

Invariant measures for some one-dimensional attractors

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Dedicated to the memory of V. M. Alexeyev to whom I shall always be grateful for his friendship and encouragement

Abstract. We consider certain non-invertible maps of the square which are extensions of the quadratic maps of the interval and their small perturbations. We show that several maps of the type possess attractors which are not hyperbolic but have invariant measures similar to Bowen–Ruelle measures for hyperbolic attractors.

1. The properties of a dynamical system $f: X \rightarrow X$ which possesses an invariant contracting foliation are closely related to the properties of the system induced on the quotient space. When (X, f) is studied from the point of view of the topology, this relation is based on the inverse limit construction (see [14]), but when f -invariant measures are studied one uses the natural extension construction (see [12]).

In this paper we consider the situation when the quotient space is an interval and the corresponding induced one-dimensional map has a single critical point. We shall see that several systems of this type possess attractors which are not hyperbolic but have similar properties.

The maps under consideration are non-invertible. We show, however, that the set where the inverse map is not defined is negligible from the point of view of any absolutely continuous measure. We shall say that an attractor Λ of the map f admits an absolutely continuous invariant Bernoulli measure if there exists a finite measure μ such that (Λ, f, μ) is a Bernoulli automorphism and there is an invariant subset $\Lambda_1 \subset \Lambda$ satisfying $\mu(\Lambda_1) = 1$, where Λ_1 is a union of rectifiable curves γ_α , forming (mod 0) a measurable partition of Λ_1 , and on every curve γ_α μ induces the conditional measure $\mu(\cdot | \gamma_\alpha)$ which is absolutely continuous with respect to normalized length.

In §§ 2–5 we consider a map $F: S \rightarrow S$ of the unit square S similar to the ‘twisted horseshoe’ map from [7] and prove the following theorem.

THEOREM 1. *The attractor of the map F admits an absolutely continuous Bernoulli measure.*

In §§ 6–10 we consider some one-parameter families of maps $F_a: S \rightarrow S$ which induce on the quotient space mappings $x \rightarrow g_a(x)$ close to $x \rightarrow ax(1-x)$. For these families we prove the following theorem using the results of [8].

THEOREM 2. *There exists a set M of parameter values of positive measure, so that for any $a \in M$ the attractor of the map F_a admits an absolutely continuous Bernoulli measure.*

The approach to the proof of these results is based on the ‘induced map’ method combined with the ‘itinerary schemes’ method due to Alexeyev [1].

In the last year of his life, despite his grave illness, V. M. Alexeyev continued his work on various mathematical problems. Several times we discussed the subject of this paper. His attentive and friendly attitude stimulated my work to a great extent.

I would like to thank B. M. Gurevich and I. P. Gornfeld for discussions and useful remarks.

2. We denote by $\mathcal{A} \subset C^2([0, 1], [0, 2])$ the set of mappings satisfying $g(0) = g(1) = 0$. Let $\varepsilon > 0$ be a small constant. We define $\mathcal{A}_\varepsilon \subset \mathcal{A}$ as the set of mappings g satisfying

$$(i) \quad \|g(x) - 4x(1-x)\|_{C^2} < \varepsilon;$$

$$(ii) \quad \max_{x \in [0, 1]} g(x) = 1.$$

\mathcal{A}_ε is a surface of codimension 1 in \mathcal{A} . The curve $\Gamma = \{x \rightarrow ax(1-x)\}$ intersects \mathcal{A}_ε transversally at $a = 4$, thus for any one-parameter family $x \rightarrow g_a(x) \in \mathcal{A}$ sufficiently close in $C^2(A \times [0, 1], [0, 2])$ to $x \rightarrow ax(1-x)$ where $a \in A = [0, 5]$, $x \in [0, 1]$, there exists an a_0 such that $g_{a_0} \in \mathcal{A}_\varepsilon$.

Let us consider the square $S = \{0 \leq x \leq 1, 0 \leq y \leq 1\}$ and the map $F: S \rightarrow S$ defined by

$$F: (x, y) \rightarrow (\varphi(y) + \lambda(x - \frac{1}{2}), g(y)) \tag{1}$$

where $\lambda > 0$ is a small constant, $g(y) \in \mathcal{A}_\varepsilon$ and $\varphi(y)$ is a C^2 function satisfying the following conditions:

$$\begin{aligned} &\varphi'(y) \leq 1; \\ \varphi(y) = &\begin{cases} \frac{1}{4} + 2\delta & \text{for } y \in [0, \frac{1}{4} + \delta] \\ y & \text{for } y \in [r_1 - \delta, 1 - r_1 + \delta] \\ \frac{3}{4} - 2\delta & \text{for } y \in [\frac{3}{4} - \delta, 1] \end{cases} \end{aligned}$$

where r_1 is defined in (3) and $\delta > 0$ is the small constant defined in (4), see § 3.

The map F is illustrated by figure 1.

Remark. We have chosen the above formula for φ in order to simplify the calculations. One can check that the following results are valid for a large class of φ .

A similar ‘twisted horseshoe’ map was analysed in [7]. Bowen in [4] considered such a map as an example whose dynamics are undecided.

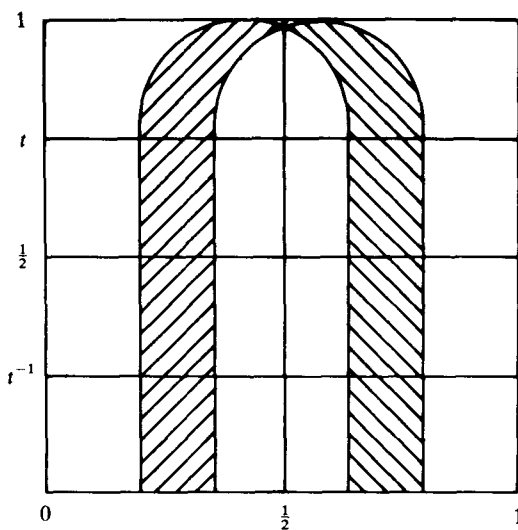


FIGURE 1

(1) implies that the horizontal foliation $\xi = \{\xi_y\}$ of S is invariant under F , and $F: \xi_y \rightarrow \xi_{g(y)}$ is a contraction with coefficient λ .

Let us consider a sequence $(y_0, y_1, \dots, y_n, \dots)$, where $y_n \in g^{-1}(y_{n-1})$, and $g^{-1}(y)$ is the inverse image of y under g . We have $F^n \xi_{y_n} \subset \xi_y$, $\text{diam } F^n \xi_{y_n} < \lambda^n$ and thus $\bigcap_{n=0}^{\infty} F^n \xi_{y_n}$ is a point which we denote by $M(y_0, y_1, \dots, y_n, \dots)$.

Following the inverse limit construction (see [14]), we introduce a space $\tilde{Y} = \{\tilde{y} = (y_0, y_1, \dots)\}$ and we define the map $\tilde{g}: \tilde{Y} \rightarrow \tilde{Y}$ by $\tilde{g}(\tilde{y}) = (g(y_0), y_0, y_1, \dots)$. We shall use the notation $\Lambda = \bigcap_{n=0}^{\infty} F^n S$, and $\pi: (y_0, y_1, \dots) \rightarrow M(y_0, y_1, \dots)$. It follows from the definition of the topology in \tilde{Y} that π is continuous, and the construction implies that $\pi: \tilde{Y} \rightarrow \Lambda$ is onto and satisfies

$$\pi \circ \tilde{g} = F \circ \pi. \tag{2}$$

3. The map π is not one-to-one because \tilde{g} is a homeomorphism and F is non-invertible on Λ . We shall see nevertheless that π induces an isomorphism between the dynamical system (Λ, F, μ_*) and the natural extension $(\tilde{Y}, \tilde{g}, \tilde{\mu})$ of the system (Y, g, μ) , where μ is a g -invariant measure absolutely continuous with respect to the Lebesgue measure dy on $[0, 1]$, and $\mu_* = \pi_* \tilde{\mu}$.

It is known that any $g \in \mathcal{A}_e$ admits an invariant measure μ_g absolutely continuous with respect to dy . Let us recall the construction of μ_g from [9].

The mapping g has a repelling fixed point $t = g(t) \neq 0$. Let t^{-1} be the second pre-image of t . Let us denote $I = [t^{-1}, t]$. Let $G: I \rightarrow I$ be the map induced by g , i.e., for $y \in I$ one sets $G(y) = g^{n(y)}(y)$, where

$$n(y) = \min\{n \geq 1: g^n(y) \in I\}.$$

We have $I = (\bigcup_{i=2}^{\infty} \Delta_{i\pm}) \cup K$, where $\Delta_{i\pm}$ are open intervals such that for

$$\begin{aligned} y \in \Delta_{i\pm}(y) = i, \quad \Delta_{i-} \subset [0, c], \quad \Delta_{i+} \subset [c, 1], \\ g \Delta_{i+} = g \Delta_{i-}, \quad \Delta_i \cap \Delta_j = \emptyset, \end{aligned}$$

and K is the union of the set of end points of $\Delta_{i\pm}$ and of the critical point c , see figure 2. Let us denote $g_0(y) = 4y(1 - y)$. Then

$$\mu_0 = \mu_{g_0} = dy/\pi\sqrt{y(1-y)}$$

is the corresponding measure. For $g = g_0$ we have $I = I_0 = [t_0, t_0^{-1}]$, where

$$t_0 = \frac{3}{4}, \quad t_0^{-1} = \frac{1}{4}, \quad \Delta_{2-} = [\frac{1}{4}, r_1], \quad \Delta_{2+} = [1 - r_1, \frac{3}{4}]. \tag{3}$$

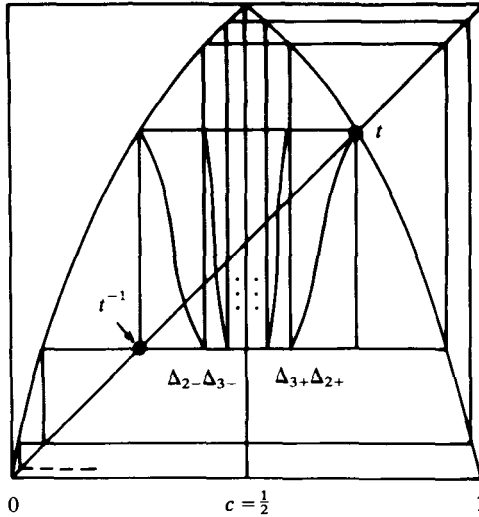


FIGURE 2

The map G_0 induced by g_0 satisfies the following conditions:

- (A) $\exists k_0: |DG_0^{k_0}(y)| > c_0 > 1$ for any y such that $G_0^{k_0}(y)$ is defined.
- (B) $\sup_{y \in \Delta_i} |D^2G_0(y)|/|DG_0(y)| \cdot |\Delta_i| < c'_0$, with c'_0 independent from i .

