# Weierstraß-Institut für Angewandte Analysis und Stochastik 

im Forschungsverbund Berlin e.V.

# On Andronov-Hopf bifurcations of two-dimensional diffeomorphisms with homoclinic tangencies. 

Sergey V. Gonchenko ${ }^{1}$, Vladimir S. Gonchenko ${ }^{1}$ submitted: February 28, 2000

1 Institute for Applied Mathematics and Cybernetics, 10 Ulyanova St.
Nizhny Novgorod, 603005
Russia
E-Mail: gonchenko@focus.nnov.ru

Preprint No. 556
Berlin 2000


1991 Mathematics Subject Classification. 58F12, 58F13.
Key words and phrases. homoclinic tangency, invariant curve, Andronov-Hopf bifurcation, strange attractors, Newhouse regions.

Edited by
Weierstraß-Institut für Angewandte Analysis und Stochastik (WIAS)
Mohrenstraße 39
D-10117 Berlin
Germany

Fax: $\quad+49302044975$
E-Mail (X.400): $\quad c=d e ; a=d 400-\mathrm{gw} ; \mathrm{p}=\mathrm{WIAS}-$ BERLIN; $\mathrm{s}=$ preprint
E-Mail (Internet): preprint@wias-berlin.de
World Wide Web: http://www.wias-berlin.de/


#### Abstract

The bifurcation of the birth of a closed invariant curve in the two-parameter unfolding of a two-dimensional diffeomorphism with a homoclinic tangency of invariant manifolds of a hyperbolic fixed point of neutral type (i.e. such that the Jacobian at the fixed point equals to 1 ) is studied. The existence of periodic orbits with multipliers $e^{ \pm i \psi}(0<\psi<\pi)$ is proved and the first Lyapunov value is computed. It is shown that, generically, the first Lyapunov value is non-zero and its sign coincides with the sign of some separatrix value (i.e. a function of coefficients of the return map near the global piece of the homoclinic orbit).


## Introduction.

Homoclinic orbits are one of the most interesting object of study in the theory of dynamical systems, because their presence leads to nontrivial dynamics. Recall that the Poincaré homoclinic orbit is an orbit which is biasymptotic to a saddle periodic orbit, i.e. it is an orbit lying in the intersection of the stable and unstable invariant manifolds of the saddle. If this intersection is transverse, the homoclinic orbit is called rough; otherwise, it is called an orbit of homoclinic tangency.

It is well-known that the set of all orbits lying entirely in a small neighborhood of a rough homoclinic orbit is hyperbolic and has a nontrivial structure which admits a complete description in terms of symbolic dynamics [1]. The situation is drastically different in the case of homoclinic tangency. Here, the complete study is proven $[2,3,4]$ to be impossible in any finite-parameter unfolding. However, some main bifurcations of periodic orbits have been studied sufficiently well $[5,6,7,4]$. It occurred that the basic feature of the bifurcation of homoclinic tangency is the appearance of a large number (even infinitely many) coexisting periodic orbits of different topological types. This is closely connected with the so-called Newhouse phenomenon: systems with homoclinic tangencies are dense in open regions (the Newhouse regions) in the space of smooth dynamical systems [8, 9, 10, 11, 12]. We note that it is the Newhouse regions to which, presumably, the most of known systems with chaotic behavior belong, e.g. systems with quasistochastic and wildhyperbolic strange attractors [13, 14, 15, 4, 16, 17]. Therefore, the question on which type periodic orbits (and more complicated invariant sets) can appear via homoclinic bifurcations is especially important. For general (codimension one) homoclinic bifurcations this question was solved (including the multidimensional case) in $[7,4,18]$ where necessary and sufficient conditions for the birth of periodic orbits
of a given topological type were obtained and the appearance of invariant tori and even infinitely many coexisting strange attractors was established in some situations.
In the two-dimensional case, it is known since [5] that, generically, the bifurcation of homoclinic tangency produces, along with saddle periodic orbits, either stable or completely unstable ones. The type of stability depends on wether the saddle value $\sigma$ (i.e. the absolute value of the product of the multipliers of the saddle fixed point) is greater or less than 1. As a result, in the Newhouse regions close to two-dimensional diffeomorphisms with a homoclinic tangency systems with infinitely many coexisting saddle and stable (if $\sigma<1$ ) or saddle and completely unstable ( $\sigma>1$ ) periodic orbits are dense [19]. It is also known [2, 3, 20] that diffeomorphisms with infinitely many arbitrarily degenerate periodic orbits are dense in these regions. Note that these degenerate periodic orbits have exactly one multiplier equal to +1 or -1 with an arbitrarily large number of Lyapunov values (may be all of them) vanishing ${ }^{1}$. When $\sigma \neq 1$ there cannot be other degeneracies $[7,4,18]$ and, in particular, no close diffeomorphism can have closed invariant curves or periodic orbits with the multipliers $e^{ \pm i \varphi}$.

In the present paper we show that if $\sigma=1$ at the moment of homoclinic tangency, then, along with such usual for systems with homoclinic tangencies bifurcations as a saddle-node and a period-doubling, the Andronov-Hopf bifurcations connected with the birth of closed invariant curves from periodic orbits with the multipliers $e^{ \pm i \psi}$ take place when the homoclinic tangency unfolds.
Note that if both multipliers of the saddle are positive, two different cases are possible depending on the sign of the separatrix value $R$ defined below (formula (5)). Namely, if $R<0$, we show that asymptotically stable invariant curves are born near the homoclinic tangency, whereas at $R>0$ we show the birth of unstable invariant curves. In the case of negative multipliers both stable and unstable invariant curves are born at $R \neq 0$.

Let us proceed to detailed formulation of the results (theorems A and B below). Let $f_{0}$ be a two-dimensional orientation-preserving $C^{r}$-diffeomorphism ( $r \geq 4$ ) satisfying the following conditions.
A) $f_{0}$ has a saddle fixed point $O$ with the multipliers $\lambda, \gamma$ such that $|\lambda|<1,|\gamma|>1$;
B) the saddle value $\sigma=\lambda \gamma$ equals to 1 (i.e. $O$ is a point of neutral type);
C) the stable and unstable manifolds $W_{0}^{s}$ and $W_{0}^{u}$ of the saddle $O$ have a quadratic homoclinic tangency at the points of some homoclinic orbit $\Gamma_{0}$ (Fig.1a).

[^0]

Figure 1

The diffeomorphisms $C^{r}$-close to $f_{0}$ and satisfying the same conditions A)-C) compose a bifurcational surface $H \in$ Diffr $^{r}$ of codimension two. Let $f_{\mu}, \mu \equiv\left(\mu_{1}, \mu_{2}\right)$, be a two-parameter family passing through $f_{0}$ at $\mu_{1}=\mu_{2}=0$. Assume that
D) the family $f_{\mu}$ is $C^{r}$-smooth with respect to all variables and $\mu$ and it is transverse to the bifurcational surface $H$ at $\mu=0$.

It is not hard to understand which should be the nature of the governing parameters $\mu$. One of them, say $\mu_{1}$, must control the position of the invariant manifolds of $O$ near the points of homoclinic tangency (i.e. we choose $\mu_{1}$ as the splitting parameter near some homoclinic point). The second parameter $\mu_{2}$ must control the saddle value at $O$, i.e.

$$
\sigma(\mu)=1+\mu_{2}
$$

Let us take a sufficiently small neighborhood $U=U\left(O \cup \Gamma_{0}\right)$ of the closure of the homoclinic orbit $\Gamma_{0}$. It is the union of a small disc $U_{0}$ around $O$ and a finite number of small neighborhoods of those points of $\Gamma_{0}$ which do not lie in $U_{0}$ (Fig.1b). We will call as a p-round periodic orbit a periodic orbit of $f_{\mu}$ which lies entirely in $U$ and visits every component of $U \backslash U_{0}$ exactly $p$ times on the period (i.e. every such orbit, before it returns to its initial point, runs $U$ exactly $p$ times).
In the present paper we study the bifurcations of single-round $(p=1)$ periodic orbits in $U$ (for families $f_{\mu}$ satisfying A$\left.)-\mathrm{D}\right)$ ). Such orbits correspond to fixed points of the first return map $T_{1} T_{0}^{k}(k=\bar{k}, \bar{k}+1, \ldots$ for some sufficiently large $\bar{k})$ near the homoclinic orbit. Here $T_{0}$ is the so-called local map which is the diffeomorphism $f_{\mu}$ restricted to the neighborhood $U_{0}$ of its fixed point $O_{\mu}$. The global map $T_{1}$ is some power of $f_{\mu}$ acting from a small neighborhood of some point of $\Gamma_{0}$ in $W_{l o c}^{u} \cap U_{0}$ into a small neighborhood of some point of $\Gamma_{0}$ in $W_{l o c}^{s} \cap U_{0}$.
The map $T_{0}$ is $C^{r}$-smooth and it has a saddle fixed point $O_{\mu}$ at all small $\mu$ (at $\mu=0$
it is the original point $O$ ). It is shown in $[21,22]$ that in some $C^{r-1}$-coordinates $(x, y)$ in $U_{0}$ the map $T_{0}$ has the following form for all small $\mu$ :

$$
\begin{equation*}
\bar{x}=\lambda(\mu) x+h(x, y, \mu) x^{2} y \quad, \quad \bar{y}=\gamma(\mu) y+g(x, y, \mu) x y^{2} \tag{1}
\end{equation*}
$$

where $h(x, y, \mu) x y \in C^{r-1}, g(x, y, \mu) x y \in C^{r-1}$. Hereafter, we assume that the map $T_{0}$ is brought to this form.
In these coordinates, the fixed point $O_{\mu}$ is in the origin and $W_{l o c}^{s}(\mu) \cap U_{0}$ and $W_{l o c}^{u}(\mu) \cap$ $U_{0}$ are segments of straight lines $y=0$ and $x=0$, respectively. By assumption C), at $\mu=0$, the homoclinic orbit $\Gamma_{0}$ have points both in $W_{l o c}^{u} \cap U_{0}$ and in $W_{\text {loc }}^{s} \cap U_{0}$. Chose a pair of such points $M^{+}\left(x^{+}, 0\right) \in W_{\text {loc }}^{s} \cap U_{0}$ and $M^{-}\left(0, y^{-}\right) \in W_{\text {loc }}^{u} \cap U_{0}$; without loss of generality we may assume $x^{+}>0, y^{-}>0$. Let $\Pi_{0}$ and $\Pi_{1}$ be sufficiently small neighborhoods of $M^{+}$and $M^{-}$, respectively. We denote the coordinates in $\Pi_{0}$ and $\Pi_{1}$ as, respectively, $\left(x_{0}, y_{0}\right)$ and $\left(x_{1}, y_{1}\right)$.
By construction, there exists such positive $n_{0}$ that $M^{+}=f_{0}^{n_{0}}\left(M^{-}\right)$. We define the global map $T_{1}$ as follows:

$$
\begin{equation*}
T_{1} \equiv f_{\mu}^{n_{0}}: \Pi_{1} \rightarrow \Pi_{0} \tag{2}
\end{equation*}
$$

It is defined at all small $\mu$ and it is, at least, $C^{r-1}$-smooth (in the coordinates for which $T_{0}$ has the form (1)).
Let us write $T_{1}$ in the following form:

$$
\begin{equation*}
\bar{x}_{0}-x^{+}=F\left(x_{1}, y_{1}-y^{-}, \mu\right), \quad \bar{y}_{0}=G\left(x_{1}, y_{1}-y^{-}, \mu\right) \tag{3}
\end{equation*}
$$

where $F(0,0,0)=0, G(0,0,0)=0$. According to the condition C), the image of the segment $\left\{x_{1}=0\right\}$ by $T_{1}$ must have a quadratic tangency with $\left\{y_{0}=0\right\}$ at $\mu=0$. Hence,

$$
\frac{\partial G(0,0,0)}{\partial y_{1}}=0 \quad, \quad \frac{\partial^{2} G(0,0,0)}{\partial y_{1}^{2}}=2 d \neq 0
$$

Thus, one can write

$$
\begin{align*}
& F \equiv a x_{1}+b\left(y_{1}-y^{-}\right)+e_{20} x_{1}^{2}+e_{11} x_{1}\left(y_{1}-y^{-}\right)+e_{02}\left(y_{1}-y^{-}\right)^{2}+ \\
& O\left[\left(\left|x_{1}\right|+\left|y_{1}-y^{-}\right|\right)^{3}\right], \\
& G \equiv \mu_{1}+c x_{1}+d\left(y_{1}-y^{-}\right)^{2}+f_{20} x_{1}^{2}+f_{11} x_{1}\left(y_{1}-y^{-}\right)+  \tag{4}\\
& +f_{03}\left(y_{1}-y^{-}\right)^{3}+f_{30} x_{1}^{3}+f_{21} x_{1}^{2}\left(y_{1}-y^{-}\right)+f_{12} x_{1}\left(y_{1}-y^{-}\right)^{2}+ \\
& +o\left[\left(\left|x_{1}\right|+\left|y_{1}-y^{-}\right|\right)^{3}\right],
\end{align*}
$$

where the coefficients $a, b, \ldots, f_{03}$ (as well as $x^{+}$and $y^{-}$) are some functions of $\mu$. Since $T_{1}$ preserves orientation,

$$
b c<0
$$

Note that we choose the parameter $\mu_{1}$ such that it enters (in the main order) the right-hand side of the equation for $\bar{y}_{0}$ additively. It means that $\mu_{1}$ is the splitting parameter for the invariant manifolds of $O$ near the homoclinic point $M^{+}$.
Let us introduce the separatrix value

$$
\begin{equation*}
R \equiv 2 a-\frac{b}{d} f_{11}-2 \frac{c}{d} e_{02} \tag{5}
\end{equation*}
$$

where the coefficients of the global map $T_{1}$ (see (4)) are taken at $\mu=0$. Note that it is important in this definition that the coordinates are chosen such that the local $\operatorname{map} T_{0}$ has the form (1)).

