemma 3.2.9)

1. $\widetilde{W}_{m}^{1} \subset \bigcup_{i=1}^{N} D_{i} \subset \widehat{W}_{m i}^{1}$,
2. 

$$
\begin{equation*}
\left|\frac{\sum_{i=1}^{N} \mu_{0}^{1}\left(f^{-m}\left(D_{i}\right)\right)}{\mu_{0}^{1}\left(f^{-m}\left(\bigcup_{i=1}^{N} D_{i}\right)\right)}-1\right| \leqslant C_{13} \varepsilon . \tag{3.2.44}
\end{equation*}
$$

Proof. The first assertion follows immediately from (3.2.43) and Lemma 3.2.8 (see assertion 1). In order to prove the second assertion we consider the cover of $\widetilde{W}_{m}^{1}$ by the sets $Q\left(z_{j},(2-\epsilon) d_{m}\right), j=1, \ldots, n_{2}$. Since the multiplicity of the cover does not exceed $L$, we have

$$
\begin{gather*}
\sum_{j=1}^{n_{2}} \mu_{m}^{1}\left(Q\left(z_{j},(2-\varepsilon) d_{m}\right)\right)  \tag{3.2.45}\\
\leqslant L \mu_{m}^{1}\left(\bigcup_{j=1}^{n_{2}} Q\left(z_{j},(2-\varepsilon) d_{m}\right)\right) \leqslant L \mu_{m}^{1}\left(\bigcup_{i=1}^{N} D_{i}\right) .
\end{gather*}
$$

We denote by $n(x)$ the number of sets $D_{i}$ containing $x$, and put

$$
D_{\varepsilon_{j}}=\left\{x \in f^{m}\left(W^{1}\right): \bar{\rho}\left(x, \partial Q\left(z_{j},(2-\varepsilon) d_{m}\right)\right) \leqslant \varepsilon d_{m}\right\}
$$

Here $\bar{\rho}$ is the distance on the submanifold $f^{m}(W)$ induced by the Riemannian metric. It is obvious that

$$
\begin{equation*}
\mu_{m}^{1}\left(D_{\varepsilon_{j}}\right) \leqslant C_{14} \varepsilon \mu\left(Q\left(z_{j},(2-\varepsilon) d_{m}\right)\right) \tag{3.2.46}
\end{equation*}
$$

where $C_{14}$ is a constant. Therefore if $D_{\epsilon}=\bigcup_{1}^{n 2} D_{\epsilon j}$, then

$$
\begin{equation*}
\mu_{m}^{1}\left(D_{\varepsilon}\right) \leqslant \sum_{j=1}^{n_{2}} \mu_{m}^{1}\left(D_{\varepsilon j}\right) \leqslant C_{14} \varepsilon \sum_{j=1}^{n_{2}} \mu_{m}^{1}\left(Q\left(z_{j},(2-\varepsilon) d_{m}\right)\right) \tag{3.2.47}
\end{equation*}
$$

It is easy to see also that

$$
\begin{equation*}
n(x)=1, \quad x \notin D_{\varepsilon} \bigcup\left(\bigcup_{i=1}^{N} \partial D_{i}\right) . \tag{3.2.48}
\end{equation*}
$$

Therefore (3.2.45)-(3.2.48) imply

$$
\begin{gathered}
\sum_{i=1}^{N} \mu\left(D_{i}\right)=\int_{\sum_{i=1}^{N} D_{i}} n(x) d \mu_{m}^{1}(x) \leqslant \mu_{m}^{1}\left(\bigcup_{i=1}^{N} D_{i}\right) \\
+(L-1) \mu_{m}^{1}\left(D_{\varepsilon}\right) \leqslant \mu_{m}^{1}\left(\bigcup_{i=1}^{N} D_{i}\right)\left(1+(L-1) L C_{14} \varepsilon\right) .
\end{gathered}
$$

Now (3.2.44) follows from Lemma 3.2.7, in which we must put $X^{2}=\bigcup_{1}^{N} D_{i}$ and $X^{1}=$ $X^{2} \backslash D_{\epsilon}$. The lemma is proved.