If ε is sufficiently small any $g \in \mathcal{A}_\varepsilon$ induces G which also satisfies (A) and (B) on the interval $I = I_g = [t^{-1}, t]$. Let us denote

$$\Delta_{2-}(g) = [t^{-1}, \rho_1], \quad \Delta_{2+}(g) = [\rho_2, t],$$

and let δ be defined by

$$\sup_{g \in \mathcal{A}_\varepsilon} \max \{ |t^{-1} - \frac{1}{4}|, |\rho_1 - r_1|, |\rho_2 - (1 - r_1)|, |t - \frac{3}{4}| \} = \delta/2. \tag{4}$$

Here δ is the constant from the definition of $\varphi(y)$ in § 2.

Conditions (A) and (B) imply ([2], [13]) that G admits an invariant measure

$$\nu_g(dy) = \chi(y) dy$$

on I with $\chi(y)$ continuous and bounded away from 0.

The f -invariant measure μ_g is constructed from ν_g according to the following formula, which holds not only for the induced map G but also in a more general

situation when G locally coincides with different powers of the map g :

$$\mu(E) = \sum_{i=1}^{\infty} \sum_{k=0}^{i-1} \nu(g^{-n}(E \cap g^n \Delta_i)). \tag{5}$$

When G is an induced map, we have $\mu \equiv \nu|I$ and $\mu(dy) = q(y) dy$ where the density $q(y)$ is continuous on $[0, 1]$, bounded away from zero and has two singularities of the type $1/\sqrt{y}$ at 0 and at 1.

Let $(\tilde{I}, \tilde{G}, \tilde{\nu})$ be the natural extension of (I, G, ν) . The construction of § 8 below implies that the automorphisms of measure spaces $(\tilde{Y}, \tilde{g}, \tilde{\mu})$ and $(\tilde{I}, \tilde{G}, \tilde{\nu})$ are related in the following way. $(\tilde{I}, \tilde{G}, \tilde{\nu})$ is isomorphic to an automorphism induced by \tilde{g} on a subset \tilde{I}' of \tilde{Y} . Let $\tilde{q}: \tilde{I}' \rightarrow \tilde{I}'$ be this isomorphism. For

$$\tilde{z} \in \tilde{I}, \quad \tilde{z} = (z_0, z_1, z_2, \dots), \quad Gz_i = z_{i-1},$$

for any $i = 1, 2, \dots$ we have $Gz_i = g^{n_i}z_i$. Then $\tilde{q}(\tilde{z}) = \tilde{y}$ is defined by

$$\tilde{y} = (z_0, g^{n_1-1}z_1, \dots, gz_1, z_1, g^{n_2-1}z_2, \dots, gz_2, z_2, \dots).$$

Furthermore \tilde{I}' coincides mod 0 (i.e., neglecting sets of zero $\tilde{\mu}$ -measure) with the set $\{\tilde{y} = (y_0, y_1, \dots): y_0 \in I\}$.

$(\tilde{I}, \tilde{G}, \tilde{\nu})$ admits a pair of invariant continuous partitions analogous to contracting and expanding foliations for Anosov systems (see [11, 12]). Let η be the decreasing partition, whose elements are

$$\eta_J = (I, \Delta_{j_1}, \Delta_{j_2j_1}, \dots, \Delta_{j_n \dots j_2j_1}, \dots)$$

where $G\Delta_{j_n j_{n-1} \dots j_1} = \Delta_{j_{n-1} \dots j_1}$; J is used to denote the entire sequence (j_1, j_2, \dots) ; I is as above, and

$$\Delta_{j_n j_{n-1} \dots j_1} = \Delta_{j_n} \cap G^{-1}\Delta_{j_{n-1}} \cap \dots \cap G^{-(n-1)}\Delta_{j_1}.$$

Lemma 2 of § 7 implies that for any η_J the conditional measure $\tilde{\nu}(\cdot|\eta_J)$ induced by $\tilde{\nu}$ on η_J satisfies the following inequalities

$$0 < c_2 < \tilde{\nu}(E|\eta_J)/\tilde{\nu}(\theta_1 E) < c_1 \tag{6}$$

where E is a measurable subset of η_J , $\theta_1: \tilde{I} \rightarrow I$ is the projection along the element of the increasing partition, and c_1, c_2 are independent from J .

When we consider the attractor Λ the elements η_J admit a natural geometric interpretation.

4. For a set $E \subset [0, 1]$ on the y -axis we shall use the symbol Π_E to denote the strip $\{(x, y): x \in [0, 1], y \in E\}$.

Let τ be the map induced by F on the strip Π_I where $I = [t^{-1}, t]$ as in § 3. The horizontal foliation $\xi = \{\xi_y\}$ is invariant under τ and we have

$$\tau\xi_y \subset \xi_{Gy}. \tag{7}$$

The partition of I into the intervals $\Delta_{i\pm}$ generates the partition of Π_I into the strips $\Pi_{i\pm}$ (we use $\Pi_{i\pm}$ instead of $\Pi_{\Delta_{i\pm}}$ in order to simplify the notation). Let us consider the images $\tau\Pi_{i\pm}$. Figure 3 illustrates the map F^2 with the images $\tau\Pi_{2\pm} = F^2\Pi_{2\pm}$ shaded in.

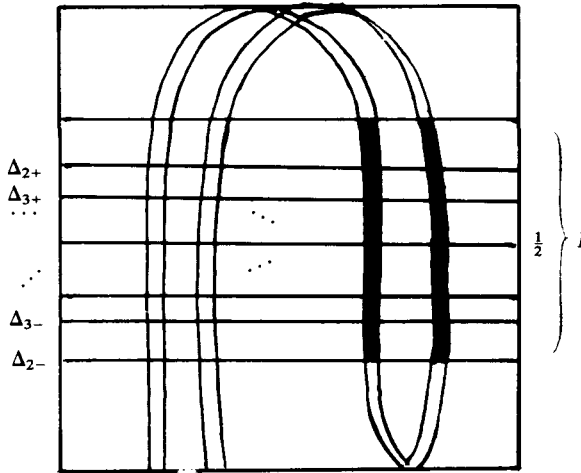


FIGURE 3

Let a be the fixed point of the linear map

$$x \rightarrow \frac{1}{4} + 2\nu + \lambda(x - \frac{1}{2}).$$

The action of F on $\Pi_{[0,1]}$ may be written in the form

$$(x - a, y) \rightarrow (\lambda(x - a), g(y)).$$

Thus $(a, 0)$ is a fixed saddle point of the map F , its stable manifold W^s coincides with the x -axis, and the intersection of its unstable manifold W^u with the strip $\Pi_{[0,1]}$ is a segment of the line $x = a$. As F acts diffeomorphically on $\Pi_{[0,1]}$, we see that the images $\tau\Pi_{n\pm}, \tau\Pi_{m\pm}$ do not intersect for $n \neq m$, that $\tau\Pi_{n+} \cap \tau\Pi_{n-} = \emptyset$ for small n but $\tau\Pi_{n+} \cap \tau\Pi_{n-} \neq \emptyset$ for large n , and $\tau\Pi_n$ accumulate to W^u when n tends to infinity (see figure 4).

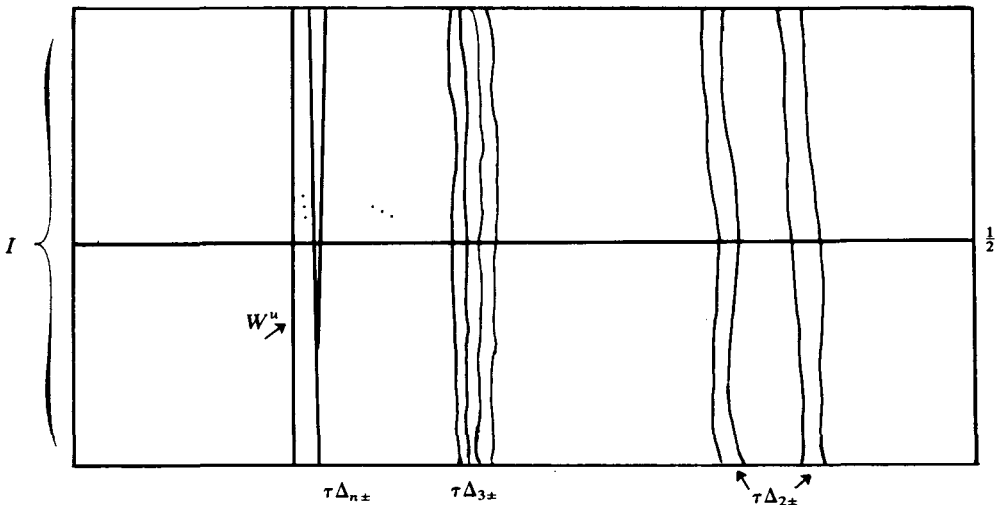


FIGURE 4

We shall denote the map $\tau|_{\Pi_{n\pm}}$ by $\tau_{n\pm}$. Let us consider the action of $D\tau_{n\pm}$. The derivative of F equals

$$DF(x, y) = \begin{pmatrix} \lambda & \varphi'(y) \\ 0 & g'(y) \end{pmatrix}.$$

The choice of δ in the definition of $\varphi(y)$ implies that all intervals $\Delta_{n\pm}$ with $n \geq 3$ are inside the domain where $\varphi'(y) \equiv 1$, and besides

$$\varphi'(y) \equiv 0 \quad \text{for } y \in [0, t^{-1}] \cup [t, 1].$$

Thus for $i \neq 2$ we have $\varphi'(y) = 1$, for any $y \in \Delta_i$, and for all $i = 2, 3, \dots$ we have

$$\varphi'(g^k y) = 0, \quad k = 1, 2, \dots, n_i - 1.$$

Therefore we obtain for $n \geq 3$

$$\begin{aligned} D\tau_{n\pm}(x, y) &= \begin{pmatrix} \lambda^{n-1} & 0 \\ 0 & \prod_{k=1}^{n-1} g'(g^k y) \end{pmatrix} \cdot \begin{pmatrix} \lambda & 1 \\ 0 & g'(y) \end{pmatrix} \\ &= \begin{pmatrix} \lambda^n & \lambda^{n-1} \\ 0 & (g^n)'(y) \end{pmatrix}. \end{aligned} \tag{8}$$