Theorem A. If $R \neq 0$, in the parameter plane $\left(\mu_{1}, \mu_{2}\right)$ there exists a sequence of open regions $\Delta_{k}$, accumulating at $\mu=0$ as $k \rightarrow+\infty$, such that the diffeomorphism $f_{\mu}$ has a closed invariant curve at all $\mu \in \Delta_{k}$. At $\lambda>0, \gamma>0$, the invariant curves are asymptotically stable at $R<0$ and unstable at $R>0$. If $\lambda<0, \gamma<0$, then the invariant curves at $\mu \in \Delta_{k}$ are stable or unstable depending on parity of $k$.

We prove this theorem by means of the study of bifurcations of single-round periodic orbits or, what is the same, of bifurcations of fixed points of the maps $T_{k}=T_{1} T_{0}^{k}$ at all sufficiently large $k: k=\bar{k}, \bar{k}+1, \ldots$. By definition, a single-round periodic orbit has exactly one point in each of the neighborhoods $\Pi_{0}$ and $\Pi_{1}$. Let $M_{0} \in \Pi_{0}$ and $M_{1} \in \Pi_{1}$ be such points. Then, $M_{0}=T_{1}\left(M_{1}\right)$ and there exists such an integer $k$ that $M_{1}=T_{0}^{k}\left(M_{0}\right)$. Thus, the point $M_{0} \in \Pi_{0}$ is a fixed point of $T_{k} \equiv T_{1} T_{0}^{k}$ (the period of the corresponding orbit of $f_{\mu}$ equals $k+n_{0}$, see (2)).

Theorem B. 1. In the plane of parameters $\left(\mu_{1}, \mu_{2}\right)$, for every sufficiently large $k$ there exist bifurcational curves $L_{k}^{+}, L_{k}^{-}$and $L_{k}^{\varphi}$, corresponding to single-round periodic orbits (fixed points of $T_{k}$ ) with multipliers $+1,-1$ and $e^{ \pm i \psi}(0<\psi<\pi)$, respectively. The curves $L_{k}^{+}$and $L_{k}^{-}$accumulate to the line $\mu_{1}=0$ as $k \rightarrow+\infty$. The curves $L_{k}^{\varphi}$ connect points $B_{k}^{++}$and $B_{k}^{--}$on, respectively, $L_{k}^{+}$and $L_{k}^{-}$and accumulate at the point $\mu_{1}=\mu_{2}=0$.
2. At $\mu$ from the region $D_{k}$ between the curves $L_{k}^{+}$and $L_{k}^{-}$the diffeomorphism $f_{\mu}$ has two single-round periodic orbits one of which, $Q_{k}$, is saddle and the other, $P_{k}$ is asymptotically stable at $\mu \in D_{k}^{s}$ and completely unstable at $\mu \in D_{k}^{u}$ where $D_{k}^{s}$ $\left(D_{k}^{u}\right)$ is the region in $D_{k}$ to the left (resp., to the right) of $L_{k}^{\varphi}$. The transitions into the region $D_{k}$ across the curves $L_{k}^{+}$(without the point $B_{k}^{++}=L_{k}^{+} \cap L_{k}^{\varphi}$ ) and $L_{k}^{-}$ (without the point $B_{k}^{--}=L_{k}^{-} \cap L_{k}^{\varphi}$ ) correspond, respectively, to generic saddle-node and period-doubling bifurcations of $P_{k}$ (on $L_{k}^{+}$the orbits $P_{k}$ and $Q_{k}$ merge together). At $\mu=B_{k}^{++}$the orbit $P_{k}$ has two multipliers equal to 1 , and both the multipliers are equal to -1 at $\mu=B_{k}^{--}$.
3. If $R \lambda^{k}<0$, the boundary $L_{k}^{\varphi}$ of stability of $P_{k}$ is "safe": the first Lyapunov value is negative, so at the transition across $L_{k}^{\varphi}$ (except for two points for which $\psi=\pi / 2,2 \pi / 3)$ towards the increase of $\mu_{2}$ the orbit $P_{k}$ becomes unstable and $a$ stable invariant curve is born from it.
If $R \lambda^{k}>0$, the boundary $L_{k}^{\varphi}$ is "dangerous": the first Lyapunov value is positive, so at the transition across $L_{k}^{\varphi}$ (except for two points for which $\psi=\pi / 2,2 \pi / 3$ ) towards the decrease of $\mu_{2}$ an unstable invariant curve is born from $P_{k}$.

See figure 2 as an illustration to theorem B.


Figure 2
An illustration to theorem B

It is obvious that theorem A follows from theorem B. Here, the region $\Delta_{k}$ is some part of $D_{k}$ adjoining to $L_{k}^{e}$ from the left if $R \lambda^{k}>0$ and from the right if $R \lambda^{k}<0$, corresponding to the existence of the invariant curve.
The content of the paper is as follows. In section 2 we study properties of iterations of the local map $T_{0}$. In section 3 the first return maps $T_{k}$ are constructed. Here we prove that the map $T_{k}$ for sufficiently large $k$ may be brought to a certain form close to the quadratic Henon family. Unlike some standard results [23, 2, 4], we take into account also small terms of the order $O\left(\lambda^{k}\right)$. In section 4 bifurcations of fixed points of $T_{k}$ are studied, and the first Lyapunov value is calculated for the fixed point undergoing the Andronov-Hopf bifurcation. In section 5 theorems A and B are finally proved.

## 1 Properties of the local map $T_{0}(\mu)$.

The map $T_{0}$ is defined as the restriction of diffeomorphism $f_{\mu}$ of the neighbourhood $U_{0}$, i.e., $T_{0}(\mu) \equiv f_{\mu \mid U_{0}}$. The map $T_{0}$ has, at all sufficiently small $\mu$, the fixed saddle point $O_{\mu}$. It is well-known that in some $C^{r}$-coordinates $(x, y)$ on $U_{0}$ the map $T_{0}$ can be written as

$$
\begin{equation*}
\bar{x}=\lambda(\mu) x+h(x, y, \mu), \bar{y}=\gamma(\mu) y+g(x, y, \mu) \tag{6}
\end{equation*}
$$

where $h(0, y, \mu)=g(x, 0, \mu) \equiv 0$. Here, the axes $x$ and $y$ are eigendirections of the Jacobi matrix of $T_{0}$ at $O_{\mu}$, and the local stable and unstable manifolds of $O_{\mu}$ are straightened.

Note that the form (6) of $T_{0}$ is not very convenient from the technical point of view since the right-hand sides of (6) contain too many non-resonant terms. For example, if to write functions $h$ and $g$ in the following extended form

$$
\begin{aligned}
h(x, y, \mu) & \equiv \varphi_{1}(x, \mu) x+\varphi_{2}(y, \mu) x+\tilde{h}(x, y, \mu) x^{2} y \\
g(x, y, \mu) & \equiv \phi_{1}(y, \mu) y+\phi_{2}(x, \mu) y+\tilde{g}(x, y, \mu) x y^{2}
\end{aligned}
$$

where $\varphi_{\alpha}(0, \mu) \equiv 0, \phi_{\alpha}(0, \mu) \equiv 0, \alpha=1,2$, then one can see that functions $\varphi$ and $\phi$ contain only nonresonant monomials ${ }^{2}$. It was shown in [21, 22] (see also [24] for the general multidimensional case) that such "always nonresonant monomials" can all be nullified by a sufficiently smooth change of variables. Namely, the following result is valid [21, 22] :

There exist such $\delta_{1}>0$ and $\delta_{2}>0$ that, at $\|(x, y)\| \leq \delta_{1}$ and $\|\mu\| \leq \delta_{2}$, the map $T_{0}(\mu)$ is brought to the form

$$
\begin{equation*}
\bar{x}=\lambda(\mu) x+h_{1}(x, y, \mu) x^{2} y \quad, \quad \bar{y}=\gamma(\mu) y+g_{1}(x, y, \mu) x y^{2} \tag{7}
\end{equation*}
$$

where $h_{1}(x, y, \mu) x y \in C^{r-1}, g_{1}(x, y, \mu) x y \in C^{r-1}$, by means of a $C^{r-1}$-smooth transformation of coordinates (this transformation is $C^{r-2}$ with respect to parameters).

By (7), the point $O_{\mu}$ is in the origin (at all sufficiently small $\mu$ ) and $W_{l o c}^{s}(\mu) \cap U_{0}$ and $W_{l o c}^{u}(\mu) \cap U_{0}$ have, respectively, equations $y=0$ and $x=0$. At $\mu=0$ let us choose in $U_{0}$ a pair of points of the orbit $\Gamma_{0}: M^{+}\left(x^{+}, 0\right)$ and $M^{-}\left(0, y^{-}\right)$, and take their sufficiently small rectangular neighbourhoods $\Pi_{0}$ and $\Pi_{1}$. We denote the coordinates $(x, y)$ in $\Pi_{0}$ and $\Pi_{1}$ as $\left(x_{0}, y_{0}\right)$ and $\left(x_{1}, y_{1}\right)$, respectively. Without loss of generality, we assume that $x^{+}>0, y^{-}>0$. The neighborhoods $\Pi_{0}$ and $\Pi_{1}$ are defined as follows

$$
\begin{align*}
& \Pi_{0}=\left\{\left(x_{0}, y_{0}\right)| | x_{0}-x^{+}\left|\leq \varepsilon_{0},\left|y_{0}\right| \leq \varepsilon_{0}\right\}\right. \\
& \Pi_{1}=\left\{\left(x_{1}, y_{1}\right)| | x_{1}\left|\leq \varepsilon_{1},\left|y_{1}-y^{-}\right| \leq \varepsilon_{1}\right\}\right. \tag{8}
\end{align*}
$$

where $\varepsilon_{0}$ and $\varepsilon_{1}$ are sufficiently small (so $T_{0}\left(\Pi_{0}\right) \cap \Pi_{0}=\emptyset, T_{0}^{-1}\left(\Pi_{1}\right) \cap \Pi_{1}=\emptyset$, in particular).
To study the maps $T_{k} \equiv T_{1} T_{0}^{k}$ it is necessary, first of all, to have appropriate formulas and estimates for the maps $T_{0}^{k}: \Pi_{0} \rightarrow \Pi_{1}$ for all sufficiently large $k$. To this aim, the form (7) of $T_{0}(\mu)$ is very convenient because the iterations of the map $T_{0}$ in form (7) are asymptotically close (as $k \rightarrow \infty$ ) to those in the linear case ${ }^{3}$. Namely, the following lemma holds.