We complete the proof of the theorem. Choose $\alpha>0$ and consider the set

$$
\hat{D}_{i}=\exp _{w_{i}}^{-1}\left\{\left(\psi_{m}^{2}(v), v\right): v \in\left(\widetilde{D}_{i}\right)_{a}\right\} \subset W_{m}^{2}\left(w_{j}, y_{j}^{2}, \delta_{0}\right) .
$$

Since

$$
\bigcup_{j=1}^{N} W_{m}^{2}\left(w_{i}, y_{j}^{2}, \delta_{0}+\alpha\right) \subset \bigcup_{i=1}^{N} \hat{D}_{l}
$$

on the basis of Lemma 3.2.4 we have for $m \geqslant m_{1}$

$$
\begin{equation*}
\bigcup_{i=1}^{N} \hat{D}_{i} \supseteq f^{m}(p(A)) \tag{3.2.49}
\end{equation*}
$$

In addition, from the continuity of the Riemannian metric on $M$, the inequalities (3.2.19) and (3.2.40) and the definition of $\mu$ (see [10]), it follows that there exists a constant $C_{16}=C_{16}(l, s, r)$ such that

$$
\begin{equation*}
\left|\frac{\mu_{m}^{1}\left(D_{i}\right)}{\mu_{m}^{2}\left(\hat{D}_{i}\right)}-1\right| \leqslant C_{16}\left(\frac{\alpha}{\varepsilon d_{0}}+\left|W^{1}\right|+\left|W^{2}\right|\right) \tag{3.2.50}
\end{equation*}
$$

It follows from (3.2.49), (3.2.50), (3.2.36), (3.2.44), (3.2.37) and Lemmas 3.2.6, 3.2.7, 3.2.9 and 3.2.10 that for $m \geqslant \max \left\{m_{1}(\alpha), m_{2}(\epsilon)\right\}$

$$
\begin{gathered}
\mu_{0}^{2}(p(A)) \leqslant \sum_{i=1}^{N} \mu_{0}^{2}\left(f^{-m}\left(\hat{D}_{i}\right)\right) \\
\leqslant\left(1+C_{11}\left(\left|W^{1}\right|+\left|W^{2}\right|\right)+C_{16}\left(\frac{\alpha}{\varepsilon d_{0}}+\left|W^{1}\right|+\left|W^{2}\right|\right)\right) \\
\times \sum^{N} \mu_{0}^{1}\left(f^{-m}\left(D_{i}\right)\right) \leqslant\left(1+C_{17}\left(\left|W^{1}\right|+\left|W^{2}\right|+\frac{\alpha}{\varepsilon d_{0}}+\varepsilon\right)\right) \\
\times \mu_{0}\left(B_{t+\beta}\right) \leqslant\left(1+C_{18}\left(\left|W^{1}\right|+\left|W^{2}\right|+\frac{\alpha}{\varepsilon d_{0}}+\varepsilon\right)\right) \mu_{0}^{1}\left(B_{t}\right) .
\end{gathered}
$$

Here $C_{17}=C_{17}(l, s, r)$ and $C_{18}=C_{18}(l, s, r)$ are constants. Since the ball $B_{t}$ is chosen arbitrarily, the mapping $p$ is absolutely continuous. Now let $z$ be a point of density of $B=W^{1} \cap \Lambda_{k, s, r}^{I}$, and $B_{t}(z)$ the ball in $W^{1}$ with center at $z$ and radius $t$. For sufficiently small $t$

$$
\left|\frac{\mu^{1}\left(B \cap B_{t}(z)\right)}{\mu^{1}\left(B_{t}(z)\right)}-1\right| \leqslant \min \left\{\frac{1}{2}, \frac{1}{2}\left(\left|W^{1}\right|+\left|W^{2}\right|\right)\right\}
$$

Therefore from the above it follows that

$$
\mu_{0}^{2}\left(p\left(B \cap B_{t}(z)\right)\right) \leqslant\left(1+C_{19}\left(\left|W^{\mathbf{1}}\right|+\left|W^{2}\right|+\frac{\alpha}{\varepsilon d_{j}}+\varepsilon\right)\right) \mu_{0}^{1}\left(B \cap B_{t}(z)\right)
$$

where $C_{19}=C_{19}(l, s, r)$ is a constant. Since the numbers $\alpha, \epsilon$ and $d_{0}$ are chosen arbitrarily, where $\alpha$ can be chosen after $d_{0}$ and $\epsilon$, it follows that

$$
\begin{equation*}
\mu_{0}^{2}\left(p\left(B \cap B_{t}(z)\right)\right) \leqslant\left(1+C_{19}\left(\left|W^{1}\right|+\left|W^{2}\right|\right)\right) \mu_{0}^{1}\left(B \cap B_{t}(z)\right) . \tag{3.2.51}
\end{equation*}
$$