For $n = 2$ we have

$$\begin{aligned} D\tau_{2\pm}(x, y) &= \begin{pmatrix} \lambda & 0 \\ 0 & g'(gy) \end{pmatrix} \cdot \begin{pmatrix} \lambda & \varphi'(y) \\ 0 & g'(y) \end{pmatrix} \\ &= \begin{pmatrix} \lambda^2 & \lambda\varphi'(y) \\ 0 & (g^2)'(y) \end{pmatrix}. \end{aligned} \tag{9}$$

Let $A = g'(0)$ (for $g \in \mathcal{A}_\varepsilon A$ is close to 4). It is easy to check that there are constants $c_3, c_4 > 0$, such that for $n = 2, 3, \dots$ and for any $y \in \Delta_{n\pm}$ the following inequality holds

$$c_3(\sqrt{A})^n < |G'(y)| = |(g^n)'(y)| < c_4(\sqrt{A})^n. \tag{10}$$

Let us consider the images $D\tau_n K$ of the cone

$$K = \{(\xi, \eta) : |\xi/\eta| < 1\}$$

in the tangent space TS . We have

$$D\tau_2(\xi, \eta) = (\xi_1, \eta_1) = (\lambda^2 \xi + \lambda \varphi'(y)\eta, ((g^2)'(y))\eta).$$

Thus

$$|\xi_1/\eta_1| = \left| \frac{\lambda\varphi'(y)}{(g^2)'(y)} + \frac{\lambda^2}{(g^2)'(y)} \frac{\xi}{\eta} \right| < \lambda \frac{1+\lambda}{Ac_3}. \tag{11}$$

For $n > 2$ we have

$$D\tau_n(\xi, \eta) = (\xi_1, \eta_1) = (\lambda^n \xi + \lambda^{n-1} \eta, (g^n)'(y)\eta).$$

Thus

$$|\xi_1/\eta_1| < \lambda^{n-1}((1+\lambda)/(\sqrt{A})^n \cdot c_3). \tag{12}$$

If λ is sufficiently small the right-hand parts of (11) and (12) are less than 1 and we get $D\tau_{n\pm} K \subset K$ for all n .

Let us denote $\Lambda' = \bigcap_{p=0}^{\infty} \tau^p \Pi_I$ the attractor of the map $\tau: \Pi_I \rightarrow \Pi_I$. In order to study the structure of Λ' we use the ‘itinerary schemes’ method developed by Alexeyev [1]. In our situation Alexeyev’s theorem (theorem 1⁻ from [1]) may be formulated as follows.

Let us denote

$$D\tau_n = \begin{pmatrix} A_{11n} & A_{12n} \\ A_{21n} & A_{22n} \end{pmatrix}$$

$$\|A_{11n} - A_{12n} \cdot A_{22n}^{-1} \cdot A_{21n}\| = \mu_{1n}, \quad \|A_{22n}^{-1}\| = \mu_{2n}$$

$$\|A_{12n} \cdot A_{22n}^{-1}\| = a_{1n}, \quad \|A_{22n}^{-1} \cdot A_{21n}\| = a_{2n}.$$

Suppose the following conditions hold. There are two families of cones $K_1(x, y)$ and $K_2(x, y)$ in the tangent space so that

$$K_1 \ni (1, 0), \quad K_2 \ni (0, 1), \quad K_1 \cap K_2 = \emptyset$$

and for any n

$$(C_1) \quad D\tau_n K_1 \supset K_1, \quad D\tau_n K_2 \subset K_2.$$

There exist constants μ_1, μ_2, a_1, a_2 such that $\mu_{1n} < \mu_1, \mu_{2n} < \mu_2, a_{1n} < a_1, a_{2n} < a_2$ and the following inequalities hold:

$$(C_2) \quad \mu_1 < 1;$$

$$(C_3) \quad \sqrt{\mu_1 \cdot \mu_2} + \sqrt{a_1 \cdot a_2} < 1.$$

Then any itinerary $\Pi_I \xleftarrow{\tau} \Pi_{i_1} \xleftarrow{\tau} \Pi_{i_2} \xleftarrow{\tau} \dots$ determines the smooth curve

$$\Lambda'_{i_1 i_2 \dots i_n \dots} = \bigcap_{k=1}^{\infty} \tau^k \Pi_{i_k}$$

and Λ' coincides with the union of such curves. Besides any tangent vector $v = (\xi, \eta)$ to $\Lambda'_{i_1 i_2 \dots i_n \dots}$ belongs to $\tau_{i_1} K$.

If we set $K_1 = \{\mathbb{R}, 0\}, K_2 = K$ then (C₁) will hold for K_1 because of the invariance of the horizontal foliation and for K_2 because of (11) and (12). Now (8), (9), (10) imply that (C₂), (C₃) hold with

$$\mu_1 = \lambda^2, \quad \mu_2 = (c_3 A)^{-1}, \quad a_1 = 0, \quad a_2 = \lambda (c_3 A)^{-1}.$$

Thus Alexeyev’s theorem is applicable. The scheme $\{\Pi_n, \tau \Pi_n\}$ is not separable in the sense of [1] because

$$\tau(\Pi_{n+}) \cap \tau(\Pi_{n-}) \neq \emptyset \quad \text{for large } n.$$

However, the images $\tau_n K$ and $\tau_{n-} K$ do not intersect for all n . It follows from the subsequent inequalities where we use the notations of (11) and (12).

For $(x, y) \in \Pi_{n-}$ we have

$$\xi_1 / \eta_1 = [\lambda^{n-1} / (g^n(y))'] (1 + \lambda (\xi / \eta)) < -\frac{\lambda^{n-1}}{c_4 \cdot 2^n} (1 - \lambda) \tag{13}$$

while for $(x, y) \in \Pi_{n+}$ we have

$$\xi_1 / \eta_1 > \frac{\lambda^{n-1}}{c_4 \cdot 2^n} (1 - \lambda). \tag{14}$$

5. Let us denote by $\pi': \tilde{I} \rightarrow \Lambda'$ the map analogous to π , namely for $\tilde{x} = (x_0, x_1, x_2, \dots)$ where $x_n \in G^{-1}x_{n-1}$,

$$\pi'(x) = \xi_{x_0} \cap \tau \xi_{x_1} \cap \tau^2 \xi_{x_2} \cap \dots$$

The map π' induces a homomorphism of dynamical systems

$$\pi'_* : (\tilde{I}, \tilde{G}, \tilde{\nu}) \rightarrow (\Lambda', \tau, \pi'_* \tilde{\nu}).$$

LEMMA 1. π'_* is an isomorphism.

Proof. For a given $J = (j_1, j_2, \dots, j_n \dots)$ consider the element $\eta_J \in \eta$ of the partition defined in § 3. It follows from the definitions of η_J and Λ'_J that $\pi'(\eta_J) = \Lambda'_K$ and for $\tilde{z} = (z_0, z_1, \dots) \in \eta_J$, $\pi'(\tilde{z})$ coincides with the projection of $z_0 \in I$ to the curve Λ'_J along the x axis. In consequence of (6) and of estimates (11) and (12) we obtain that the induced measures $\pi'_* \tilde{\nu}(\cdot | \eta_J)$ are absolutely continuous with respect to the Lebesgue measure (normalized length) on Λ'_J .

Let

$$\sigma_J = \{\tilde{z} : \text{card } \pi'^{-1}(\pi' \tilde{z}) > 1\}$$

and let

$$\pi'(\sigma_J) = S_J \subset \Lambda'_J.$$

We show that $\text{mes } S_J = 0$ (mes is the Lebesgue measure Λ'_J). We have

$$S_J = \bigcup_{I \neq J} \Lambda'_I \cap \Lambda'_J = \bigcup_{k=1}^{\infty} \bigcup_{i_k \neq j_k} (\Lambda'_{i_k i_{k+1} \dots} \cap \Lambda'_{j_k j_{k+1} \dots}). \tag{15}$$

Let us denote $\tau(\Lambda' \cap \Pi_{n\pm})$ by $\Lambda'_{n\pm}$. As τ is non-singular with respect to Lebesgue measure on Λ'_K we obtain taking into account that $\Lambda'_m \cap \Lambda'_k \neq \emptyset$ only for $m = n_{\pm}$, $k = n_{\pm}$, that it suffices to check the equality

$$\text{mes}(\Lambda'_{n_+ j_2 j_3 \dots} \cap \Lambda'_{n_-}) = 0. \tag{16}$$

For $\lambda < \frac{1}{2}$ the intersection $\Lambda \cap \xi_y$ is a Cantor set of zero Lebesgue measure for any ξ_y . Now it follows from (13) and (14) that the angles between the curves $\Lambda'_{n_+ j_2 j_3 \dots}$ and $\Lambda'_{n_- i_2 i_3 \dots}$ are uniformly bounded away from zero. This implies (16). Thus $\text{mes } S_J = 0$. Thereby we have

$$\tilde{\nu}\{y(\tilde{z}_J | \eta_J) = 0 \text{ for any } \eta_J\}$$

and the Fubini theorem gives

$$\tilde{\nu}\{\tilde{z} \in \tilde{I} | \text{card } \pi'^{-1}(\pi' \tilde{z}) > 1\} = 0,$$

and thus π'_* is an isomorphism. □

Let $\tilde{\Delta}_{n\pm} = \{\tilde{y} = (y_0, y_1, \dots) \in \tilde{Y} | y_0 \in \Delta_{n\pm}\}$. If we neglect a set of zero $\tilde{\mu}$ -measure we can represent \tilde{Y} in the following form

$$\tilde{Y} = \bigcup_n \bigcup_{0 \leq k \leq n-1} \tilde{g}^k \tilde{\Delta}_{n\pm}. \tag{17}$$

As $\pi' \tilde{I}$ coincides with $\pi\{\tilde{y} = (y_0, y_1, \dots) | y_0 \in I\}$ we have

$$\Lambda = \bigcup_n \bigcup_{0 \leq k \leq n-1} F^k(\Lambda' \cap \Pi_{n\pm}). \tag{18}$$

Since F^k are non-singular with respect to Lebesgue measure on Λ'_j we obtain an isomorphism $(\tilde{Y}, \tilde{g}, \tilde{\mu}) \approx (\Lambda, F, \pi_*\tilde{\mu})$ where $\pi_*\tilde{\mu}$ has absolutely continuous conditional measures on the curves $F^i\Lambda'_j$ constituting the attractor Λ . According to the results of Ledrappier [10] the automorphism $(\tilde{Y}, \tilde{g}, \tilde{\mu})$ is Bernoulli. Thus $(\Lambda, F, \pi_*\tilde{\mu})$ is also Bernoulli and theorem 1 is proved.