[^1]Lemma 1 For any $\varepsilon>0$ there exists $\bar{k}>0$ such that for any $k \geq \bar{k}$ and $\|\mu\| \leq \varepsilon$ the map $T_{0}^{k}: \Pi_{0} \rightarrow \Pi_{1}$ can be represented as follows:

$$
\begin{align*}
& x_{1}=\lambda^{k}(\mu) x_{0}\left(1+k \hat{\gamma}^{-k} \hat{\xi}_{k}\left(x_{0}, y_{1}, \mu\right)\right) \\
& y_{0}=\gamma^{-k}(\mu) y_{1}\left(1+k \hat{\gamma}^{-k} \hat{\eta}_{k}\left(x_{0}, y_{1}, \mu\right)\right) \tag{9}
\end{align*}
$$

where $\hat{\gamma}=\gamma(0) /(1+\varepsilon)$ and functions $\hat{\xi}_{k} \equiv \tilde{\xi}_{k} \cdot x_{0} y_{1}$ and $\hat{\eta}_{k} \equiv \tilde{\eta}_{k} \cdot x_{0} y_{1}$ and their derivatives (along with derivatives with respect to $\mu$ ) up to the order $(r-2)$ are bounded, uniformly in $k$. The derivatives of order $(r-1)$ from the right-hand sides of (9) tend to zero as $k \rightarrow \infty$.

Proof. The proof of this lemma repeats closely the proof of an analogous statement (lemma 1.2 in [22]). Therefore, we prove here only the boundedness for the functions $\xi_{k}$ and $\eta_{k}$ themselves; the boundedness of derivatives is verified along the same lines (for more detail see [22, 24]).
We will use the method of the boundary-value problem [1, 25] in a modification of [26]. For the sake of simplicity, we write the map $T_{0}(\mu)$ in the form

$$
\begin{equation*}
\bar{x}=\lambda(\mu) x+\hat{h}(x, y, \mu), \quad \bar{y}=\gamma(\mu) y+\hat{g}(x, y, \mu) \tag{10}
\end{equation*}
$$

where $\lambda(0) \gamma(0)=1$ and, by (7), $\hat{h}(x, y, \mu) \equiv x^{2} y\left(\beta_{1}(\mu)+O(|x|+|y|)\right), \hat{g}(x, y, \mu) \equiv$ $x y^{2}\left(\beta_{2}(\mu)+O(|x|+|y|)\right)$. Also, in the proof, we will use notations $\lambda(\mu) \equiv$ $\lambda, \gamma(\mu) \equiv \gamma$.
Let us consider the following operator $\Phi:\left[\left(x_{j}, y_{j}\right)\right]_{j=0}^{k} \mapsto\left[\left(\bar{x}_{j}, \bar{y}_{j}\right)\right]_{j=0}^{k}$ where

$$
\begin{align*}
& \bar{x}_{j}=\lambda^{j} x_{0}+\sum_{s=0}^{j-1} \lambda^{j-s-1} \hat{h}\left(x_{s}, y_{s}, \mu\right)  \tag{11}\\
& \bar{y}_{j}=\gamma^{j-k} y_{k}-\sum_{s=j}^{k-1} \gamma^{j-s-1} \hat{g}\left(x_{s}, y_{s}, \mu\right)
\end{align*}
$$

$j=0,1, \ldots, k$. The operator $\Phi$ is defined on the set

$$
Z(\delta)=\left\{z=\left[\left(x_{j}, y_{j}\right)\right]_{j=0}^{k},\|z\| \leq \delta\right\}
$$

where the norm $\|\cdot\|$ is defined as the maximum of the absolute values of the components $x_{j}, y_{j}$ of the vector $z$. Note that if $z_{0}=\left[\left(x_{j}^{0}, y_{j}^{0}\right)\right]_{j=0}^{k}$ is a fixed point of $\Phi$, then

$$
\left(x_{0}^{0}, y_{0}^{0}\right) \xrightarrow{T_{0}}\left(x_{1}^{0}, y_{1}^{0}\right) \xrightarrow{T_{0}} \ldots \xrightarrow{T_{0}}\left(x_{k}^{0}, y_{k}^{0}\right),
$$

i.e. the fixed point of $\Phi$ is the orbit of the map $T_{0}$.

For sufficiently small $\delta=\delta_{0}$ and $\left\|x_{0}\right\| \leq \delta_{0} / 2,\left|y_{k}\right| \leq \delta_{0} / 2$ the operator $\Phi$ maps the set $Z\left(\delta_{0}\right)$ into itself and is contracting on this set (see the proof in [26]). Thus, map (11) has a unique fixed point $z_{0}=\left[\left(x_{j}^{0}\left(x_{0}, y_{k}\right), y_{j}^{0}\left(x_{0}, y_{k}\right)\right]_{j=0}^{k}\right.$ which is the limit of iterations by $\Phi$ of any initial sequence in $Z\left(\delta_{0}\right)$, i.e., the coordinates $x_{j}^{0}$ and $y_{j}^{0}$ can
be found by the successive approximations method. As an initial approximation we take the following (the solution of the linear problem)

$$
x_{j}^{0(1)}=\lambda^{j} x_{0} \quad, \quad y_{j}^{0(1)}=\gamma^{j-k} y_{k}
$$

As it follows from (11), the second approximation has the form

$$
\begin{align*}
& x_{j}^{0(2)}=\lambda^{j} x_{0}+\sum_{s=0}^{j-1} \lambda^{j-s-1} \lambda^{2 s} \gamma^{s-k} x_{0}^{2} y_{k}\left(\beta_{1}(\mu)+O\left(|\lambda|^{s}\left|x_{0}\right|+|\gamma|^{s-k}\left|y_{k}\right|\right)\right)= \\
& =\lambda^{j} x_{0}+\lambda^{j} \gamma^{-k} \sum_{s=0}^{j-1} \lambda^{-1} \lambda^{s} \gamma^{s} x_{0}^{2} y_{k}\left(\beta_{1}(\mu)+O\left(|\lambda|^{s}\left|x_{0}\right|+|\gamma|^{s-k}\left|y_{k}\right|\right)\right)= \\
& =\lambda^{j} x_{0}+\lambda^{j} \gamma^{-k} \lambda^{-1} \sum_{s=0}^{j-1}\left(1+\mu_{2}\right)^{s} x_{0}^{2} y_{k}\left(\beta_{1}(\mu)+O\left(|\lambda|^{s}\left|x_{0}\right|+|\gamma|^{s-k}\left|y_{k}\right|\right)\right), \\
& y_{j}^{0(2)}=\gamma^{j-k} y_{k}+\sum_{s=j}^{k-1} \gamma^{j-s-1} \lambda^{s} \gamma^{2(s-k)} x_{0} y_{k}^{2}\left(\beta_{2}(\mu)+O\left(|\lambda|^{s}\left|x_{0}\right|+|\gamma|^{s-k}\left|y_{k}\right|\right)\right)= \\
& =\gamma^{j-k} y_{k}+\gamma^{j-2 k} \sum_{s=j}^{k-1} \gamma^{-1} \lambda^{s} \gamma^{s} x_{0} y_{k}^{2}\left(\beta_{2}(\mu)+O\left(|\lambda|^{s}\left|x_{0}\right|+|\gamma|^{s-k}\left|y_{k}\right|\right)\right)= \\
& =\gamma^{j-k} y_{k}+\gamma^{j-2 k-1} \sum_{s=j}^{k-1}\left(1+\mu_{2}\right)^{s} x_{0} y_{k}^{2}\left(\beta_{2}(\mu)+O\left(|\lambda|^{s}\left|x_{0}\right|+|\gamma|^{s-k}\left|y_{k}\right|\right)\right) \tag{12}
\end{align*}
$$

(in the last step, in both formulas here, we use the relation $\lambda \gamma=1+\mu_{2}$ ). Let $\left|\mu_{2}\right| \leq \varepsilon$. Then

$$
\sum_{s=0}^{j-1}\left(1+\mu_{2}\right)^{s} \leq \sum_{s=0}^{j-1}(1+\varepsilon)^{j-1}=j(1+\varepsilon)^{j-1}
$$

and, analogously,

$$
\sum_{s=j}^{k-1}\left(1+\mu_{2}\right)^{s} \leq(k-j)(1+\varepsilon)^{k-1} .
$$

It follows from (12) that

$$
\begin{align*}
& \left|x_{j}^{0(2)}-\lambda(\mu)^{j} x_{0}\right| \leq L_{1} j \lambda^{j} \gamma^{-k}(1+\varepsilon)^{j}  \tag{13}\\
& \left|y_{j}^{0(2)}-\gamma(\mu)^{j-k} y_{k}\right| \leq L_{2}(k-j) \gamma^{j-2 k}(1+\varepsilon)^{k}
\end{align*}
$$

where $L_{1}$ and $L_{2}$ are some positive constants independent of $j, k$ and $\mu$. When we plug (13) in (11) as the initial guess for the next approximation, it is not hard to see that inequalities (13) will be valid for the next approximation too, with the same constants $L_{1}$ and $L_{2}$, etc.. Thus, we obtain the following formula for the coordinates $x_{j}^{0}$ and $y_{j}^{0}$ of the fixed point of $\Phi$

$$
\begin{align*}
& x_{j}^{0}=\lambda(\mu)^{j} x_{0}\left(1+j \gamma(\mu)^{-k}(1+\varepsilon)^{j}\right) F_{j k}\left(x_{0}, y_{k}, \mu\right),  \tag{14}\\
& y_{j}^{0}=\gamma(\mu)^{j-k} y_{k}\left(1+(k-j) \gamma(\mu)^{-k}(1+\varepsilon)^{k}\right) G_{j k}\left(x_{0}, y_{k}, \mu\right),
\end{align*}
$$

where functions $F_{j k}$ and $G_{j k}$ are uniformly bounded in $j$ and $k$. Assuming $j=k$ for $x_{j}^{0}$ and $j=0$ for $y_{j}^{0}$ in (14) and taking into account $\gamma(\mu)=\gamma(0)+O \mu$, we obtain formula (9).


Figure 3

Hereafter, we will write (9) as

$$
\begin{align*}
& x_{1}=\lambda^{k}(\mu) x_{0}\left(1+\tilde{\gamma}^{-k} \xi_{k}\left(x_{0}, y_{1}, \mu\right)\right), \\
& y_{0}=\gamma^{-k}(\mu) y_{1}\left(1+\tilde{\gamma}^{-k} \eta_{k}\left(x_{0}, y_{1}, \mu\right)\right) \tag{15}
\end{align*}
$$

where $\tilde{\gamma}$ depends on $\varepsilon$ only and $\tilde{\gamma}=\max \left\{|\lambda(\mu)|^{-1},|\gamma(\mu)|\right\}$ for $\||\mu \|| \leq \varepsilon$.