Applying similar assertions to the mapping $p^{-1}$, given on the set $p\left(B \cap B_{t}(z)\right)$, we obtain that

$$
\begin{equation*}
\mu_{0}^{1}\left(B \cap B_{t}(z)\right) \leqslant\left(1+C_{19}\left(\left|W^{\mathbf{1}}\right|+\left|W^{2}\right|\right)\right) \mu_{0}^{2}\left(p\left(B \cap B_{t}(z)\right)\right) . \tag{3.2.52}
\end{equation*}
$$

The inequality (3.2.1) follows from (3.2.51) and (3.2.52). The theorem is proved.
3.3. Consider a set $\Lambda_{k, s, r}^{l}$ of positive measure, and let $x$ be a point of density of $\Lambda_{k, s, r}^{l}$. There exists a partition $\xi^{s m}$ of the neighborhood $U\left(x, q_{k, s, r}^{l}\right)$ (the number $q_{k, s, r}^{l}$ was constructed in Theorem 3.2.1), each element of which is an ( $n-k$ )-dimensional submanifold transversal to the submanifold $V(x)$. We denote the elements of this partition by $C_{\xi \operatorname{sm}}(y)$, $y \in V(x)$. Let $A \subset C_{\xi \mathrm{sm}}(x) \cap \hat{\Lambda}_{k, s, r}^{l}$ be a measurable set. Put

$$
\begin{equation*}
\hat{A}=\bigcup_{z \in \Lambda_{k, s, r}^{l}, V(z) \cap A \neq \varnothing}\left(V(z) \cap U\left(x, q_{k, s, r}^{l}\right)\right) \tag{3.3.1}
\end{equation*}
$$

Using the absolute continuity of the families of local stable manifolds, we will show that the partition $\zeta$ of $\hat{A}$ into the submanifolds $V(z), z \in \Lambda_{k, s, r}^{l} V(z) \cap A \neq \varnothing$, has a system of conditional measures absolutely continuous relative to the measure on $V(z)$ induced by the Riemannian metric. The corresponding argument is due to Ja. G. Sinail (see [2], Russian p. 153, English pp. 147-148).

Let $\mu_{y}$ be the measure on $C_{\xi \mathrm{sm}}(y)$ induced by the restriction of the Riemannian metric on $M$ to $C_{\xi \mathrm{sm}}(y)$. In it we introduce another measure $\hat{\mu}_{y}$, putting for any measurable set $A$

$$
\hat{\mu}_{y}(A)=v(\hat{A}) .
$$

Let $\nu_{z}$ denote the measure on the submanifold $V(z)$ induced by the Riemannian metric.
Proposition 3.3.1. 1. $\hat{\mu}_{y}$ is absolutely continuous relative to $\mu_{y}$.
2. The partition $\zeta$ is measurable, and the conditional measure $\hat{\nu}_{z}$ on an element of the partition is absolutely continuous relative to the measure $\nu_{z}, z \in \Lambda_{k, s, r}^{l} V(z) \cap A \neq \varnothing$.
3. For almost all $y \in A, \nu_{2}(V(z))>0$, where $z \in \Lambda_{k, s, r}^{l}$ and $V(z) \cap A=y$.

Proof. It follows from Theorem 2.3.1 that $\zeta$ is measurable. Denote by $\tilde{\mu}_{y}$ the conditional measure induced by $\nu$ on the element $C_{\xi s m}(y)$ and by $\kappa$ the measure in the quotient space $\hat{A} /_{\xi s m}=X$. It is obvious that there exist a measurable function $t_{y}=t_{y}(s)$, $s \in C_{\xi s m}(y), \phi=\phi(y), y \in V(z), z \in \Lambda_{k, s, r}^{l}$, such that

$$
d \tilde{\mu}_{y}(s)=t_{y}(s) d \mu_{y}(s), \quad d x(y)=\varphi(y) d v_{z}(y)
$$

Let $B \subset \hat{A}$ be an arbitrary set of positive measure. Then from the above and Fubini's theorem it follows that

$$
\begin{aligned}
v(B)= & \int_{X} d x(y) \int_{c_{\xi} s \mathrm{sm}(g)} \chi_{B}(s) d \tilde{\mu}_{y}(s)=\int_{X} d x(y) \int_{C_{\xi} s \mathrm{sm}(y)} \chi_{B}(s) \dot{t}_{y}(s) d \mu_{y}(s) \\
= & \int_{X} d x(y) \int_{C_{\xi} \mathrm{sm}\left(q_{0}\right)} J_{q_{0}}(q) \chi_{B}\left(p_{q_{0}}(q)\right) t\left(p_{q_{0}}(q)\right) d \mu_{q_{0}}(q) \\
& =\int_{C_{\xi s \mathrm{~m}}\left(q_{0}\right)} d \mu_{q_{0}}(q) \int_{X} J_{u}(q) \chi_{B}\left(p_{u}(q)\right) t\left(p_{u}(q)\right) d x(u) \\
= & \int_{C_{\xi} \mathrm{sm}^{\left(q_{0}\right)}} d \mu_{q_{0}}(q) \int_{V(q)} J_{u}(q) \chi_{B}\left(p_{u}(q)\right) t\left(p_{u}(q)\right) \varphi(u) d v_{q}(u) .
\end{aligned}
$$