Remark Although (18) defines Λ as a union of Λ'_j and their images $F^i\Lambda'_j$, it is easy to verify that every leaf of Λ is a smooth curve which connects the top of the square S to its bottom. These curves may be identified as

$$\Lambda_{p_1 p_2 p_3 \dots} = \pi\{\tilde{y} = (y_0, y_1, y_2, \dots) \mid y_0 \in [0, 1], y_1 = g_{p_1}^{-1}(y_0), y_2 = g_{p_2}^{-1}(y_1), \dots\}$$

where $p_i \in \{0, 1\}$, $g_0^{-1}(x) \in [0, \frac{1}{2}]$, $g_1^{-1}(x) \in [\frac{1}{2}, 1]$.

6. Now we turn to the proof of theorem 2. We consider again a map of the square given by (1)

$$F:(x, y) \rightarrow (\varphi(y) + \lambda(x - \frac{1}{2}), g(y))$$

where $\varphi(y)$ is the same as in § 2 and $g \in C^3([0, 1], [0, 1])$ which satisfies (i) but instead of (ii) we assume $\max g(y) = g(c) < 1$.

Invariant measures for such mappings were studied in [8]. It follows from [8] that for the family $y \rightarrow ay(1 - y)$ and for any family $y \rightarrow g_a(y)$ sufficiently close to $y \rightarrow ay(1 - y)$ in $C^3([0, 4] \times [0, 1], [0, 1])$ there is a set of parameter values \mathfrak{M} such that $\text{mes } \mathfrak{M} > 0$ and for $a \in \mathfrak{M}$ the map g_a has an absolutely continuous invariant measure. In contrast with the situation considered in §§ 2–5 we cannot prove the existence of an absolutely continuous Bernoulli measure for any $\lambda \in [0, \lambda_0]$ although we think it is true. Instead we will prove it for λ belonging to some subset of full measure. Let a family $\{g_a\}$ and a set \mathfrak{M} be as above.

THEOREM 2. *For the family*

$$F_a:(x, y) \rightarrow (\varphi(y) + \lambda(x - \frac{1}{2}), g_a(y))$$

there is a set $\Xi \subset [0, \lambda_0]$ such that $\text{mes } \Xi = \lambda_0$ and for any $\lambda \in \Xi$ the subset \mathfrak{M}' of \mathfrak{M} defined by

$$\mathfrak{M}' = \{a : \text{the attractor of } F_a \text{ admits an a.c.B. measure}\}$$

satisfies $\text{mes } (\mathfrak{M} \setminus \mathfrak{M}') = 0$.

We begin with constructing the induced map $G:I \rightarrow I$ and the corresponding two-dimensional map $\tau: \Pi_I \rightarrow \Pi_I$ induced by F on the strip Π_I . This is done exactly as in §§ 3 and 4 but now G has only a finite number of monotone branches $G_k:D_k \rightarrow I$ and one central parabolic branch. Correspondingly the image $\tau\Pi_I$ consists of a finite number of nearly vertical strips $\tau\Pi_{D_k}$ and of a narrow horseshoe (figure 5). Certainly it suffices to construct the desired attractor for τ .

An absolutely continuous invariant measure $\mu(dy)$ for G was constructed in [8] with the help of an auxiliary map $T:I \rightarrow I$ satisfying the following conditions. There

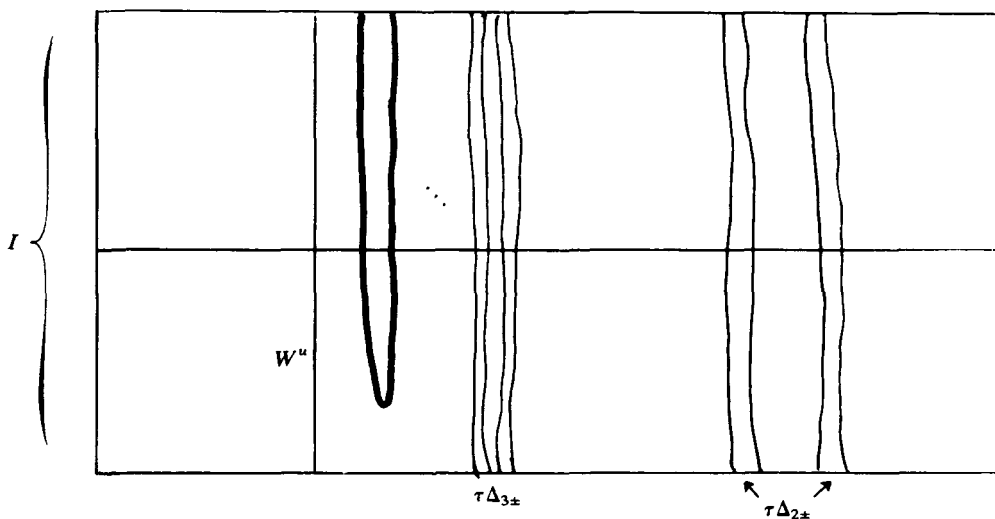


FIGURE 5

is a countable number of intervals $\Delta_i \subset I$ such that:

- 1) $\text{int } \Delta_i \cap \text{int } \Delta_j = \emptyset$;
- 2) the Lebesgue measure of $\bigcup \Delta_i$ equals $\text{mes } I$;
- 3) $T|\Delta_i = T_i = G^{n_i}$ are diffeomorphisms, $\text{Im } T_i = I$;
- 4) $\sup_i \max_{x \in \Delta_i} |D^2T(x)/DT(x)|(\text{mes } \Delta_i) < c_0$;
- 5) $|DT(x)| > L \gg 1$;
- 6) $\sum_i n_i \text{mes } \Delta_i < \infty$.

According to [13] there exists a T -invariant Bernoulli measure $\nu(dy)$ with the continuous density $\chi(y) > 0$ with respect to dy . ν generates the G -invariant measure μ by formula (5).

For $Z_0 = I \setminus \bigcup \Delta_i$ set

$$Z = \bigcup_{n \geq 0} G^{-n} \left(\bigcup_{k \geq 0} G^k Z_0 \right).$$

Then $\mu(Z) = 0$, $G^{-1}(Z) = G(Z) = Z$. Consider $X = I \setminus Z$. Then the endomorphisms (I, G, μ) and (X, G, μ) are isomorphic and for any $x \in X$, $k \in \mathbb{Z}$ and $x' \in G^k(x)$ $T(x)$ is defined. We shall use the notation $(\tilde{I}, \tilde{T}, \tilde{\mu})$ for the natural extension of T , and $(\tilde{X}, \tilde{G}, \tilde{\mu})$ for the natural extension of G . Let \mathcal{C} be the corresponding two-dimensional map which coincides with $\mathcal{C}_i = \tau^{n_i}$ on Π_{Δ_i} .

Let us denote by $\Delta = \bigcap_{n=0}^{\infty} \tau^n \Pi_I$ the attractor for the map τ ; by $\Lambda' = \bigcap_{n=0}^{\infty} \mathcal{C}^n \Pi_I$ the attractor for \mathcal{C} . We define the map $\pi: \tilde{X} \rightarrow \Lambda$ by

$$\pi(\tilde{x} = (x_0, x_1, x_2, \dots)) = \xi_{x_0} \cap \tau \xi_{x_1} \cap \tau^2 \xi_{x_2} \cap \dots$$

and the map $\pi': \tilde{I} \rightarrow \Lambda'$ by

$$\pi'(\tilde{y} = (y_0, y_1, y_2, \dots)) = \xi_{y_0} \cap \mathcal{C} \xi_{y_1} \cap \mathcal{C}^2 \xi_{y_2} \cap \dots$$

7. Let \mathcal{A} be the partition of \tilde{I} into cylinders $(\Delta_i, T^{-1}\Delta_i, T^{-2}\Delta_i, \dots)$. Consider the corresponding decreasing partition $\bigvee_{n=1}^{\infty} \tilde{T}^n = \eta$ and the increasing one

$\bigvee_{n=0}^{\infty} T^{-n}\mathcal{A} = \gamma$. Let us denote by α the measure induced by $\tilde{\nu}$ on the factor space \tilde{I}/η . Let C_n be an element of η ; $\nu(\cdot|C_n)$ the conditional measure induced by $\tilde{\nu}$ on C_n ; C_γ an element of γ ; $\nu(\cdot|C_\gamma)$ the conditional measure on C_γ . The following assertion may be considered as a simple version of several theorems about systems with absolutely continuous foliations (see e.g. [3, 11]).