## 2 Construction of the first return maps $T_{k}$

The domain of the first return map is the set of those points in $\Pi_{0}$ whose iterations by $T_{0}$ reach $\Pi_{1}$. One can easily see that it is a union of a countable number of horizontal strips $\sigma_{k}^{0}=\Pi_{0} \cap T_{0}^{-k} \Pi_{1}, k=\bar{k}, \bar{k}+1, \ldots$ which accumulate on the segment $\Pi_{0} \cap W_{\text {loc }}^{s}$ as $k \rightarrow \infty$ (Figure 3a). In turn, the images of the strips $\sigma_{k}^{0}$ by $T_{0}^{k}$ are vertical strips $\sigma_{k}^{1}=\Pi_{1} \cap T_{0}^{k} \Pi_{0}$ lying in $\Pi_{1}$. These strips accumulate on the segment $\Pi_{1} \cap W_{\text {loc }}^{u}$ as $k \rightarrow \infty$ (Figure 3b).
Indeed, if we define $\Pi_{0}$ and $\Pi_{1}$ by formula (8) for some sufficiently small $\varepsilon_{0}$ and $\varepsilon_{1}$, then it follows immediately from (15) that

$$
\begin{align*}
& \sigma_{k}^{0}=\left\{\left(x_{0}, y_{0}\right)| | x_{0}-x^{+} \mid \leq \varepsilon_{0}\right. \\
& \left.\gamma^{-k}\left(y^{-}-\varepsilon_{1}+O\left(\gamma^{-k}\right)\right) \leq y_{0} \leq \gamma^{-k}\left(y^{-}+\varepsilon_{1}+O\left(\gamma^{-k}\right)\right)\right\}  \tag{16}\\
& \sigma_{k}^{1}=\left\{\left(x_{1}, y_{1}\right) \mid \lambda^{k}\left(x^{+}-\varepsilon_{0}+O\left(\gamma^{-k}\right)\right) \leq x_{1} \leq \lambda^{k}\left(x^{+}+\varepsilon_{0}+O\left(\gamma^{-k}\right)\right)\right. \\
& \left.\left|y_{1}-y^{-}\right| \leq \varepsilon_{1}\right\}
\end{align*}
$$

The map $T_{k} \equiv T_{1} T_{0}^{k}$ is defined on the strip $\sigma_{k}^{0}$. Using (4) and (15) this map can be
represented in the following form

$$
\begin{align*}
& \bar{x}_{0}=x^{+}+a \lambda^{k} x_{0}+b\left(y_{1}-y^{-}\right)+e_{02}\left(y_{1}-y^{-}\right)^{2}+e_{11} \lambda^{k} x_{0}\left(y_{1}-y^{-}\right)+ \\
& +E_{1}\left(x_{0}, y_{1}-y^{-}, \mu\right), \\
& \gamma^{-k} \bar{y}_{1}\left(1+\tilde{\gamma}^{-k} \eta_{k}\left(\bar{x}_{0}, \bar{y}_{1}, \mu\right)\right)=\mu_{1}+ \\
& +\left[c \lambda^{k} x_{0}+f_{11} \lambda^{k} x_{0}\left(y_{1}-y^{-}\right)\right]\left(1+\tilde{\gamma}^{-k} \xi_{k}\left(x_{0}, y_{1}, \mu\right)\right)+d\left(y_{1}-y^{-}\right)^{2}+  \tag{17}\\
& +f_{12} \lambda^{k} x_{0}\left(y_{1}-y^{-}\right)+f_{20} \lambda^{2 k} x_{0}^{2}+f_{21} \lambda^{2 k} x_{0}^{2}\left(y_{1}-y^{-}\right)+f_{03}\left(y_{1}-y^{-}\right)^{3}+ \\
& +E_{2}\left(x_{0}, y_{1}-y^{-}, \mu\right),
\end{align*}
$$

where

$$
\begin{align*}
& E_{1}=O\left(|\lambda|^{k} \tilde{\gamma}^{-k}\right)+O\left(\left(y-y^{-}\right)^{3}\right)+|\lambda|^{k} O\left(\left(y-y^{-}\right)^{2}\right)+|\lambda|^{2 k} O\left(y-y^{-}\right),  \tag{18}\\
& \left.E_{2}=O\left(\lambda^{2 k} \tilde{\gamma}^{-k}\right)\right)+o\left(\left(y-y^{-}\right)^{3}\right)+|\lambda|^{k} o\left(\left(y-y^{-}\right)^{2}\right)+|\lambda|^{2 k} o\left(y-y^{-}\right) .
\end{align*}
$$

Hereafter we will adhere the following notation: the symbols $O(\varphi(k))$ and $o(\varphi(k))$ for some function $\varphi(k)$ denote functions (of all variables) which are bounded by a constant (resp., an infinitely small factor) times $\|\varphi(k)\|$ as $k \rightarrow+\infty$, and the same asymptotic behavior is shown by all their derivatives up to the order $(r-2)$; the symbol $o\left(x^{m} y^{n}\right)$ denotes a function vanishing, if $n \neq 0$, at $y=0$ along with the first $n$ derivatives with respect to $y$ and, if $m \neq 0$, vanishing at $x=0$ along with the first $m$ derivatives with respect to $x$; the symbol $O\left(x^{m} y^{n}\right)$ denotes a function of the form const $\cdot x^{m} y^{n}+o\left(x^{m} y^{n}\right)$; below, we will also write $O(|f|+|g|)$ instead of $O(f)+O(g)$ and $o(|f|+|g|)$ instead of $o(f)+o(g)$, by esthetic reasons.
To study map (17), we use the rescaling technique. Thus, the following result takes place.

Rescaling lemma. By means of an affine transformation of coordinates and parameters the map $T_{k}$ can be brought to the form ${ }^{4}$

$$
\begin{align*}
& \bar{X}=Y-\frac{1}{b d}\left(e_{02}+\ldots\right) \gamma^{-k} Y^{2}+O\left(\tilde{\gamma}^{-2 k}\right), \\
& \bar{Y}=M_{1}-M_{2} X-Y^{2}+\left(2 a-\frac{b}{d} f_{11}+\ldots\right) \lambda^{k} X Y+  \tag{19}\\
& +\frac{1}{d^{2}}\left(f_{03}+\ldots\right) \gamma^{-k} Y^{3}+o\left(\gamma^{-k}\right) o\left(Y^{3}\right)+O\left(\tilde{\gamma}^{-2 k}\right)
\end{align*}
$$

where

$$
\begin{align*}
& M_{1}=-d \gamma^{2 k}\left[\mu_{1}-\gamma^{-k}\left(y^{-}+\ldots\right)+c \lambda^{k}\left(x^{+}+\ldots\right)\right]  \tag{20}\\
& M_{2}=-b c\left(1+\mu_{2}\right)^{k}(1+\ldots)
\end{align*}
$$

the dots stand for the terms, independent of $(X, Y)$, which tend to zero as $k \rightarrow+\infty$.
The new coordinates $(X, Y)$, as well as the new parameter $M_{1}$, run all finite values at sufficiently large $k$ and $M_{2}$ runs all finite positive values ${ }^{5}$. Namely, there exist positive constants $C_{1}, \ldots, C_{5}$ such that $X, Y, M_{1}, M_{2}$ take values in the following domain

$$
\begin{align*}
& |X| \leq C_{1} \varepsilon_{0}|\gamma|^{k},|Y| \leq C_{2} \varepsilon_{1}|\gamma|^{k} \\
& \left|M_{1}\right| \leq C_{3} \varepsilon \gamma^{2 k}, C_{4}(1-\varepsilon)^{k} \leq M_{2} \leq C_{5}(1+\varepsilon)^{k} \tag{21}
\end{align*}
$$

[^2]where $\varepsilon_{0}$ and $\varepsilon_{1}$ are the diameters of $\Pi_{0}$ and $\Pi_{1}$, respectively, and $\varepsilon$ is the size of the interval of variation of $\mu$.

Proof. First, we shift the origin to the point $\left(x^{+}, y^{-}\right)$, i.e. we introduce the coordinates $x=x_{0}-x^{+}, y=y_{1}-y^{-}$. The map (17) takes the form

$$
\begin{align*}
& \bar{x}=a \lambda^{k} x+a \lambda^{k} x^{+}+\left(b+e_{11} \lambda^{k} x^{+}\right) y+e_{11} \lambda^{k} x y+e_{02} y^{2}+ \\
& +O\left(|\lambda|^{k} \tilde{\gamma}^{-k}+\left(|y|+|\lambda|^{k}\right)^{3}\right), \\
& \left(\bar{y}+y^{-}\right)\left(1+\tilde{\gamma}^{-k} \eta_{k}\left(\bar{x}+x^{+}, \bar{y}+y^{-}, \mu\right)\right)= \\
& =\mu_{1} \gamma^{k}+\lambda^{k} \gamma^{k}\left(x+x^{+}\right)\left[c+f_{11} y\right]\left[\left(1+\tilde{\gamma}^{-k} \xi_{k}\left(x+x^{+}, y+y^{-}, \mu\right)\right]+\right.  \tag{22}\\
& +d \gamma^{k} y^{2}\left(1+\frac{f_{12}}{d} \lambda^{k}\left(x+x^{+}\right)\right) \lambda^{2 k} \gamma^{k} x^{2}\left(f_{20}+f_{21} y\right)+f_{03} \gamma^{k} y^{3}+ \\
& +O\left(\lambda^{k} \tilde{\gamma}^{-k}\left(1+\mu_{2}\right)^{k}\right)+\gamma^{k} o\left(\left(\left|y_{1}-y^{-}\right|+|\lambda|^{k}\right)^{3}\right)
\end{align*}
$$

By an additional shift of coordinates

$$
x \rightarrow x_{\text {new }}+a \lambda^{k} x^{+}-\frac{b f_{11} x^{+}}{2 d} \lambda^{k}+o\left(\lambda^{k}\right), y \rightarrow y_{\text {new }}-\frac{b f_{11} x^{+}}{2 d} \lambda^{k}+o\left(\lambda^{k}\right)
$$

we can bring the map to such form where the right-hand side of the first equation does not contain the constant term and the right-hand side of the second equation does not contain the term linear in $y$. Then, (22) is rewritten as follows

$$
\begin{align*}
& \bar{x}=(a+\ldots) \lambda^{k} x+b y\left(1+\nu_{1} \lambda^{k}\right)++\left(e_{02}+\ldots\right) y^{2}+ \\
& +O\left(|\lambda|^{k}|x|(|y|+|x|)+|y|^{3}\right) \\
& \bar{y}+\tilde{\gamma}^{-k} O(|\bar{y}|+|\bar{x}|)=\gamma^{k}\left[\mu_{1}-\gamma^{-k}\left(y^{-}+\ldots\right)+c \lambda^{k}\left(x^{+}+\ldots\right)\right]+  \tag{23}\\
& (c+\ldots) \lambda^{k} \gamma^{k} x+\left(f_{11}+\ldots\right) \lambda^{k} \gamma^{k} x y+d \gamma^{k} y^{2}\left(1+\nu_{2} \lambda^{k}\right)+f_{03} \gamma^{k} y^{3}+ \\
& +O\left(|\lambda \gamma|^{k}|x| y^{2}+\lambda^{2 k}|\gamma|^{k} x^{2}\right)+\gamma^{k} o\left(y^{3}\right)
\end{align*}
$$

where $\nu_{1}=e_{11} x^{+} / b-e_{02} f_{11} x^{+} /(b d)+\ldots, \nu_{2}=f_{12} x^{+} / d-3 f_{11} f_{03} x^{+} / d+\ldots$; the dots stand for terms asymptotically vanishing as $k \rightarrow \infty$ and independent of $(x, y)$. Let us make the following rescaling of the variables

$$
x=-\frac{b\left(1+\nu_{1} \lambda^{k}\right)}{d\left(1+\nu_{2} \lambda^{k}\right)} \gamma^{-k} X, y=-\frac{1}{d\left(1+\nu_{1} \lambda^{k}\right)} \gamma^{-k} Y
$$

Then, map (23) is brought to the form

$$
\begin{align*}
& \bar{X}=Y+(a+\ldots) \lambda^{k} X-\frac{1}{b d}\left(e_{02}+\ldots\right) \gamma^{-k} Y^{2}+O\left(|\lambda|^{k} \tilde{\gamma}^{-k}\right) \\
& \bar{Y}+\rho_{k}^{1} \bar{Y}+\rho_{k}^{2} \bar{X}=\hat{M}_{1}+b c(\lambda \gamma)^{k} X(1+\ldots)-Y^{2}-\frac{b}{d}\left(f_{11}+\ldots\right) \lambda^{k} X Y+  \tag{24}\\
& +\left(\frac{f_{03}}{d^{2}}+\ldots\right) \gamma^{-k} Y^{3}+o\left(\gamma^{-k}\right) o\left(Y^{3}\right)+O\left(\tilde{\gamma}^{-2 k}\right)
\end{align*}
$$

where $\rho_{k}^{1,2}$ are some coefficients of order $\tilde{\gamma}^{-k}$, and

$$
\hat{M}_{1}=-d \gamma^{2 k}\left[\mu_{1}-\gamma^{-k}\left(y^{-}+\ldots\right)+c \lambda^{k}\left(x^{+}+\ldots\right)\right] .
$$