Here $J_{q_{0}}(q)$ is the positive measurable function arising from the absolute continuity of the successor mapping $p_{q 0}: C_{\xi \mathrm{sm}}\left(q_{0}\right) \rightarrow C_{\xi \mathrm{sm}}(y), p_{q 0}(q)=s$. (We note that $J_{q 0}(q)=0$ for $q \notin C_{\xi \mathrm{sm}}\left(q_{0}\right) \cap \hat{A}$.) For fixed $q$ and variable $q_{0}$ the points $p_{q 0}(q)$ run through $V(q)$. From the latter equality it follows that on the elements $C_{\zeta}(q)$ the conditional measure is absolutely continuous relative to the measure $\nu_{q}$, and the measure on $\hat{A} / \zeta$ is absolutely continuous relative to $\mu_{y}$. Assertions 1 and 2 are proved.

In order to prove assertion 3 we denote

$$
Y=\left\{z \in \hat{A} \cap \Lambda_{k, s, r}^{l}: v_{z}\left(C_{\zeta}(z)\right)=0\right\}
$$

and let $Z=\bigcup_{y \in Y} C_{\zeta}(y)$. Obviously $Y \subset Z$. From 1 and 2 it follows that

$$
v(Y) \leqslant v(Z)=\int_{z / \zeta} \tilde{v}_{z}\left(C_{\zeta}(z)\right) d \hat{\mu}_{x}(z)=\int_{z / \zeta} l(z) v_{z}\left(C_{\zeta}(z)\right) d \hat{\mu}_{x}(z)=0
$$

Here $l(z)$ is the measurable function arising from the absolute continuity of $\hat{\nu}_{z}$ relative to $\nu_{z}$. The proposition is proved.
3.4. Proposition 3.4.1. Let $\nu\left(\Lambda_{k, s, r}\right)>0$. There exists a set $N \subset M$ of measure zero such that for any $l \in \mathbf{Z}^{+}$and $x \in \Lambda_{k, s, r}^{l} \backslash N$

$$
\boldsymbol{v}_{x}\left(V(x) \cap \Lambda_{k, s, r}\right)=v_{x}(V(x))
$$

(recall that $\nu_{x}$ is the Riemannian volume on $V(x)$ induced by the Riemannian metric on $M$ ).
Proof. Let $N_{1}$ be the set of nonregular points in $M$. We have $\nu\left(N_{1}\right)=0$ (see Theorem 0.3). Let $x \in \Lambda_{k, s,{ }^{\prime}}$ Put $A_{x}=\left\{y \in V(x): y \in N_{1}\right\}$ and consider the set $N_{2}$ of those points $x \in \Lambda_{k, s, r}$ for which $\nu_{x}\left(A_{x}\right)>0$. The set $N_{2}$ is measurable. We will show that $\nu\left(N_{2}\right)=0$. If this is not so, then for some $l \in \mathbf{Z}^{+}$the set $N_{2} \cap \Lambda_{k, s, r}^{l}$ has positive measure. Let $x$ be a point of density of this set. By Proposition 3.3.1, $v_{y}\left(V(y) \cap N_{1}\right)=0$ for almost all $y \in \Lambda_{k, s, r}^{l} \cap U\left(x, q_{k, s, r}^{l}\right)$. Thus the set of those points $y \in \Lambda_{k, s, r}^{l} \cap$ $U\left(x, q_{k, s, r}^{l}\right)$ for which $\nu_{y}\left(V(y) \cap N_{1}\right)>0$ has zero measure. This contradiction shows that $\nu\left(N_{2}\right)=0$. If $x \in \Lambda_{k, s,}^{l} \quad \cup N_{2}$ ), then almost any (relative to $\nu_{x}$ ) point $y \in V(x)$ is a regular point, and, by $\operatorname{Pr}$ anon 2.3.1, $y \in \Lambda_{k, s, r}$. The proposition is proved.

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