LEMMA 2. $\tilde{\nu}$ is equivalent to the direct product $\nu \times \alpha$. There are constants $c_1, c_2 > 0$ so that for any ν -measurable set $\Delta \subset I$ and for any α -measurable set $\Gamma \subset \tilde{I}/\eta$ one has

$$c_2 < \frac{\tilde{\nu}(\Delta \times \Gamma)}{\nu(\Delta) \cdot \alpha(\Gamma)} < c_1. \tag{19}$$

Proof. As $\tilde{T}^{-n}\mathcal{A} = \{(T^{-n}\Delta_i, T^{-(n+1)}\Delta_i, \dots)\}$ we have

$$\bigvee_0^N \tilde{T}^{-n}\mathcal{A} = \{\Delta_{i_0 i_1 \dots i_N}, T^{-1}\Delta_{i_0 i_1 \dots i_N}, \dots\}$$

where

$$\Delta_{i_0 i_1 \dots i_N} = \Delta_{i_0} \cap T^{-1}\Delta_{i_1} \cap \dots \cap T^{-N}\Delta_{i_N}.$$

The intersection $\Delta_{i_0} \cap T^{-1}\Delta_{i_1} \cap \dots \cap T^{-N}\Delta_{i_N} \cap \dots$ is a point for any sequence $(i_0, i_1, i_2, \dots, i_N, \dots)$ because of 5). Thus we obtain

$$\gamma = \bigvee_0^{\infty} \tilde{T}^{-n}\mathcal{A} = \{(y, T^{-1}y, T^{-2}y, \dots)\}.$$

On the other hand

$$\tilde{T}^n\mathcal{A} = \{(I, I, \dots, I, \Delta_i, T^{-1}\Delta_i, \dots)\}.$$

n times

Hence

$$\bigvee_1^N \tilde{T}^n\mathcal{A} = \{(I, \Delta_{i_1}, \Delta_{i_2 i_1}, \dots, \Delta_{i_N i_{N-1} \dots i_1}, T^{-1}\Delta_{i_N i_{N-1} \dots i_1}, \dots)\}$$

and we obtain

$$\eta = \bigvee_1^{\infty} \tilde{T}^n\mathcal{A} = \{(I, \Delta_{i_1}, \Delta_{i_2 i_1}, \dots, \Delta_{i_N i_{N-1} \dots i_1}, \Delta_{i_{N+1} i_N \dots i_1}, \dots)\}.$$

Any point $\tilde{y} = (y_0, y_1, y_2, \dots, y_n, \dots) \in \tilde{I}$ may be uniquely represented as $\tilde{y} = C_n \subset C_\gamma$ where the indices i_1, i_2, \dots in the definition of C_n are determined by $y_n \in \Delta_{i_n \dots i_1}$ and $C_\gamma = (y_0, T^{-1}y_0, \dots)$. This gives the direct product structure in \tilde{I} . Let $\Delta \subset I$ be a measurable set, $\Gamma_\Delta = (\Delta, T^{-1}\Delta, T^{-2}\Delta, \dots)$, $\Gamma_{i_1 i_2 \dots i_n}$ the cylinder in \tilde{I} defined by

$$\Gamma_{i_1 i_2 \dots i_n} = (I, \Delta_{i_1}, \Delta_{i_2 i_1}, \dots, \Delta_{i_n \dots i_2 i_1}, T^{-1}\Delta_{i_n \dots i_2 i_1}, T^{-2}\Delta_{i_n \dots i_2 i_1}, \dots).$$

Let us denote the intersection $\Gamma_\Delta \cap \Gamma_{i_1 i_2 \dots i_n}$ by $P(\Delta, [i_1, \dots, i_n])$. We shall use the notation $T_{i_1 i_2 \dots i_k}$ for the unique monotone branch of T^k which maps

$$\Delta_{i_k} \cap T^{-1}\Delta_{i_{k-1}} \cap \dots \cap T^{-(k-1)}\Delta_{i_1}$$

onto I . The assertion of lemma 2 will be proved if we show that independently from Δ and $[i_1 i_2 \dots i_n]$ the following inequality holds:

$$c_2 < \frac{\tilde{\nu}(P(\Delta, [i_1 i_2 \dots i_n]))}{\tilde{\nu}(\Gamma_\Delta) \cdot \tilde{\nu}(\Gamma_{i_1 i_2 \dots i_n})} < c_1. \tag{20}$$

According to the definition of $\tilde{\nu}$ we have

$$\begin{aligned} \tilde{\nu}(\Gamma_\Delta) &= \nu(\Delta), \\ \tilde{\nu}(\Gamma_{i_1 i_2 \dots i_n}) &= \nu(T_{i_1 i_2 \dots i_n}^{-n}(I)), \\ \tilde{\nu}(P(\Delta, [i_1 i_2 \dots i_n])) &= \nu(T_{i_1 i_2 \dots i_n}^{-n}(\Delta)). \end{aligned}$$

Thus (20) is equivalent to

$$c_2 < \frac{\nu(T_{i_1 i_2 \dots i_n}^{-n}(\Delta))}{\nu(T_{i_1 i_2 \dots i_n}^{-n}(I))} \cdot \frac{\nu(\Delta)}{\nu(I)} < c_1. \tag{21}$$

Now it follows from properties 4) and 5) of T that 4) holds not only for T_i but also for all their compositions $T_{i_1} \circ T_{i_2} \circ \dots \circ T_{i_n}$ with a constant c_5 independent from $[i_1 i_2 \dots i_n]$, see for example lemma 5 of [8]. Thus $b_1 < DT^n(x)/DT^n(y) < b_2$ for some $b_1, b_2 > 0$; and for any $x, y \in \Delta_{i_n \dots i_1}$. This implies an inequality similar to (21) but for Lebesgue measure. Since $0 < d_1 = \min_{y \in I} \chi(y) < \nu(A)/\text{mes } A < \max_{y \in I} \chi(y) = d_2$ for any measurable set A we obtain (21) which finishes the proof of lemma 2. □

COROLLARY. Let $\theta_1: \tilde{I} \rightarrow I$ be the projection along C_γ and $\theta_2: \tilde{I} \rightarrow \tilde{I}/\eta$ the projection along C_η . Then for any element C_η and for any measurable $\Delta \subset C_\eta$

$$c_2 < \frac{\tilde{\nu}(\Delta|C_\eta)}{\nu(\theta_1 \Delta)} < c_1.$$

Similarly for any measurable $\Gamma \subset \tilde{I}/\eta$

$$c_2 < \frac{\tilde{\nu}(\Gamma|C_\gamma)}{\alpha(\theta_2 \Gamma)} < c_1.$$

In particular we obtain that η and γ are absolutely continuous partitions and the intersection $\eta \wedge \gamma$ is trivial.

8. The following construction which clarifies the relation between $(\tilde{X}, \tilde{G}, \tilde{\mu})$ and $(\tilde{I}, \tilde{T}, \tilde{\nu})$ is similar to the ‘tower’ construction for automorphisms. For any Δ_i consider the sets

$$A_{ik} = G^k \Delta_i, \quad 0 \leq k \leq n_i,$$

and let $\mathcal{X} = \bigcup_{i,k} A_{ik}$ be the disjoint sum of A_{ik} . We define the projection $p: \mathcal{X} \rightarrow I$ by assigning to any $u \in A_{ik}$ its image under the natural inclusion $p(u) \in G^k \Delta_i \subset I$. Then we define a map $\mathcal{F}: \mathcal{X} \rightarrow \mathcal{X}$. If $u \in A_{ik}, k < n_i - 1$, then

$$\mathcal{F}(u) = p^{-1}(G(u)) \cap A_{i, k+1};$$

if $u \in A_{i, n_i - 1}$ then

$$\mathcal{F}(u) = p^{-1}(G(u)) \cap \bigcup_j A_{j0}.$$

The construction implies $G \circ p = p \circ \mathcal{F}$. Notice that if we identify $\bigcup_i A_{i0}$ with I , then the map induced by \mathcal{F} on $I \subset \mathcal{X}$ coincides with T .

Let us define the measure ρ on \mathcal{X} by

$$\rho(U \subset A_{ik}) = \nu(\mathcal{F}^{-k}U).$$

Since ν is T -invariant ρ is \mathcal{F} -invariant, and the definition of μ implies that $p_*\rho = \mu$. Consider the natural extension $(\tilde{\mathcal{X}}, \tilde{\mathcal{F}}, \tilde{\rho})$. It follows from

$$G \circ p = p \circ \mathcal{F}$$

that the map $\tilde{p}: \tilde{\mathcal{X}} \rightarrow \tilde{X}$ defined by $\tilde{p}(u_0, u_1, \dots) = (pu_0, pu_1, \dots)$ satisfies

$$\tilde{p} \circ \tilde{\mathcal{F}} = \tilde{G} \circ \tilde{p}.$$

LEMMA 3. \tilde{P} is an isomorphism of dynamical systems $(\tilde{\mathcal{X}}, \tilde{\mathcal{F}}, \tilde{\rho})$ and $(\tilde{X}, \tilde{G}, \tilde{\mu})$.

Proof. We show that there is a subset \tilde{M} of \tilde{X} such that $\tilde{\mu}(\tilde{M}) = 0$ and $\tilde{p}: \tilde{\mathcal{X}} \rightarrow \tilde{X} \setminus \tilde{M}$ is one-to-one. For $x \in \Delta_k$ set $n(x) = n_k$. We define $\tilde{M} = \bigcup_{i=0}^{\infty} \tilde{M}_i$ where the set \tilde{M}_i is defined by

$$\tilde{x} = (x_0, x_1, \dots, x_i, x_{i+1}, \dots) \in \tilde{M}_i$$

iff there is a sequence $j_n \rightarrow \infty$ so that $n(x_{j_n}) > j_n - i$. Since $n(x_{j_n}) > j_n - i \rightarrow \infty$ we obtain for $\tilde{x} \in \tilde{M}_i$ that $x_{j_n} \in \Delta_{p_n}$ where $p_n \rightarrow \infty$ when $n \rightarrow \infty$. Now

$$\tilde{\nu}\{\tilde{x}: x_{j_n} \in \Delta_{p_n}\} = \mu(\Delta_{p_n})$$

hence it is sufficient to show that

$$\lim_{N \rightarrow \infty} \{x: n(x) > N\} = 0.$$

But this follows from the convergence of $\sum_{k=1}^{\infty} \mu(\Delta_k)$. Thus $\mu(\tilde{M}_i) = 0$ and consequently $\tilde{\mu}(\tilde{M}) = 0$.

We shall use the following decomposition property which follows from the inductive construction of the maps T_i in [8].

PROPERTY D. Let $T|\Delta_i = T_i = G^{n_i}$. Then for every $k < n_i$ there are $T_{j_1}, \dots, T_{j_r}, r \geq 1$, so that the restriction of G^{n_i-k} to $G^k \Delta_i$ may be represented as

$$G^{n_i-k} = T_{j_1} \circ T_{j_2} \circ 1 \cdots \circ T_{j_r}.$$

Consider a point $\tilde{x} = (x_0, x_1, \dots) \in \tilde{X} \setminus \tilde{M}$. From $\tilde{x} \notin \tilde{M}_0$ we obtain that there exists

$$m_0(\tilde{x}) = \max \{m: n(x_m) > m\}.$$

Then property D implies $n(x_{m_0+k}) \leq k, k = 1, 2, \dots$. In particular, $n(x_{m_0+1}) = 1$.