Substituting expressions for $\bar{X}$ into the left-hand side of the second equation of (24) we obtain

$$
\begin{align*}
& \bar{X}=Y+(a+\ldots) \lambda^{k} X-\frac{1}{b d}\left(e_{02}+\ldots\right) \gamma^{-k} Y^{2}+O\left(\tilde{\gamma}^{-2 k}\right), \\
& \bar{Y}=-\rho_{k}^{2} Y+\hat{M}_{1} /\left(1+\rho_{k}^{1}\right)+b c(\lambda \gamma)^{k} X(1+\ldots)-Y^{2}\left(1+\rho_{k}^{1}\right)^{-1}-  \tag{25}\\
& -\frac{b}{d}\left(f_{11}+\ldots\right) \lambda^{k} X Y+\frac{f_{03+\ldots}}{d^{2}} \gamma^{-k} Y^{3}+o\left(\gamma^{-k}\right) o\left(Y^{3}\right)+O\left(\tilde{\gamma}^{-2 k}\right) .
\end{align*}
$$

Denote the coefficient of the linear in $X$ term (taken with the sign "minus") in the second equation of (25) as $M_{2}$, i.e.

$$
M_{2}=-b c\left(\lambda^{k} \gamma\right)^{k}(1+\ldots) .
$$

Make the new change

$$
X=X_{\text {new }}\left(1+\rho_{k}^{1}\right), Y=Y_{\text {new }}\left(1+\rho_{k}^{1}\right)-(a+\ldots) \lambda^{k} X
$$

which brings the map to the form

$$
\begin{align*}
& \bar{X}=Y-\frac{1}{b d}\left(e_{02}+\ldots\right) \gamma^{-k} Y^{2}+O\left(\tilde{\gamma}^{-2 k}\right), \\
& \bar{Y}=-\rho_{k}^{3} Y+\hat{M}_{1} /\left(1+\rho_{k}^{1}\right)^{2}-M_{2} X-Y^{2}+\left(2 a-\frac{b}{d} f_{11}+\ldots\right) \lambda^{k} X Y+  \tag{26}\\
& +\frac{f_{03+}+\ldots \gamma^{-k} Y^{3}+o\left(\gamma^{-k}\right) o\left(Y^{3}\right)+O\left(\tilde{\gamma}^{-2 k}\right)}{d^{2}}
\end{align*}
$$

where the coefficient $\rho_{k}^{3}$ is of order $\tilde{\gamma}^{-k}$. One more shift of coordinates

$$
X \rightarrow X_{\text {new }}+\frac{1}{2} \rho_{k}^{3}, Y \rightarrow Y_{\text {new }}+\frac{1}{2} \rho_{k}^{3}
$$

eliminates the linear in $Y$ term in the second equation and brings, finally, the map to the requried form (19) (with the new constant term $M_{1}$ satisfying (20)). This completes the proof of the lemma.

## 3 Bifurcations of the fixed points of the first return map.

By the rescaling lemma, the map $T_{k}$ in the rescaled coordinates $X, Y$ is close to the standard quadratic Henon map

$$
\begin{equation*}
\bar{X}=Y, \quad \bar{Y}=M_{1}-M_{2}-Y^{2} \tag{27}
\end{equation*}
$$

where the Jacobian $M_{2}$ of (27) may take arbitrary finite positive values (see inequalities (20)). The bifurcations of the fixed points of the Henon family are well known. The corresponding bifurcational diagram on the half-plane $M_{2}>0$ is represented in Figure 4. It contains three bifurcational curves: $L^{+}: M_{1}=-\frac{1}{4}\left(1+M_{2}\right)^{2}$; $L^{-}: M_{1}=\frac{3}{4}\left(1+M_{2}\right)^{2}$ and $L^{\varphi}: \hat{M}_{2}=-1,-1<\hat{M}_{1}<3$. At $\left(M_{1}, M_{2}\right) \in L^{+}$
map (27) possesses a fixed point with a multiplier $\nu_{1}=+1$; at $M_{2} \neq 1$ this point is a saddle-node either with the stable sector at $M_{2}<1$ or with the unstable sector at $M_{2}>1$. At $\left(M_{1}, M_{2}\right) \in L^{-}$map (27) possesses a fixed point with a multiplier $\nu_{1}=-1$; this point is stable at $M_{2}<1$ and unstable at $M_{2}>1$. At $M_{2}=1$ the Henon map is area-preserving. It is known that this map has a parabolic fixed point ( $\nu_{1}=\nu_{2}=1$ ) of unstable type at $M_{1}=-1$; a parabolic fixed point ( $\nu_{1}=\nu_{2}=-1$ ) of stable type at $M_{1}=3$ and an elliptic fixed point ( $\nu_{1,2}=e^{ \pm i \psi}$ ) at $-1<M_{1}<3$ (i.e., at $\left(M_{1}, M_{2}\right) \in L^{\varphi}$ ). The elliptic fixed point is generic if $\psi \notin\{\pi / 2,2 \pi / 3, \arccos (-1 / 4)\}[29]$.


Figure 4
Bifurcation diagram for fixed
points in the Henon map for $\mathrm{M}_{2}>0$

Since the rescaled map $T_{k}$ given by (19) is $C^{r-2}$-close to Henon map, it also has in the half-plane (defined by inequalities (20)) three bifurcational curves $L_{(k)}^{+}, L_{(k)}^{-}$and $L_{(k)}^{\varphi}$ close to the curves $L^{+}, L^{-}$and $L^{\varphi}$, respectively. Direct calculations give us the following equations for $L_{(k)}^{+}$and $L_{(k)}^{-}$:

$$
\begin{align*}
& L_{(k)}^{+}: M_{1}=-\frac{\left(1+M_{2}\right)^{2}}{4}\left(1+O\left(|\lambda|^{k}+|\gamma|^{-k}\right)\right), \\
& L_{(k)}^{-}: M_{1}=\frac{3\left(1+M_{2}\right)^{2}}{4}\left(1+O\left(|\lambda|^{k}+|\gamma|^{-k}\right)\right) . \tag{28}
\end{align*}
$$

The equation of $L_{(k)}^{\varphi}$ in a parametric form (where the argument $\psi$ of the multiplier is taken as a parameter) can be written as follows

$$
\begin{align*}
& M_{1}=\cos ^{2} \psi-2 \cos \psi+O\left(|\lambda|^{k}+|\gamma|^{-k}\right) \\
& M_{2}=1+\cos \psi\left[2 \gamma^{-k}\left(e_{02}+\ldots\right)-\left(2 a-\frac{b}{d} f_{11}+\ldots\right) \lambda^{k}\right] \tag{29}
\end{align*}
$$

where $0<\psi<\pi$.
In order to prove that a closed invariant curve is born while crossing the curve $L_{(k)}^{\varphi}$, it is necessary to calculate the first Lyapunov value $G_{1}$ at the weak focus (the fixed point of $T_{k}$ with the multipliers $e^{ \pm i \psi}$ ). We show that the following result takes place

Lemma 2 The following formula is valid for the first Lyapunov value $G_{1}$

$$
\begin{equation*}
G_{1}=\frac{R}{16(1-\cos \psi)} \cdot \lambda^{k}+o\left(\lambda^{k}\right) \tag{30}
\end{equation*}
$$

where $R$ is the separatrix value given by (5).
Proof. The first step of the proof is to write the Taylor expansion for the map $T_{k}$ given by (19) at the weak focus. Let us fix the values of parameters $M_{1}$ and $M_{2}$ so that $\left(M_{1}, M_{2}\right) \in L_{(k)}^{\varphi}$, i.e. $T_{k}$ has a fixed point $P_{k}$ with the multipliers $e^{ \pm i \psi}$, for some $\psi \in(0, \pi)$. We will denote the corresponding values of $M_{1}$ and $M_{2}$ as $M_{1}^{\psi}$ and $M_{2}^{\psi}$. By (20), we have $M_{2}^{\psi}=-b c \lambda^{k} \gamma^{k}(1+\ldots)$, so it follows from the second equation of (29) that

$$
\begin{equation*}
\gamma^{-k}=-b c \lambda^{k}(1+\ldots) \tag{31}
\end{equation*}
$$

for values of the parameters at the curve $L_{(k)}^{\varphi}$. Thus, map (19) near the weak focus can be written as follows

$$
\begin{align*}
& \bar{X}=Y+h_{02} \lambda^{k} Y^{2}+O\left(\tilde{\gamma}^{-2 k}\right), \\
& \bar{Y}=M_{1}-M_{2} X-Y^{2}+s_{11} \lambda^{k} X Y+s_{03} \lambda^{k} Y^{3}+  \tag{32}\\
& +o\left(\lambda^{k}\right) o\left(Y^{3}\right)+O\left(\tilde{\gamma}^{-2 k}\right)
\end{align*}
$$

where

$$
\begin{equation*}
h_{02}=\frac{c}{d} e_{02}+\ldots, s_{11}=2 a-\frac{b}{d} f_{11}+\ldots, s_{03}=-\frac{b c}{d^{2}} f_{03}+\ldots ; \tag{33}
\end{equation*}
$$

the dots stand for the terms, independent of $(X, Y)$, which tend to zero as $k \rightarrow+\infty$. Let us now find coordinates $\left(X^{\psi}, Y^{\psi}\right)$ of the fixed point of (32) at $M_{1}=M_{1}^{\psi}, M_{2}=$ $M_{2}^{\psi}$. We use that the trace of the characteristic matrix

$$
\left(\begin{array}{ll}
\frac{\partial \bar{X}}{\partial X} & \frac{\partial \bar{X}}{\partial Y}  \tag{34}\\
\frac{\partial \bar{Y}}{\partial X} & \frac{\partial \bar{Y}}{\partial Y}
\end{array}\right)
$$

of $T_{k}$ at the fixed point is equal to the sum of the multipliers, i.e. at the weak focus it equals to $2 \cos \psi$. We have

$$
-2 Y^{\psi}+s_{11} \lambda^{k} X^{\psi}+3 s_{03} \lambda^{k}\left(Y^{\psi}\right)^{2}+o\left(\lambda^{k}\right)=2 \cos \psi
$$

or, since $X^{\psi}=Y^{\psi}+h_{02} \lambda^{k}\left(Y^{\psi}\right)^{2}+o\left(\lambda^{k}\right)$,

$$
Y^{\psi}\left(1-\frac{1}{2} s_{11} \lambda^{k}-\frac{3}{2} s_{03} \lambda^{k} Y^{\psi}\right)=-\cos \psi+o\left(\lambda^{k}\right) .
$$

We obtain from this that

$$
\begin{align*}
& Y^{\psi}=-\cos \psi \cdot\left(1+\frac{1}{2} s_{11} \lambda^{k}+\frac{3}{2} s_{03} \lambda^{k} \cos \psi\right)+o\left(\lambda^{k}\right)  \tag{35}\\
& X^{\psi}=Y^{\psi}+h_{02} \lambda^{k} \cos ^{2} \psi+o\left(\lambda^{k}\right)
\end{align*}
$$

We shift now the origin to the point $\left(X^{\psi}, Y^{\psi}\right)$ and consider map (32) at $M_{1}=$ $M_{1}^{\psi}, M_{2}=M_{2}^{\psi}$. We obtain

$$
\begin{align*}
& \bar{X}=Y+h_{02} \lambda^{k} Y^{2}+o\left(\lambda^{k}\right), \\
& \bar{Y}=-X+2 Y \cos \psi+s_{11} \lambda^{k} X Y-  \tag{36}\\
& -Y^{2}\left(1+2 h_{02} \lambda^{k} \cos \psi+3 s_{03} \lambda^{k} \cos \psi\right)+s_{03} \lambda^{k} Y^{3}+o\left(\lambda^{k}\right) .
\end{align*}
$$

The following change of variables

$$
X=\alpha \cdot X_{\text {new }} \quad, \quad Y=\alpha \cdot Y_{\text {new }}
$$

where $\alpha=\left(1+2 h_{02} \lambda^{k} \cos \psi+3 s_{03} \lambda^{k} \cos \psi\right)^{-1}$, brings map (36) to the form

$$
\begin{align*}
& \bar{X}=Y+h_{02} \lambda^{k} Y^{2}+o\left(\lambda^{k}\right) \\
& \bar{Y}=-X+2 Y \cos \psi+s_{11} \lambda^{k} X Y-  \tag{37}\\
& -Y^{2}+s_{03} \lambda^{k} Y^{3}+o\left(\lambda^{k}\right) .
\end{align*}
$$

