

# No Smooth Julia Sets for Polynomial Diffeomorphisms of $\mathbb{C}^2$ with Positive Entropy

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**Abstract** For any polynomial diffeomorphism  $f$  of  $\mathbb{C}^2$  with positive entropy, the Julia set of  $f$  is never  $C^1$  smooth as a manifold-with-boundary.

**Keywords** Polynomial diffeomorphisms of  $\mathbb{C}^2$  · Julia set · Generalized Hénon maps

**Mathematics Subject Classification** 37F10

## 1 Introduction

There are several reasons why the polynomial diffeomorphisms of  $\mathbb{C}^2$  form an interesting family of dynamical systems. One of these is the fact that there are connections with two other areas of dynamics: polynomial maps of  $\mathbb{C}$  and diffeomorphisms of  $\mathbb{R}^2$ , which have each received a great deal of attention. Among the endomorphisms of  $\mathbb{P}^k$ , certain ones have more special, and regular, geometric structure.

The question arises whether, among the polynomial diffeomorphisms of  $\mathbb{C}^2$ , are there analogous special maps with special geometry? The Julia set of such a special map would be expected to have some smoothness. Here we show that this does not happen.

More generally, we consider a holomorphic mapping  $f : X \rightarrow X$  of a complex manifold  $X$ . The Fatou set of  $f$  is defined as the set of points  $x \in X$  where the iterates  $f^n := f \circ \cdots \circ f$  are locally equicontinuous. If  $X$  is not compact, then in

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the definition of equicontinuity, we consider the one point compactification of  $X$ ; in this case, a sequence which diverges uniformly to infinity is equicontinuous. By the nature of equicontinuity, the dynamics of  $f$  is regular on the Fatou set. The Julia set is defined as the complement of the Fatou set, and this is where any chaotic dynamics of  $f$  will take place. The first nontrivial case is where  $X = \mathbb{P}^1$  is the Riemann sphere, and in this case Fatou (see [17]) showed that if the Julia set  $J$  is a smooth curve, then either  $J$  is the unit circle or  $J$  is a real interval. If  $J$  is the circle, then  $f$  is equivalent to  $z \mapsto z^d$ , where  $d$  is an integer with  $|d| \geq 2$ ; if  $J$  is the interval, then  $f$  is equivalent to a Chebyshev polynomial. These maps with smooth  $J$  play special roles, and this sparked our interest to look for smooth Julia sets in other cases. (The higher dimensional case is discussed, for instance, in Nakane [19] and Uchimura [23, 24].)

Here we address the case where  $X = \mathbb{C}^2$ , and  $f$  is a polynomial automorphism, which means that  $f$  is biholomorphic, and the coordinates are polynomials. Since  $f$  is invertible, there are two Julia sets:  $J^+$  for iterates in forward time, and  $J^-$  for iterates in backward time. Polynomial automorphisms have been classified by Friedland and Milnor [12]; every such automorphism is conjugate to a map which is either affine or elementary, or it belongs to the family  $\mathcal{H}$ . The affine and elementary maps have simple dynamics, and  $J^\pm$  are (possibly empty) algebraic sets (see [12]).

Thus we will restrict our attention to the maps in  $\mathcal{H}$ , which are finite compositions  $f := f_k \circ \cdots \circ f_1$ , where each  $f_j$  is a generalized Hénon map, which by definition has the form  $f_j(x, y) = (y, p_j(y) - \delta_j x)$ , where  $\delta_j \in \mathbb{C}$  is nonzero, and  $p_j(y)$  is a monic polynomial of degree  $d_j \geq 2$ . The degree of  $f$  is  $d := d_1 \cdots d_k$ , and the complex Jacobian of  $f$  is  $\delta := \delta_1 \cdots \delta_k$ . In [12] and [22], it is shown that the topological entropy of  $f$  is  $\log d > 0$ . The dynamics of such maps is complicated and has received much study, starting with the papers [3, 11, 14, 15].