For $k \geq 1$ we define

$$m_k = \max_{i > m_{k-1}} \{i: n(x_i) = i - m_{k-1}\}.$$

Property D implies

$$i - n(x_i) \geq m_{k-1} \quad \text{for } m_{k-1} < i \leq m_k.$$

Let $\tilde{x} = (x_0, x_1, \dots, x_{n_0}, \dots, x_{n_1}, \dots, x_{n_2}, \dots)$ where (n_0, n_1, n_2, \dots) is the sequence defined above, $x_{n_i} \in \Delta_{m_i}$. Then $\tilde{u} = (u_0, u_1, \dots, u_{n_0}, \dots, u_{n_1}, \dots, u_{n_2}, \dots) \in \mathcal{X}$ defined by

$$u_{n_i} = x_{n_i} \in A_{i,0}, \quad u_{n_i-k} = \mathcal{F}^k u_{n_i}, \quad k = 1, \dots, n_i - n_{i-1}$$

satisfies $\tilde{p}\tilde{u} = \tilde{x}$, hence \tilde{p} is onto $\tilde{X} \setminus \tilde{M}$. Let $\tilde{p}\tilde{u}' = \tilde{x}$, where $\tilde{u}' = (j'_0, j'_1, \dots)$. Let $n'_0 = \min\{n: u'_n \in I\}$ and for $k \geq 1$

$$n'_k = \min\{n: n > n_{k-1}, u'_n \in I\}.$$

The definition of the sequence (n_0, n_1, \dots) implies $Tx_{n_i} = x_{n_i}$ and the construction of $(\mathcal{X}, \mathcal{F})$ implies $Tx_{n_i} = x'_{n_i-1}$. Besides it follows from the definitions of n_0 and of n'_0 that $n(x_{n_0}) > n_0$ and $n(x_{n'_0}) > n'_0$. Suppose $n'_0 < n_0$. Then $n_0 \in [n'_k, n'_{k+1}]$ and using property D we obtain $n(x_{n_0}) \leq n_0 - n'_k < n_0$, a contradiction. The same arguments show that $n_0 < n'_0$ contradicts to $n(x_{n'_0}) > n'_0$. Therefore, $n_0 = n'_0$. Analogously using property D one obtains $n_k = n'_k$ for all k . Since

$$\mathcal{F}: A_{ik} \rightarrow A_{i+ik} \quad \text{and} \quad G: G^k \Delta_i \rightarrow G^{k+1} \Delta_i$$

are homeomorphisms for $0 \leq k \leq n_i - 1$ we obtain $\tilde{u}' = \tilde{u}$, and lemma 3 is proved. \square

Let us denote

$$\tilde{B} = \left\{ \tilde{u} = (u_0, u_1, \dots) \in \tilde{\mathcal{X}} \mid u_0 \in \bigcup_i A_{i0} \right\}.$$

It follows from the proof of lemma 3 that $\tilde{p}\tilde{B}$ coincides with $\{\tilde{x} = (x_0, x_1, \dots) : n_0(\tilde{x}) = 0\}$. As the automorphism induced by \mathcal{F} on $(\tilde{B}, \tilde{\rho})$ coincides with $(\tilde{I}, \tilde{T}, \tilde{\nu})$ we obtain that \tilde{p} is an isomorphism between $(\tilde{I}, \tilde{T}, \tilde{\nu})$ and the automorphism induced by \tilde{G} on $(\tilde{p}\tilde{B}, \tilde{\mu})$. For

$$\tilde{y} = (y_0, y_1, y_2, \dots) \in \tilde{I}, \quad \text{where } y_i = G^{n_{i+1}} y_{i+1},$$

we have

$$\tilde{p}\tilde{y} = (y_0, G^{n_1-1} y_1, \dots, G y_1, y_1 = G^{n_2} y_2, G^{n_2-1} y_2, \dots, y_2 = G^{n_3} y_3, \dots).$$

Remark. The representation of $(\tilde{I}, \tilde{T}, \tilde{\nu})$ as an induced automorphism, the triviality of the intersection $\eta \wedge \gamma$ proved in § 7, and the K -property of T allow us to apply the results of Gurevich [6] and Blanchard [5] which give the following corollary.

COROLLARY. $(\tilde{X}, \tilde{G}, \tilde{\mu})$ is a K -system.

This follows of course from the B-property proved by Ledrappier in [10]. The above construction gives another way to construct the unstable foliation.

9. Let us consider a point

$$x = (x_0 = x_{n_0}, x_1, \dots, x_{n_1}, \dots, x_{n_2}, \dots) \in \tilde{X}$$

and its pre-image

$$\tilde{p}^{-1} \tilde{x} = (x_{n_0}, x_{n_1}, x_{n_2}, \dots) \in \tilde{I}.$$

For the projections π and π' defined in § 6 we have

$$\begin{aligned} \pi(\tilde{x}) &= \xi_{x_0} \cap \tau \xi_{x_1} \cap \dots \cap \tau^{n_1} \xi_{x_{n_1}} \cap \dots \cap \tau^{n_1+n_2} \xi_{x_{n_2}} \cap \dots \\ &= \xi_{x_{n_0}} \cap \tau^{n_1} \xi_{x_{n_1}} \cap \tau^{n_1+n_2} \xi_{x_{n_2}} \cap \dots \\ &= \pi'(\tilde{p}^{-1} \tilde{x}). \end{aligned}$$

Thus the image of \tilde{I} under π' coincides with its image under π when we consider \tilde{I} as the subset of \tilde{X} . Using the result of the previous section we represent $(\Lambda', \mathcal{C}, \pi'_* \tilde{\nu})$ as an induced system with respect to $(\Lambda, \tau, \pi_* \mu)$.

Now we check that $\{\Pi_{\Delta_i}, \mathcal{C}_i\}$ satisfy the conditions of Alexeyev's theorem. For $T: \Delta \rightarrow I$ and the corresponding $\mathcal{C}: \Pi_{\Delta} \rightarrow \Pi_I$ where $T = G'$, $\mathcal{C} = \tau'$, we can rewrite (8), (9) as

$$D\tau(x, y) = \begin{pmatrix} \lambda^m & \lambda^{m-1}\varphi'(y) \\ 0 & G'(y) \end{pmatrix}$$

where $y \in \Pi_{m\pm}$, $G(y) = g^m(y)$ and $\varphi'(y) \equiv 1$ for $m > 2$. Consider a point $(x_r, y_r) \in \Pi_{\Delta}$ and for $k \in [0, r]$ set

$$(x_k, y_k) = \tau^{r-k}(x_r, y_r)$$

in particular

$$(x_0, y_0) = \mathcal{C}(x_r, y_r) = \tau^r(x_r, y_r).$$

Then

$$D\mathcal{C}(x_r, y_r) = \prod_{i=1}^r \begin{pmatrix} \lambda^{m_i} & \lambda^{m_i-1}\varphi'(y_i) \\ 0 & G'(y_i) \end{pmatrix} = \begin{pmatrix} \lambda^{\sum_{i=1}^r m_i} & \lambda^{-1} \left[\lambda^{m_1}\varphi'(y_1) \prod_{i=2}^r G'(y_i) + \dots + \lambda^{\sum_{i=1}^{r-1} m_i}\varphi'(y_{r-1})G'(y_r) + \lambda^{\sum_{i=1}^r m_i}\varphi'(y_r) \right] \\ 0 & \prod_{i=1}^r G'(y_i) \end{pmatrix}$$

Consequently

$$D\mathcal{C}(x_r, y_r)(\xi, \eta) = (\xi_1, \eta_1) = \left(\lambda^{\sum_{i=1}^r m_i}\xi + \lambda^{-1}(\lambda^{m_1}\varphi'(y_1) \times \prod_{i=2}^r G'(y_i) + \dots + \lambda^{\sum_{i=1}^r m_i}\varphi'(y_r))\eta, \prod_{i=1}^r G'(y_i)\eta \right).$$

Hence

$$\frac{\xi_1}{\eta_1} = \lambda^{-1} \left(\frac{\lambda^{m_1}\varphi'(y_1)}{G'(y_1)} + \frac{\lambda^{m_1+m_2}\varphi'(y_2)}{G'(y_1)G'(y_2)} + \dots + \frac{\lambda^{\sum_{i=1}^r m_i}\varphi'(y_r)}{\prod_{i=1}^r G'(y_i)} \right) + \frac{\lambda^{\sum_{i=1}^r m_i}}{\prod_{i=1}^r G'(y_i)} \cdot \frac{\xi}{\eta}. \tag{22}$$

According to property D $G^k|G^{r-k}\Delta$ may be represented as a composition

$$G^k = T_{s_1} \circ \dots \circ T_{s_l}.$$

Since $|DT_j| > L$ for any j we have

$$|G'(y_1) \circ \dots \circ G'(y_k)| = |(G^k)'| > L^l.$$

Thus $|\xi_1/\eta_1|$ may be estimated as

$$|\xi_1/\eta_1| < \lambda^{m_1-1}/L(1-\lambda). \tag{23}$$

It follows from (23) that for any $\mathcal{C}_i = \mathcal{C}|_{\Pi_{\Delta_i}}$ the cone $|\xi/\eta| < 1$ is invariant under $D\mathcal{C}_i$. The remaining conditions of Alexeyev's theorem are also satisfied with

$$\mu_1 = \lambda^2, \quad \mu_2 = L^{-1}, \quad a_1 = 0, \quad a_2 = \lambda L^{-1}.$$

Therefore $\{\Pi_{\Delta_i}, \mathcal{C}_i\}$ form an itinerary scheme. This implies just as in § 4 that Λ' is the union of smooth curves $\Lambda'_J = \Lambda'_{j_1 j_2 \dots j_n \dots}$. Any Λ'_J is a graph of a function $x(y)$ satisfying $|dx/dy| < \lambda/L(1-\lambda)$. For an element C_J of the partition η defined by

$$C_J = (I, \Delta_{j_1}, \Delta_{j_2 j_1}, \dots) = \{(y, T_{j_1}^{-1}y, T_{j_2}^{-1} \circ T_{j_1}^{-1}y, \dots)\}, \quad y \in I,$$

we have $\pi'(C_J) = \Lambda'_J$ and $\pi'(y, T_{j_1}^{-1}y, \dots)$ coincides with the projection of y on Λ'_J along the x -axis. Since the conditional measure $\tilde{\nu}(\cdot|C_J)$ is equivalent to dy , the conditional measure $\pi'_*(\cdot|\Lambda'_J)$ is equivalent to Lebesgue measure on Λ'_J . Set $\Lambda'_{\Delta_i} = \Lambda' \cap \Pi_{\Delta_i}$.

The representation of (Λ', \mathcal{C}) as an induced system with respect to (Λ, τ) gives

$$\Lambda = \bigcup_i \bigcup_{k=0}^{n_i-1} \tau^k(\Lambda'_{\Delta_i}). \tag{24}$$

Notice that contrary to (18) all the summands in (24) are contained in Π_I . One can check that the leaves of Λ are smooth curves but their projections on the y -axis depend now on the curves; the turning points lie on ξ_{c_n} where c_n belong to the trajectory of the critical point c .