This map has a fixed point with multipliers $e^{ \pm i \psi}$ at the origin. Evidently, the first Lyapunov value of this point will coincide with that for map (36) up to terms of order $o\left(\lambda^{k}\right)$. It is also obvious that if we omit the $o\left(\lambda^{k}\right)$ terms in the right-hand side of (37), the first Lyapunov value will get only some $o\left(\lambda^{k}\right)$ corrections.
Thus, to prove the lemma it is enough to show that the first Lyapunov value of the fixed point (in the origin) of the map

$$
\begin{align*}
& \bar{X}=Y+h_{02} \lambda^{k} Y^{2} \\
& \bar{Y}=-X+2 Y \cos \psi-Y^{2}+s_{11} \lambda^{k} X Y+s_{03} \lambda^{k} Y^{3} . \tag{38}
\end{align*}
$$

satisfies formula (30)
The first Lyapunov value of the weak focus in the map (38) is a polynomial with respect to the coefficients $h_{02} \lambda^{k}, s_{11} \lambda^{k}, s_{03} \lambda^{k}$. Hence, it has the form

$$
\begin{equation*}
G_{1}=F_{0}+\lambda^{k} h_{02} \cdot F_{1}+\lambda^{k} s_{11} \cdot F_{2}+\lambda^{k} s_{03} \cdot F_{3}+O\left(\lambda^{2 k}\right) \tag{39}
\end{equation*}
$$

where $F_{0}, F_{1}, F_{2}$ and $F_{3}$ are some coefficients, depending on $\psi$. We note that if $h_{02}=s_{11}=0$, then map (37) is conservative, i.e. $G_{1}=0$ in this case. It means that $F_{0}=F_{3} \equiv 0$ in (39). Hence, the first Lyapunov value is independent of $s_{03}$ in the main order. Therefore, it is sufficient to prove (30) for the first Lyapunov value of the following quadratic map (i.e., at $s_{03}=0$ )

$$
\begin{align*}
& \bar{X}=Y+h_{02} \lambda^{k} Y^{2}, \\
& \bar{Y}=-X+2 Y \cos \psi-Y^{2}+s_{11} \lambda^{k} X Y . \tag{40}
\end{align*}
$$

By means of the linear change

$$
\begin{equation*}
X=\xi, Y=\cos \psi \cdot \xi-\sin \psi \cdot \eta \tag{41}
\end{equation*}
$$

map (40) is brought to the form where its linear part is the rotation

$$
\begin{align*}
& \bar{\xi}=\cos \psi \cdot \xi-\sin \psi \cdot \eta+h_{02} \lambda^{k}\left[\cos ^{2} \psi \cdot \xi^{2}-\sin 2 \psi \cdot \xi \eta+\sin ^{2} \psi \cdot \eta^{2}\right] \\
& \bar{\eta}=\sin \psi \cdot \xi+\cos \psi \cdot \eta+\xi^{2}\left[\frac{1+h_{02} \lambda^{k} \cos \psi}{\sin \psi} \cos ^{2} \psi-s_{11} \lambda^{k} \frac{\cos \psi}{\sin \psi}\right]+  \tag{42}\\
& +\xi \eta\left[-2 \cos \psi \cdot\left(1+h_{02} \lambda^{k} \cos \psi\right)+s_{11} \lambda^{k}\right]+ \\
& +\eta^{2} \sin \psi \cdot\left(1+h_{02} \lambda^{k} \cos \psi\right)
\end{align*}
$$

In the complex coordinates $z=\xi+i \eta, z^{*}=\xi-i \eta$ map (42) takes the form

$$
\begin{equation*}
\bar{z}=e^{i \psi} z+C_{20} z^{2}+C_{11} z z^{*}+C_{02}\left(z^{*}\right)^{2} \tag{43}
\end{equation*}
$$

where

$$
\begin{aligned}
& C_{20}=\left[\left(A+B^{\prime}-C\right)+i\left(A^{\prime}-B-C^{\prime}\right)\right] / 4 \\
& C_{11}=\left[2(A+C)+i 2\left(A^{\prime}+C^{\prime}\right)\right] / 4 \\
& C_{02}=\left[\left(A-B^{\prime}-C\right)+i\left(A^{\prime}+B-C^{\prime}\right)\right] / 4
\end{aligned}
$$

with $A, B, C$ and $A^{\prime}, B^{\prime}, C^{\prime}$ denoting the coefficients of the quadratic terms $\xi^{2}, \xi \eta, \eta^{2}$ in the right-hand sides of the equations for $\bar{\xi}$ and $\bar{\eta}$, respectively; i.e.

$$
\begin{aligned}
& A=h_{02} \lambda^{k} \cos ^{2} \psi, B=-h_{02} \lambda^{k} \sin 2 \psi, C=h_{02} \lambda^{k} \sin ^{2} \psi \\
& A^{\prime}=\left[\frac{1+h_{02} \lambda^{k} \cos \psi}{\sin \psi} \cos ^{2} \psi-s_{11} \lambda^{k} \frac{\cos \psi}{\sin \psi}\right] \\
& B^{\prime}=\left[-2 \cos \psi \cdot\left(1+h_{02} \lambda^{k} \cos \psi\right)+s_{11} \lambda^{k}\right] \\
& C^{\prime}=\sin \psi \cdot\left(1+h_{02} \lambda^{k} \cos \psi\right)
\end{aligned}
$$

Thus, we have

$$
\begin{align*}
& C_{20}=\frac{1}{4}\left[-2 \cos \psi-h_{02} \lambda^{k}+s_{11} \lambda^{k}\right]+ \\
& +\frac{i}{4}\left[\frac{\cos 2 \psi}{\sin \psi}+\lambda^{k} \frac{\cos \psi}{\sin \psi}\left(h_{02}-s_{11}\right)\right] \\
& C_{11}=\lambda^{k} \frac{h_{02}}{2}+\frac{i}{2}\left[\frac{1}{\sin \psi}+\lambda^{k} \frac{\cos \psi}{\sin \psi}\left(h_{02}-s_{11}\right)\right]  \tag{44}\\
& C_{02}=\frac{1}{4}\left[2 \cos \psi+h_{02} \lambda^{k}\left(3 \cos ^{2} \psi-\sin ^{2} \psi\right)-s_{11} \lambda^{k}\right]+ \\
& +\frac{i}{4}\left[\frac{\cos 2 \psi}{\sin \psi}+h_{02} \lambda^{k} \frac{\cos \psi}{\sin \psi}\left(\cos ^{2} \psi-3 \sin ^{2} \psi\right)-s_{11} \lambda^{k} \frac{\cos \psi}{\sin \psi}\right] .
\end{align*}
$$

Make the following normalizing change of variables in (43)

$$
\begin{equation*}
w=z+A_{20} z^{2}+A_{11} z z^{*}+A_{02}\left(z^{*}\right)^{2} \tag{45}
\end{equation*}
$$

in order to nullify all quadratic terms in (43). It is not hard to see that the coefficients $A_{20}, A_{11}$ and $A_{02}$ should be taken as follows

$$
\begin{equation*}
A_{20}=\frac{C_{20}}{e^{i \psi}-e^{2 i \psi}}, A_{11}=\frac{C_{11}}{e^{i \psi}-1}, A_{02}=\frac{C_{02}}{e^{i \psi}-e^{-2 i \psi}} \tag{46}
\end{equation*}
$$

in this case all quadratic terms will be eliminated indeed, provided $\psi \neq 2 \pi / 3$ (in the latter case, the term $\left(z^{*}\right)^{2}$ is resonant and cannot be killed by smooth coordinate changes). Thus, map (43) is brought to the form (if $\psi \neq 2 \pi / 3$ )

$$
\begin{equation*}
\bar{w}=e^{i \psi}\left(w+G_{30} w^{3}+G_{21} w^{2} w^{*}+G_{12} w\left(w^{*}\right)^{2}+G_{03}\left(w^{*}\right)^{3}+O\left(|w|^{4}\right)\right. \tag{47}
\end{equation*}
$$

We note that, among the cubic terms, the term $w^{2} w^{*}$ is always resonant, the terms $w^{3}$ and $w\left(w^{*}\right)^{2}$ are always nonresonant, while the term $\left(w^{*}\right)^{3}$ is resonant only in the case where $\psi=\pi / 2$. Thus, if $\psi \neq \pi / 2,2 \pi / 3$, the last three terms may be eliminated by some cubic change of variables. But the coefficient $G_{21}$ is not changed and, hence, the map (47) can be brought to the form

$$
\begin{equation*}
\bar{w}=e^{i \psi}\left(w+G_{21} w^{2} w^{*}\right)+O\left(|w|^{4}\right) \tag{48}
\end{equation*}
$$

In the polar coordinates $(\rho, \alpha)$, where $w=\rho e^{i \alpha}$, the map (48) takes the following form

$$
\begin{equation*}
\bar{\rho}=\rho\left(1+G_{1} \cdot \rho^{2}\right)+O\left(\rho^{4}\right), \quad \bar{\alpha}=\alpha+\psi+B \rho^{2}+O\left(\rho^{3}\right) \tag{49}
\end{equation*}
$$

where $G_{1}=\operatorname{Re}\left(G_{21}\right), B=\operatorname{Im}\left(G_{21}\right)$. Thus, $\operatorname{Re}\left(G_{21}\right)$ is the first Lyapunov value.
Now we calculate the coefficient $G_{21}$. Using (43),(45) and (46) we obtain for $G_{21}$ the following expression

$$
G_{21}=2 A_{20} C_{11}+A_{11} C_{11}^{*}+A_{11} C_{20} e^{-2 i \psi}+2 A_{02} C_{02}^{*} e^{-i \psi}
$$

By virtue of (46)

$$
\begin{equation*}
G_{21}=C_{20} C_{11} \frac{2 e^{i \psi}-1}{\left(1-e^{i \psi}\right.} e^{-2 i \psi}--2 \frac{\left|C_{02}\right|^{2}}{\left(1-e^{3 i \psi}\right.}-\frac{\left|C_{11}\right|^{2}}{\left(1-e^{i \psi}\right.} \tag{50}
\end{equation*}
$$

and, hence,
or

$$
\begin{align*}
& G_{1}=\operatorname{Re}\left(C_{20} C_{11}\right) \frac{\cos 3 \psi-3 \cos 2 \psi+2 \cos \psi}{2(1-\cos \psi)}+ \\
& +\operatorname{Im}\left(C_{20} C_{11}\right) \frac{\sin 3 \psi-3 \sin 2 \psi+2 \sin \psi}{2(1-\cos \psi)}-\left|C_{02}\right|^{2}-\frac{1}{2}\left|C_{11}\right|^{2} . \tag{52}
\end{align*}
$$

Finally, (52) can be rewritten as

$$
\begin{align*}
& G_{1}=\frac{1}{2} \operatorname{Re}\left(C_{20} C_{11}\right)\left(3+\cos \psi-2 \cos ^{2} \psi\right)+ \\
& +\frac{1}{2} \operatorname{Im}\left(C_{20} C_{11}\right)\left(2 \sin \psi-2 \sin 2 \psi-\frac{1+\cos \psi}{\sin \psi}\right)--\left|C_{02}\right|^{2}-\frac{1}{2}\left|C_{11}\right|^{2} \tag{53}
\end{align*}
$$

where $C_{20}, C_{11}$ and $C_{02}$ are given by (44).
In the case where $s_{11}=0, h_{02}=0$ one can check directly that $G_{1} \equiv 0$. Indeed, it follows from (44) that in this case

$$
\begin{align*}
& \operatorname{Re}\left(C_{20} C_{11}\right)=-\frac{\cos 2 \psi}{8 \sin ^{2} \psi}, \operatorname{Im}\left(C_{20} C_{11}\right)=-\frac{\cos \psi}{4 \sin \psi} \\
& \left|C_{11}\right|^{2}=\frac{1}{4 \sin ^{2} \psi},\left|C_{02}\right|^{2}=\frac{1}{4} \cos ^{2} \psi+\frac{\cos ^{2} 2 \psi}{16 \sin ^{2} \psi} \tag{54}
\end{align*}
$$