10. Now we take into account the dependence on a and λ . For a map

$$F_{a,\lambda}(x, y) = (\varphi(y) + \lambda(x - \frac{1}{2}), g_a(y))$$

let $(\tilde{X}_a, \tilde{G}_a, \tilde{\mu}_a)$ and $(\Lambda_{a,\lambda}, \tau_{a,\lambda}, \xi_{a\lambda}, \tilde{\mu}_a)$ be the corresponding systems. The preceding results were obtained for any $a \in \mathfrak{M}$ and any $0 < \lambda \leq \lambda_0$. Let

$$\sigma_{a,\lambda} = \{\tilde{x} \in \tilde{X}_a : \text{card } \pi^{-1}(\pi(x)) > 1\}.$$

LEMMA 4. (i) For any $a \in \mathfrak{M}$ and for any $\lambda \in (0, \lambda_0]$ either

$$\tilde{\mu}_a(\sigma_{a\lambda}) = 0 \quad \text{or} \quad \tilde{\mu}_a(\tilde{\sigma}_{a\lambda}) = \tilde{\mu}_a(\tilde{X}_a);$$

(ii) for any $a \in \mathfrak{M}$

$$\text{mes } \{\lambda : \tilde{\mu}_a(\sigma_{a\lambda}) = \tilde{\mu}_a(\tilde{X}_a)\} = 0.$$

Let us show that lemma 4 implies theorem 2. Consider the measure $da \times d\lambda$ on the direct product $\mathfrak{M} \times (0, \lambda_0]$. Let

$$\Xi = \{\lambda \in (0, \lambda_0] : \text{mes } \{a \in \mathfrak{M} : \tilde{\mu}_a(\sigma_{a\lambda}) > 0\} = 0\}.$$

It follows from the Fubini theorem that $\text{mes } \Xi = \lambda_0$. For $\lambda \in \Xi$ the set of parameter values a such that $F_{a,\lambda}$ admits an attractor with an absolutely continuous Bernoulli measure differs from \mathfrak{M} by a set of measure 0.

Proof of lemma 4. (1) First we fix a and λ . For any $\tilde{y} \in \sigma$ there exists some $\tilde{z} \in \sigma$ such that

$$\tilde{y} = (y_0, y_1, \dots) \neq \tilde{z} = (z_0, z_1, \dots)$$

but $\pi(\tilde{y}) = \pi(\tilde{z}) = M \in \Lambda$. Let

$$\begin{aligned} \sigma_n = \{y : \text{for some } \tilde{z} \pi\tilde{y} = \pi\tilde{z} \text{ and } y_i = z_i \\ \text{for } i = 0, 1, \dots, n-1, \text{ but } y_n \neq z_n\}. \end{aligned}$$

Consider the points

$$\tilde{y}^{-n} = (y_n, y_{n+1} \dots), \quad \tilde{z}^{-n} = (z_n, z_{n+1} \dots) \in \tilde{X},$$

$$M_1^{-n} = \pi(\tilde{y}^{-n}), \quad M_2^{-n} = \pi(\tilde{z}^{-n}) \in \Lambda.$$

Then we have $M_1^{-n} \in \xi_{y_n}, M_2^{-n} \in \xi_{z_n} \rightarrow M_k^{-n} \neq M_2^{-n}$, but $\tau(M_1^{-n}) = \tau(M_2^{-n})$. Let $\sigma_{n,k}$ be a subset of σ_n defined by $y_n \in \Delta_{k\pm}$ (the corresponding $z_n \in \Delta_{k\mp}$). Let $\tilde{y}_* = \lim_{m \rightarrow \infty} \tilde{y}_m$ for a sequence $\tilde{y}_m \in \sigma_{n,k}$, and let $y_{*n} = \lim_{m \rightarrow \infty} y_{mn}$ be the n th coordinate of \tilde{y}_* . For $\tilde{z}_m \in \sigma_{n,k}$ satisfying $\pi(\tilde{z}_m) = \pi(\tilde{y}_m)$ consider some limit point \tilde{z}_* . As $\text{dist}(z_{nm}, y_{nm}) > c$ for any m and $\tau(z_{nm}) = \tau(y_{nm})$ we obtain using the continuity of τ that $z_{*n} = \lim_{m \rightarrow \infty} z_{nm}$ is different from y_{*n} and $\tau(z_{*n}) = \tau(y_{*n})$. Hence $y_* \in \sigma_{nk}, \sigma_{nk}$ is closed, and $\sigma = \bigcup_{n,k} \sigma_{nk}$ is measurable.

It follows from the definition that $\tilde{G}\sigma_n = \sigma_{n+1}$. Thus σ is invariant. As $\tilde{\mu}$ is ergodic $\tilde{\mu}(\sigma) = 0$ or $\tilde{\mu}(\sigma) = \tilde{\mu}(\tilde{X})$ which proves (i).

(2) Now a is fixed but λ varies in $(0, \lambda_0]$. Then the sets $\sigma, \sigma_n, \sigma_{nk}$ depend on λ . Let \mathfrak{B} be the set of λ such that $\tilde{\mu}(\sigma_\lambda) = \tilde{\mu}(\tilde{X})$. As $\tilde{\mu}(\sigma_{\lambda n}) = \tilde{\mu}(\sigma_{\lambda n+1})$ for $\lambda \in \mathfrak{B}$ we have $\tilde{\mu}(\sigma_{\lambda_1}) > 0$. Let us define

$$\mathfrak{B}_{mk} = \{\lambda \in \mathfrak{B} : \tilde{\mu}(\sigma_{\lambda 1k}) \geq 1/m\}.$$

Then $\mathfrak{B} = \bigcup_{m,k} \mathfrak{B}_{mk}$. Suppose $\lambda_* = \lim_{n \rightarrow \infty} \lambda_n$ where $\lambda_n \in \mathfrak{B}_{mk}$ but $\lambda_* \notin \mathfrak{B}_{mk}$. Then

$$\tilde{\mu}(\sigma_{\lambda_* 1k}) = c_* \quad \text{for some } 0 \leq c_* < 1/m.$$

Let $d_{k\pm} = \{\tilde{x} = (x_0, x_1, x_2 \dots) : x_1 \in \Delta_{k\pm}\}$. If $\tilde{x} \in d_{k\pm} \setminus \sigma_{1k}$ then

$$\text{dist}[\tau(\pi(x_1, x_2 \dots)), \tau\Delta_{k\mp}] > 0.$$

Let

$$V_p = \{\tilde{x} : \text{dist}(\tau\pi(x_1, x_2, \dots), \tau\Delta_{k\mp}) \geq 1/p\}.$$

Then V_p is closed and

$$\tilde{\mu}\left(\bigcup_p V_p\right) = \mu(\Delta_k) = c_* > \mu(\Delta_k) - 1/m.$$

For any ε we can choose P_0 such that

$$\tilde{\mu}\left(\bigcup_{p \leq P_0} V_p\right) > \mu(\Delta_k) - c_* - \varepsilon.$$

Using the continuous dependence of π and τ from λ we can choose $t > 0$ such that any $\lambda \in (\lambda_* - t, \lambda_* + t)$ satisfies

$$\tilde{\mu}(\sigma_{\lambda 1k}) < \mu(\Delta_k) - (\mu(\Delta_k) - c_* - \varepsilon) < c_* + \varepsilon.$$

For ε sufficiently small and $\lambda = \lambda_n$ this contradicts to $\tilde{\mu}(\sigma_{\lambda_n 1k}) \geq 1/m$. Hence the sets \mathfrak{B}_{mk} are closed and \mathfrak{B} is measurable.

(3) Let $\text{mes } \mathfrak{B} > 0$. Then for some m, k $\text{mes } \mathfrak{B}_{mk} > 0$. For $\tilde{x} \in \sigma_{\lambda 1k+}$ we define

$$\Gamma_{\tilde{x}} = \{\lambda : \pi_\lambda(\tilde{x}) \in \pi_\lambda \sigma_{\lambda 1k-}\}.$$

Applying the Fubini theorem to the product of the measure spaces

$$(\tilde{X}, \tilde{\mu}) \times ([0, \lambda_0], d\lambda)$$

we see that $\tilde{\mu}\{\tilde{x} \in \sigma_{\lambda 1k+} : \text{mes } \Gamma_{\tilde{x}} > 0\} > 0$. We fix some $\tilde{p} = (p_0, p_{1+}, p_2, \dots)$ such that $\text{mes } \Gamma_{\tilde{p}} > 0$. Then to any $\lambda \in \Gamma_{\tilde{p}}$ there corresponds some \tilde{x} satisfying $\tilde{x} =$

$(x_0 = p_0, x_1 = p_{1-}, x_2, \dots), \pi_\lambda(\tilde{x}) = \pi_\lambda(\tilde{p})$. Now it is convenient to return to the inverse limit $([0, \bar{1}], \tilde{g})$ of the initial one-dimensional map $g: [0, 1] \rightarrow [0, 1]$ and to the initial map $F: (x, y) \rightarrow (\varphi(y) + \lambda(x - \frac{1}{2}), g(y))$ of the unit square. Let $q_{1+} = p_{1+}, q_2, q_3, \dots$ be the coordinates of the point $(p_{1+}, p_2, p_3, \dots) \in \tilde{X}$ considered as the point of $[\bar{0}, \bar{1}]$. Set $q_0 = g(q_{1+})$ and let the point $\tilde{q} = (q_0, q_{1+}, q_2, \dots) \in [\bar{0}, \bar{1}]$ correspond to $\tilde{p} \in \tilde{X}$. Similarly to any $\tilde{x} = (p_0, p_{1-}, x_2, x_3, \dots) \in \tilde{X}$ we associate $\tilde{y} = (q_0, q_{1-}, y_2, y_3, \dots) \in [\bar{0}, \bar{1}]$. For uniformity we shall use $\Gamma_{\tilde{q}}$ instead of $\Gamma_{\tilde{p}}$. According to our assumption $\text{mes } \Gamma_{\tilde{q}} > 0$ and for any $\lambda \in \Gamma_{\tilde{q}}$ there exists $\tilde{y} = (q_0, q_{1-}, y_2, \dots)$ satisfying $\pi_\lambda(\tilde{y}) = \pi_\lambda(\tilde{q})$.