Then, we obtain from (53) that

$$
\begin{aligned}
& G_{1}=-\frac{1}{16} \frac{\cos ^{2} 2 \psi}{\sin ^{2} \psi}\left(3+2 \cos \psi-4 \cos ^{2} \psi\right)- \\
& -\frac{1}{8} \frac{\cos \psi}{\sin \psi}\left(2 \sin \psi-2 \sin 2 \psi-\frac{1+\cos \psi}{\sin \psi}\right)- \\
& -\frac{1}{4} \cos ^{2} \psi-\frac{1}{16} \frac{\cos ^{2} 2 \psi}{\sin ^{2} \psi}-\frac{1}{8 \sin ^{2} \psi} \equiv 0 .
\end{aligned}
$$

We are now in the position to compute the coefficients $F_{1}$ and $F_{2}$ in formula (39) for $G_{1}$. To compute $F_{2}$, we may assume $h_{02}=0$. It follows from (44), that in this case

$$
\begin{align*}
& C_{20}=\frac{1}{4}\left[-2 \cos \psi+s_{11} \lambda^{k}\right]+\frac{i}{4}\left[\frac{\cos 2 \psi}{\sin \psi}-\frac{\cos \psi}{\sin \psi} s_{11} \lambda^{k}\right] \\
& C_{11}=\frac{i}{2}\left[\frac{1}{\sin \psi}-\frac{\cos \psi}{\sin \psi} s_{11} \lambda^{k}\right]  \tag{55}\\
& C_{02}=\frac{1}{4}\left[2 \cos \psi-s_{11} \lambda^{k}\right]++\frac{i}{4}\left[\frac{\cos 2 \psi}{\sin \psi}-\frac{\cos \psi}{\sin \psi} s_{11} \lambda^{k}\right] .
\end{align*}
$$

We obtain from here

$$
\begin{aligned}
& \operatorname{Re}\left(C_{20} C_{11}\right)=\left[\operatorname{Re}\left(C_{20} C_{11}\right)\right]_{f i n}+\frac{(\cos 2 \psi+1) \cos \psi}{8 \sin ^{2} \psi} \cdot s_{11} \lambda^{k}+O\left(\lambda^{2 k}\right) \\
& \operatorname{Im}\left(C_{20} C_{11}\right)=\left[\operatorname{Im}\left(C_{20} C_{11}\right)\right]_{f i n}+\frac{1+2 \cos ^{2} \psi}{8 \sin \psi} \cdot s_{11} \lambda^{k}+O\left(\lambda^{2 k}\right) \\
& \left|C_{11}\right|^{2}=\left[\left|C_{11}\right|^{2}\right]_{f i n}-\frac{\cos \psi}{2 \sin ^{2} \psi} \cdot s_{11} \lambda^{k}+O\left(\lambda^{2 k}\right) \\
& \left|C_{02}\right|^{2}=\left[\left|C_{02}\right|^{2}\right]_{f i n}-\frac{1}{4}\left(\cos \psi+\frac{\cos 2 \psi \cos \psi}{2 \sin ^{2} \psi}\right) \cdot s_{11} \lambda^{k}+O\left(\lambda^{2 k}\right)
\end{aligned}
$$

where $[\cdot]_{\text {fin }}$ denotes the finite part of the corresponding coefficients, i.e., its value at $\lambda^{k}=0$ according to (54). Substituting these expressions in formula (53) and collecting the terms of order $\lambda^{k}$, we find

$$
\begin{aligned}
& 16 F_{2}=\frac{(\cos 2 \psi+1) \cos \psi}{\sin ^{2} \psi}\left(3+2 \cos \psi-4 \cos ^{2} \psi\right)+ \\
& +\frac{1+2 \cos ^{2} \psi}{\sin \psi}\left(2 \sin \psi-2 \sin 2 \psi-\frac{1+\cos \psi}{\sin \psi}\right)+ \\
& +4 \cos \psi+\frac{2 \cos 2 \psi \cdot \cos \psi}{\sin ^{2} \psi}+\frac{4 \cos \psi}{\sin ^{2} \psi} .
\end{aligned}
$$

It is easy to check that

$$
F_{2}=\frac{1}{16(1-\cos \psi)}
$$

Analogously we compute that

$$
F_{1}=-\frac{2}{16(1-\cos \psi)}
$$

Thus, the following formula

$$
G_{1}=\frac{s_{11}-2 h_{02}}{16(1-\cos \psi)} \lambda^{k}+O\left(\lambda^{2 k}\right)
$$

defines the first Lyapunov value. Since $R=s_{11}-2 h_{02}$ (see (5) and (33)), this completes the proof.

Now we are able to describe main bifurcations of fixed points of the rescaled first return map $T_{k}$. Such a map has three bifurcational curves $L_{(k)}^{+}, L_{(k)}^{-}$and $L_{(k)}^{\varphi}$ which divide the half-plane $M_{2}>0$ of parameters $M_{1}$ and $M_{2}$ into four parts $D_{0}, D_{1}^{s}$, $D_{1}^{u}$ and $D_{2}$ (see figure 5). The map $T_{k}$ has no fixed points for $\left(M_{1}, M_{2}\right) \in D_{0}$. Transitions from $D_{0}$ into regions $D_{1}^{s}$ and $D_{1}^{u}$ across the curve $L_{(k)}^{+}$corresponds usual saddle-node bifurcations which are the same as in the Henon map (27). Thus, map $T_{k}$ has exactly two fixed points $Q_{k}$ and $P_{k}$ for values of the parameters from the region from above the curve $L_{(k)}^{+}$. The point $Q_{k}$ is a saddle and $P_{k}$ is a stable point when $\left(M_{1}, M_{2}\right) \in D_{1}^{s}$ and a completely unstable point when $\left(M_{1}, M_{2}\right) \in D_{1}^{u}$; it is a

[^3]weak focus when $\left(M_{1}, M_{2}\right) \in L_{(k)}^{\varphi}$. Thus, the curve $L_{(k)}^{\varphi}$ is the boundary of stability of $P_{k}$. In the case where $R \neq 0$ we can define the type of this boundary, i.e., whether the boundary is "safe" or "dangerous" [32]. Namely, it depends on the sign of the value $R \lambda^{k}<0$ as follows.


Figure 5
Elements of the bifurcation diagram for the rescalling map in the case $\mathrm{R} \lambda^{k}>0$

If $R \lambda^{k}<0$, the boundary $L_{(k)}^{\varphi}$ is "safe": the first Lyapunov value $G_{1}$ is negative, so at the transition across $L_{(k)}^{\varphi}$ (except for two points on $L_{(k)}^{\varphi}$ where $\psi=\pi / 2,2 \pi / 3$, see formula (29)) in the direction from $D_{1}^{s}$ to $D_{1}^{u}$ (towards the increase of $M_{2}$ ) the orbit $P_{k}$ becomes unstable and a stable invariant curve is born from it. Moreover, $P_{k}$ is the stable weak focus at $\left(M_{1}, M_{2}\right) \in L_{(k)}^{\varphi} \backslash\{\pi / 2,2 \pi / 3\}$. This boundary is "safe" because, just after the loss of stability, iterations of any initial point close to $P_{k}$ approach the
stable invariant curve and, hence, do not escape a small neighbourhood of the fixed point. The corresponding bifurcation picture is shown in figure 5 .
If $R \lambda^{k}>0$, the boundary $L_{(k)}^{\varphi}$ is "dangerous": the first Lyapunov value $G_{1}$ is positive and ,here, at the transition across $L_{(k)}^{\varphi}$ (except for two points where $\psi=\pi / 2,2 \pi / 3$ ) in the direction from $D_{1}^{u}$ to $D_{1}^{s}$ an unstable invariant curve is born from $P_{k}$. Moreover, $P_{k}$ is the unstable weak focus at $\left(M_{1}, M_{2}\right) \in L_{(k)}^{\varphi} \backslash\{\pi / 2,2 \pi / 3\}$. This boundary is "dangerous" because the loss of stability is happened when the unstable invariant curve "merges" into the stable point $P_{k}$ and, as a result, iterations of any initial point close to $P_{k}$ (except for $P_{k}$ itself) escape any sufficiently small neighbourhood of the fixed point.

To conclude this section, we note that the importance of the separatrix value $R$ goes beyond the fact that its sign defines the sign of the first Lyapunov value at the weak focus.

Indeed, for example, the Jacobian $J$ of the map

$$
\begin{align*}
& \bar{X}=Y+h_{02} \lambda^{k} Y^{2}, \\
& \bar{Y}=M_{1}-M_{2} X-Y^{2}+s_{11} \lambda^{k} X Y+s_{03} \lambda^{k} Y^{3} \tag{56}
\end{align*}
$$

calculated in the weak focus (or, which is the same, the Jacobian of map (38) in the origin) is given by

$$
\begin{equation*}
J=1-R \lambda^{k} Y+O\left(\lambda^{2 k}\right) \tag{57}
\end{equation*}
$$

That is, $J$ differs from 1 on a value which is proportional to $R \lambda^{k}$ in the main order. The second observation is that the shape of the bifurcational curve $L_{(k)}^{\varphi}$ of map (56) depends essentially on the coefficient $R \lambda^{k}$. Indeed, the equation of this curve has the form (put $\gamma^{-k}=-b c \lambda^{k}(1+\ldots$ ) in (29)):

$$
\begin{align*}
& M_{1}=\cos ^{2} \psi-2 \cos \psi+O\left(\lambda^{k}\right) \\
& M_{2}=1+\cos \psi \cdot R \lambda^{k}+O\left(\lambda^{2 k}\right) \tag{58}
\end{align*}
$$

where $\psi$ is the parameter, $0<\psi<\pi$. We see that the curve $L_{(k)}^{\varphi}$ at $R \neq 0$ is not $C^{1}$-close to the curve $L^{\varphi}\left(M_{2}=1\right)$ for the Henon map, see figure.
An interesting curve (nonbifurcational) $L_{(k)}^{s}$ starts with the point $B_{(k)}^{++}$, which corresponds to the existence of a saddle fixed point of (56) of the neutral type, (i.e., the fixed point with multipliers $\lambda_{s}>0$ and $\lambda_{u}>0$ such that $\lambda_{s} \lambda_{u}=1$ ). This curve is drawn in figure as the dotted line, its equation is

$$
\begin{align*}
& M_{1}=\alpha^{2}-2 \alpha+O\left(\lambda^{k}\right) \\
& M_{2}=1+\alpha \cdot R \lambda^{k}+O\left(\lambda^{2 k}\right) \tag{59}
\end{align*}
$$

. where $\alpha=\left(\lambda_{s}+\lambda_{u}\right) / 2$ is the parameter, and $\alpha>1 .{ }^{7}$

[^4]We note that the curve $L_{(k)}^{s}$ may play an important role for the answer to the following principal question: how does the found invariant curve disappear? It is naturally to assume that this curve exists only for values of parameters $M_{1}$ and $M_{2}$ from some part of the region bounded by the curves $L_{(k)^{\varphi}}, L_{(k)}^{s}$ and $L_{(k)}^{-}$(the dashed region in figure). In any case, it follows from [5, 7] that no invariant closed curve exists for values of the parameters which are sufficiently far from this region. If to assume that our map can be embedded into the flow (this models the situation reasonably near the point $\left.B_{(k)}^{++}[31,24]\right)$, then the invariant curve should disappear merging into a homoclinic loop of the saddle. In this case, the homoclinic loop should have the same type of stability as the invariant curve [33, 34]. Thus this phenomenon can occur only for values of the parameters in the dashed region, because the saddle value $\lambda_{s} \lambda_{u}$ of the saddle fixed point is less or greater than one for values of the parameters from the left of $L_{(k)}^{s}$ or from the right of $L_{(k)}^{s}$, respectively.
In fact, the general mechanism of disappearance of the invariant curve is connected with its break-down [30]: the invariant curve becomes resonant one (in this case it contains alternating saddle and stable (or completely unstable) periodic orbits of the same period) and it is destroyed, typically, by one of the ways given in [30]. In this connection, the questions related to the existence of the resonant zones seem to be very interesting.

Another important question which we will consider in a forthcoming paper addreses the bifurcational phenomena accompanying the transition across the strong resonances $\psi=\pi / 2$ and $\psi=2 \pi / 3$ (and $\psi=0, \pi$ also).

All these questions may be studied both in the map (56) itself and in the map

$$
\begin{align*}
& \bar{x}=y+A \varepsilon y^{2} \\
& \bar{y}=M_{1}-M_{2} x-y^{2}+B \varepsilon x y+C \varepsilon y^{3} \tag{60}
\end{align*}
$$

where the parameters $M_{1}, M_{2}, A, B, C$ are arbitrary and the parameter $\varepsilon$ is sufficiently small. Map (60) can be considered as a practically interesting small perturbation of the standard Henon map.

## 4 The proof of theorems A and B.

We note, first of all, that the rescaling lemma allows to compute immediately the equations of bifurcational curves $L_{k}^{+}, L_{k}^{-}$and $L_{k}^{\varphi}, k=\bar{k}, \bar{k}+1, \ldots$, for the first return maps $T_{k}$ on the plane of the initial parameters ( $\mu_{1}, \mu_{2}$ ). Namely, using (31) and the relation $\lambda \gamma=1+\mu_{2}$, formulas (28) and (29) are transformed as follows

$$
\begin{aligned}
& L_{k}^{+}: \mu_{1}=\gamma-k\left(y^{-}-\left(1+\mu_{2}\right)^{k} c x^{+}\right)+r_{k}+\quad+\gamma^{-2 k} \frac{\left(1+b c\left(1+\mu_{2}\right)^{2}\right.}{4 d}(1+\ldots) \\
& L_{k}^{-}: \mu_{1}=\gamma-k\left(y^{-}-\left(1+\mu_{2}\right)^{k} c x^{+}\right)+r_{k}-\quad-\gamma^{-2 k} \frac{3\left(1+b c\left(1+\mu_{2}\right)^{2}\right.}{4 d}(1+\ldots) \\
& L_{k}^{\varphi}: \mu_{1}=\gamma-k\left(y^{-}+\frac{x^{+}}{b}+\ldots\right)+r_{k}-\quad-\gamma^{-2 k} \frac{\cos ^{2} \varphi-2 \cos \varphi}{4 d}(1+\ldots) \\
& \mu_{2}=-1+(-b c)^{-1 / k}\left(1-\frac{R \cos \varphi}{b c} \cdot \frac{\lambda^{k}}{k}(1+\ldots)\right)
\end{aligned}
$$

where $r_{k}=o\left(\gamma^{-k}\right)$ and the dots denote terms tending to zero as $k \rightarrow \infty$.
Evidently, curves $L_{k}^{+}$and $L_{k}^{-}$accumulate on the line $\mu_{1}=0$ corresponding to the diffeomorphisms possessing a (single-round) orbit of homoclinic tangency. Curves $L_{k}^{\varphi}$ connect points $B_{k}^{++}$and $B_{k}^{--}$on the curves $L_{k}^{+}$and $L_{k}^{-}$, respectively, and accumulate at the point $\mu_{1}=\mu_{2}=0$. The bifurcational part of theorem B (items 2 and 3 ) follows directly from our analysis of bifurcations of the first return maps (section 4 and 5). Finally, theorem A follows immediately from theorem B: the region $\Delta_{k}$ is some part of $D_{k}$ adjoining to $L_{k}^{\varphi}$ (on the segment $0<\psi<\pi / 2$, for example) from the left if $R \lambda^{k}>0$ and from the right if $R \lambda^{k}<0$.

## Acknowledgments.

This work was supported by grant INTAS No. 97-804, by grant of Russian Ministry of Education No. 97-0-1.8-59 and by DFG Schwerpunktprogramm DANSE. We are grateful to L.P.Shilnikov and D.V.Turaev for very useful discussions. The first author is grateful to Professor Klaus Schneider for his warm hospitality during the three month visit to WIAS in 1999-2000.

## References

[1] L.P. Shilnikov "On a Poincare-Birkhoff problem".- Math. USSR Sb., 1967, 3, pp.91-102.
[2] S.V.Gonchenko, D.V.Turaev, L.P.Shilnikov. "On models with a structurally unstable homoclinic curve".- Soviet Math.Dokl., 1992, 44, No.2, 422-426.
[3] S.V.Gonchenko, D.V.Turaev , L.P.Shilnikov "On models with non-rough Poincare homoclinic curves".- Physica D, 62 (1993), Nos.1-4, 1-14.
[4] S.V.Gonchenko, L.P.Shilnikov, D.V.Turaev "Dynamical phenomena in systems with structurally unstable Poincare homoclinic orbits".- Interdisc. J. CHAOS, 6 (1996), No.1, pp. 1-17.
[5] N.K.Gavrilov N.K., L.P.Shilnikov "On three-dimensional dynamical systems close to systems with a structurally unstable homoclinic curve".- Part 1, Math.USSR Sb., 1972, 17, pp.467-485 ; Part 2, Math.USSR Sb, 1973, 19, pp.139-156.
[6] S.V.Gonchenko "On stable periodic motions of systems close to a system with structurally unstable homoclinic curve".- Matematicheskie zametki, 33 (1983), No.5, pp. 745-755, in Russian.
[7] S.V.Gonchenko, D.V.Turaev, L.P.Shilnikov "Dynamical phenomena in multidimensional systems with a structurally unstable homoclinic Poincare curve".Russian Acad. Sci. Dokl. Math., 47 (1993), No.3, pp.410-415.
[8] S. Newhouse "The abundance of wild hyperbolic sets and non-smooth stable sets for diffeomorphisms".- Publ. Math. IHES, 1979. v.50, pp.101-151.
[9] J.Palis, M.Viana "High dimension diffeomorphisms displaying infinitely many periodic attractors".- Annals of Math., 1994, v.140, pp.207-250.
[10] S.V.Gonchenko, D.V.Turaev, L.P.Shilnikov "On the existence of Newhouse regions in a neighbourhood of systems with a structurally unstable homoclinic Poincaré curve (the multidimensional case)".- Russian Acad. Sci. Dokl. Math., 1993. v.47, No.2, pp.268-273.
[11] P.Duarte "Persistent homoclinic tangencies for conservative maps near the identity".- Preprint,6/98, Lisbon, March, 1998.
[12] P.Duarte "Abundance of elliptic isles at conservative bifurcations".Preprint,7/98, Lisbon, April, 1998.
[13] V.S.Afraimovich, L.P.Shilnikov "Strange attractors and quasiattractors".- in Nonlinear Dynamics and Turbulence, 1982, Pitman, Boston, pp.336-339.
[14] L.P.Shilnikov "Mathematical problems of nonlinear dynamics: a tutorial".Int.J. Bifurcation and Chaos, 1997, v.6, No.6, pp. 969-989.
[15] L.P.Shilnikov "Chua's circuit: rigorous results and future problems".- Int. J. of Bifurcation and Chaos, 1994, v.4, No.3, pp.489-519.
[16] Y.-C.Lai, C.Grebogi, J.A.Yorke, I.Kan "How often are chaotic saddles nonhyperbolic?"- Nonlinearity, 1993, v.6, No.5, pp.779-797.
[17] D.V.Turaev, L.P.Shilnikov "An example of a wild strange attractor".- Sbornik: Mathematics, 1998, v.189, N0.2, pp. 137-160.
[18] D.V.Turaev "On dimension of non-local bifurcational problems".- Bifurcation and Chaos, 1996, v.6.
[19] S.E.Newhouse "Diffeomorphisms with infinitely many sinks".- Topology, $\mathbf{1 3}$ (1974), pp. 9-18.
[20] S.V.Gonchenko, D.V.Turaev , L.P.Shilnikov "Homoclinic tangencies of arbitrary order in the Newhouse regions".- to appear.
[21] S.V.Gonchenko, L.P.Shilnikov "Invariants of $\Omega$-conjugacy of diffeomorphisms with a structurally unstable homoclinic trajectory".- Ukrainian Math.J., 1990, 42 , No.2, pp.134-140.
[22] S.V.Gonchenko, L.P.Shilnikov "On the moduli of systems with a non-rough Poincare homoclinic curve".- Russian Acad. Sci. Izv. Math., 1993, 41, No.3, pp.417-445.
[23] L.Tedeshini-Lalli, J.A.Yorke "How often do simple dynamical processes have infinitely many coexisting sinks?"- Comm.Math.Phys., 1995, v.106, pp.635-657.
[24] L.P.Shilnikov, A.L.Shilnikov, D.V.Turaev, L.O.Chua "Methods of qualitative theory in nonlinear dynamics, Part I".- World Scientific, 1998.
[25] L.P.Shilnikov "A contribution to the problem of structure of a neighbourhood of a homoclinic tube of invariant torus".- Soviet Math. Dokl., 1968, 180, No.2, pp.286-289. (Russian).
[26] V.S.Afraimovich, L.P.Shilnikov "On critical sets of Morse-Smale systems".Trans. Moscow Math. Soc., 1973, 28
[27] S.V.Gonchenko, L.P.Shilnikov "On two-dimensional analytic area-preserving diffeomorphisms with infinitely many stable elliptic periodic points".- Regular and Chaotic Dynamics, 1997, v.2,No.3/4, pp.106-123.
[28] E.Colli "Infinitely many coexisting strange attractors".- Ann. Inst. Poincare, 1998, v.15, No.5, pp.539-579.
[29] V.Biragov "Bifurcations in a two-parameter family of conservative mappings that are close to the Henon map".- Selecta Math. Sovietica, 1990, v.9, pp.247280 [Orig. publ. in Methods of the Qualitative Theory of Differential Equations, (E.A.Leontovich-Andronova, ed.), 1987, (Russian)].
[30] V.S.Afraimovich, L.P.Shilnikov "On invariant two-dimensional tori, their desintegration and stochasticity".- in Methods of the Qualitative Theory of Differential Equations, (E.A.Leontovich-Andronova, ed.), 1983, pp.3-26. (Russian). [Engl. transl. in: Amer.Math.Soc.Transl., 1991, v.149, pp.201-211].
[31] Kuznetsov Yu.A. "Elements of applied bifurcation theory".- Springer-Verlag, 1995.
[32] N.N.Bautin "Behavior of dynamical systems near the boundaries of stability regions".- OGIZ GOSTEXIZDAT, Leningrad-Moscow, 1949 (Russian).
[33] A.A.Andronov, E.A.Leontovich, I.E.Gordon, A.G.Majer "The theory of dynamical systems on a plane".- Israel Program of Scientific Translations, Jerusalem, 1973.
[34] N.N.Bautin, E.A.Leontovich "Methods and ways (examples) of the qualitative analysis of dynamical systems in a plane".- Moscow, "Nauka", 1976. (Russian)


[^0]:    ${ }^{1}$ For instance, diffeomorphisms with periodic orbits for which the first return map is locally identical on the center manifold are dense in the Newhouse regions [3, 20].

[^1]:    ${ }^{2}$ it is a consequence of the inequalities $|\lambda|<1$ and $|\gamma|>1$ : a monomial $x^{m} y^{n}$ in the first equation of (6) is resonant if $\lambda^{m-1} \gamma^{n}=1$, and in the second equation it is resonant if $\lambda^{m} \gamma^{n-1}=1$
    ${ }^{3}$ Note, that a smooth linearization is impossible here because of the resonance $\lambda \gamma=1$ at $\mu=0$.

[^2]:    ${ }^{4}$ Note that in comparison with the known results on the rescaling near homoclinic $[23,2,7,4]$ or heteroclinic [27,28] tangencies we compute also terms which are asymptotically small as $k \rightarrow+\infty$.
    ${ }^{5} M_{2}>0$ since $f_{\mu}$ preserves orientation.

[^3]:    ${ }^{6}$ Except for the point $B_{(k)}^{++}$where the map $T_{k}$ has the fixed point with two unit multipliers. It is not hard to check that the "parabolic" fixed point of $T_{k}$ is not degenerate in the sense that at $R \neq 0$ this point is the so-called Bogdanov-Takens fixed point [31]

[^4]:    ${ }^{7}$ Values $\alpha<-1$ corresponds to that part of the curve $L_{(k)}^{s}$ where the corresponding neutral type saddle fixed point of (56) has both negative multipliers; values $|\alpha|<1$ corresponds to the curve $L_{(k)}^{\varphi}$ - in this case $\alpha=\cos \psi$.