(4) Applying subsequently formula (1) we obtain that if $\tilde{y} = (y_0, y_1, y_2, \dots)$ then the second coordinate of the point $\pi_\lambda(\tilde{y})$ belonging to the unit square is y_0 and its first coordinate is given by

$$s(\tilde{y}, \lambda) = \varphi(y_1) + \lambda(\varphi(y_2) - \frac{1}{2}) + \lambda^2(\varphi(y_3) - \frac{1}{2}) + \dots \tag{25}$$

If $\lambda_1 \in \Gamma_{\tilde{q}}$ and $\tilde{y} = \tilde{y}(\lambda_1) = (q_0, q_{1-}, y_2, \dots)$ is a point satisfying $\pi_{\lambda_1}(\tilde{y}) = \pi_{\lambda_1}(\tilde{q}, \lambda_1)$ then $s(\tilde{y}, \lambda_1) = s(\tilde{q}, \lambda_1)$ but the series $s(\tilde{y}, \lambda)$ and $s(\tilde{q}, \lambda)$ do not coincide identically because their constant terms differ:

$$\varphi(y_1) = \varphi(q_{1-}) = q_{1-} < \frac{1}{2}, \quad \text{and} \quad \varphi(q_{1+}) = q_{1+} > \frac{1}{2}$$

(the map F identifies points lying on $\xi_{q_{1-}}$ and on $\xi_{q_{1+}}$ only if q_{1+}, q_{1-} are sufficiently close to $\frac{1}{2}$ and $\varphi(y) \equiv y$ for such points).

Let us denote by $s_{\lambda_1}(\tilde{y}, \delta\lambda)$ the series $s(\tilde{y}, \lambda)$ expressed as a power series in $(\lambda - \lambda_1) = \delta\lambda$. Then

$$s_{\lambda_1}(\tilde{y}, \delta\lambda) = s_0(\tilde{y}, \lambda_1) + s_1(\tilde{y}, \lambda_1)\delta\lambda + \dots + s_n(\tilde{y}, \lambda_1)\delta\lambda^n + \dots$$

Since $s_{\lambda_1}(\tilde{y}, \delta\lambda) \neq s_{\lambda_1}(\tilde{q}, \delta\lambda)$ there is an n such that $s_n(\tilde{y}, \lambda_1) \neq s_n(\tilde{q}, \lambda_1)$. For any rational $r > 0$ let us define $\Gamma_{n,r} \subset \Gamma_{\tilde{q}}$ by

$$\Gamma_{n,r} = \{\lambda_1: \text{there exists a } \tilde{y} \text{ satisfying } \pi_{\lambda_1}(\tilde{y}) = \pi_{\lambda_1}(\tilde{q}), s_i(\tilde{y}, \lambda_1) = s_i(\tilde{q}, \lambda_1) \text{ for } i \in [1, n-1], |s_n(\tilde{q}, \lambda_1) - s_n(\tilde{y}, \lambda_1)| \geq r\}.$$

One checks as above that $\Gamma_{n,r}$ are measurable. Thus $\text{mes } \Gamma_{n_1,r_1} > 0$ for some choice of indices. Set $\Gamma = \Gamma_{n_1,r_1}$. Let $\Gamma_1 \in \Gamma$ and let $\tilde{y}_{(1)}$ be the corresponding point satisfying $\pi_{\lambda_1}(\tilde{q}) = \pi_{\lambda_1}(\tilde{y}_{(1)})$. Then $s_{\lambda_1}(\tilde{q}, \delta\lambda), s_{\lambda_1}(\tilde{y}_{(1)}, \delta\lambda)$ may be written as

$$s_{\lambda_1}(\tilde{q}, \delta\lambda) = s_0(\lambda_1) + s_1(\lambda_1)\delta\lambda + \dots + s_{n_1-1}(\lambda_1)\delta\lambda^{n_1-1} + s_{n_1}(\lambda_1)\delta\lambda^{n_1} + \dots \tag{26}$$

$$s_{\lambda_1}(\tilde{y}_{(1)}, \delta\lambda) = s_0(\lambda_1) + s_1(\lambda_1)\delta\lambda + \dots + s_{n_1-1}(\lambda_1)\delta\lambda^{n_1-1} + s_{n_1}(\tilde{y}_{(1)}, \lambda_1)\delta\lambda^{n_1} + \dots$$

Let us define

$$\theta(\delta\lambda) = s_{\lambda_1}(\tilde{q}, \delta\lambda) - s_{\lambda_1}(\tilde{y}_{(1)}, \delta\lambda),$$

$$\theta(\tilde{y}, \delta\lambda) = s_{\lambda_1}(\tilde{y}, \delta\lambda) - s_{\lambda_1}(\tilde{y}_{(1)}, \delta\lambda).$$

Then

$$(\delta\lambda) = \theta_{n_1}(\delta\lambda)^{n_1} + \theta_{n_1+1}(\delta\lambda)^{n_1+1} + \dots$$

$$\theta(\tilde{y}, \delta\lambda) = \theta_0(\tilde{y}) + \theta_1(\tilde{y})\delta\lambda + \dots + \theta_{n_1-1}(\tilde{y})\delta\lambda^{n_1-1} + \theta_{n_1}(\tilde{y})\delta\lambda^{n_1} + \dots \tag{27}$$

where $|\theta_{n_1}| > r_1$.

(5) We fix some $\varepsilon_1, 0 < \varepsilon_1 < r_1/2$. Let

$$d(\tilde{y}, \tilde{y}') = \sum_{i=0}^{\infty} |y_i - y'_i|/2^i$$

be the distance in $[\widetilde{0}, \widetilde{1}]$. We denote by $B_\varepsilon(\tilde{y})$ the ball of radius ε centred in $\tilde{y} \in [\widetilde{0}, \widetilde{1}]$. There exists a $\delta_1 = \delta_1(\varepsilon_1) > 0$ satisfying the following conditions:

- (a) if $\tilde{y} \in B_{\delta_1}(\tilde{y}_{(1)})$ then $|\theta_i(\tilde{y})| < \varepsilon_{1/3}$ for $i \in [0, n_1]$;
- (b) if $|\delta\lambda| < \delta_1$ then for any \tilde{y}

$$\left| \frac{n_1+1}{2!} \theta_{n_1+1}(\tilde{y})\delta\lambda + \frac{(n_1+1)(n_1+2)}{3!} \theta_{n_1+2}(\tilde{y})\delta\lambda^2 + \dots \right| > \frac{\varepsilon_1}{3}.$$

For $\lambda_1 \in \Gamma$ we shall use $V(\lambda_1)$ to denote the non-empty set of $\tilde{y} = (q_0, q_1, \dots, y_2, \dots)$ appearing in the definition of λ_1 , i.e.,

$$\pi_{\lambda_1}(\tilde{y}) = \pi_{\lambda_1}(\tilde{q}),$$

the curves $s(\tilde{q}, \lambda), s(\tilde{y}, \lambda)$ have tangency of order $n_1 - 1$ at $\lambda = \lambda_1$, and

$$|s_n(\tilde{q}, \lambda_1) - s_n(\tilde{y}, \lambda_1)| > r_1.$$

Let

$$\Gamma(\lambda_1, \tilde{y}_{(1)}) = \{ \delta\lambda : |\delta\lambda| < \delta_1, \lambda_1 + \delta\lambda \in \Gamma, V(\lambda_1 + \delta\lambda) \cap B_{\delta_1}(\tilde{y}_{(1)}) \neq \emptyset \}.$$

Differentiating (27) $n_1 - 1$ times we obtain for

$$\begin{aligned} &\delta\lambda \in \Gamma(\lambda_1, \tilde{y}_{(1)}) \text{ and } \tilde{y} \in V(\lambda_1 + \delta\lambda) \cap B_{\delta_1}(\tilde{y}_{(1)}) \\ &(n_1!) \delta\lambda [(\theta_{n_1} - \theta_{n_1}(\tilde{y})) + \frac{n_1+1}{2!} (\theta_{n_1+1} - \theta_{n_1+1}(\tilde{y}))\delta\lambda \\ &\quad + \frac{(n_1+1)(n_1+2)}{3!} (\theta_{n_1+2} - \theta_{n_1+2}(\tilde{y}))\delta\lambda^2 + \dots] = (n_1 - 1)! \theta_{n_1-1}(\tilde{y}). \end{aligned} \tag{28}$$

Taking into account the choice of ε_1 and δ_1 we obtain from (28)

$$\delta\lambda(\tilde{y}) = \frac{1}{n_1} \theta_{n_1-1}(\tilde{y}) \frac{1}{\theta_{n_1}(1 + \varepsilon(\tilde{y}))} \tag{29}$$

while $|\varepsilon(\tilde{y})| > \frac{1}{2}$.

(6) Let

$$A_k(\lambda_i) = \left\{ \frac{d^k}{d\lambda^k} s(\tilde{y}, \lambda) \Big|_{\lambda=\lambda_1} \right\}.$$

One easily checks that for $\lambda_1 < \frac{1}{2} \text{mes } A_k(\lambda_1) = 0$ for any k (notice that $A_0(\lambda_1) = \{s(\tilde{y}, \lambda_1)\} = \Lambda \cap \xi_{a_0}$). When \tilde{y} varies in $[\widetilde{0}, \widetilde{1}]$ the values of $\theta_{n_1-1}(\tilde{y})$ belong to the set $A_{n_1-1}(\lambda_1) + c$ where

$$c = \frac{d^{n_1-1}}{d\lambda^{n_1-1}} s(\tilde{y}_{(1)}, \lambda) \Big|_{\lambda=\lambda_1}.$$

Using (29) we obtain from $\text{mes } A_{n_1-1}(\lambda_1) = 0$ that $\text{mes } \Gamma(\lambda_1, \tilde{y}_{(1)}) = 0$.

Consider a cover $\{l_i\}_{i=1}^{m_1}$ of $[0, \lambda_0]$ by intervals of $\text{diam} < \delta_1$, and a cover $\{B_{\delta_{ij}}\}_{i=1}^{m_1}$ of the set $\tilde{y} = (q_0, q_{1-}, y_2, \dots)$. Let

$$\Gamma_{ij} \subset \Gamma = \{\lambda_1 \in l_i : V(\lambda_1) \cap B_j \neq \emptyset\}.$$

The above arguments show that $\text{mes } \Gamma_{ij} = 0$ for any choice of i, j . Hence $\text{mes } \Gamma = 0$, which proves lemma 4. \square

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