For maps in  $\mathcal{H}$ , we can ask whether  $J^+$  can be a manifold. For any saddle point  $q$ , the stable manifold  $W^s(q)$  is a Riemann surface contained in  $J^+$ . Thus  $J^+$  would have to have real dimension at least two. However,  $J^+$  is also the support of a positive, closed current  $\mu^+$  with continuous potential, and such potentials cannot be supported on a Riemann surface (see [3, 11]). On the other hand, since  $J^+ = \partial K^+$  is a boundary, it cannot have interior. Thus dimension 3 is the only possibility for  $J^+$  to be a manifold. In fact, there are examples of  $f$  for which  $J^+$  has been shown to be a topological 3-manifold (see [8, 11, 16, 20]). Fornæss and Sibony [11] have shown that  $J^+$  cannot be smooth for a generic element of  $\mathcal{H}$ .

The purpose of this paper is to prove the following:

**Theorem** *For any polynomial automorphism of  $\mathbb{C}^2$  of positive entropy, neither  $J^+$  nor  $J^-$  is smooth of class  $C^1$ , in the sense of manifold-with-boundary.*

We may interchange the roles of  $J^+$  and  $J^-$  by replacing  $f$  by  $f^{-1}$ , so there is no loss of generality if we consider only  $J^+$ .

In an Appendix, we discuss the nonsmoothness of the related sets  $J$ ,  $J^*$ , and  $K$ .

## 2 No Boundary

Let us start by showing that if  $J^+$  is a  $C^1$  manifold-with-boundary, then the boundary is empty. Recall that if  $J^+$  is  $C^1$ , then for each  $q_0 \in J^+$  there is a neighborhood  $U \ni q_0$

and  $r, \rho \in C^1(U)$  with  $dr \wedge d\rho \neq 0$  on  $U$ , such that  $U \cap J^+ = \{r = 0, \rho \leq 0\}$ . If  $J^+$  has boundary, it is given locally by  $\{r = \rho = 0\}$ . For  $q \in J^+$ , the tangent space  $T_q J^+$  consists of the vectors that annihilate  $dr$ . This contains the subspace  $H_q \subset T_q J^+$  consisting of the vectors that annihilate  $\partial r$ .  $H_q$  is the unique complex subspace inside  $T_q J^+$ , so if  $M \subset J^+$  is a complex submanifold, then  $T_q M = H_q$ .

We start by showing that if  $J^+$  is  $C^1$ , then it carries a Riemann surface lamination.

**Lemma 2.1** *If  $J^+$  is  $C^1$  smooth, then  $J^+$  carries a Riemann surface foliation  $\mathcal{R}$  with the property that if  $W^s(q)$  is the stable manifold of a saddle point  $q$ , then  $W^s(q)$  is a leaf of  $\mathcal{R}$ . If  $J^+$  is a  $C^1$  smooth manifold-with-boundary, then  $\mathcal{R}$  extends to a Riemann surface lamination of  $J^+$ . In particular, any boundary component is a leaf of  $\mathcal{R}$ .*

*Proof* Given  $q_0 \in J^+$ , let us choose holomorphic coordinates  $(z, w)$  such that  $dr(q_0) = dw$ . We work in a small neighborhood which is a bidisk  $\Delta_\eta \times \Delta_\eta$ . We may choose  $\eta$  small enough that  $|r_z/r_w| < 1$ . In the  $(z, w)$ -coordinates, the tangent space  $H_q$  has slope less than 1 at every point  $\{|z|, |w| < \eta\}$ . Now let  $\hat{q}$  be a saddle point, and let  $W^s(\hat{q})$  be the stable manifold, which is a complex submanifold of  $\mathbb{C}^2$ , contained in  $J^+$ . Let  $M$  denote a connected component of  $W^s(\hat{q}) \cap (\Delta_\eta \times \Delta_{\eta/2})$ . Since the slope is  $< 1$ , it follows that there is an analytic function  $\varphi : \Delta_\eta \rightarrow \Delta_\eta$  such that  $M \subset \Gamma_\varphi := \{(z, \varphi(z)) : z \in \Delta_\eta\}$ . Let  $\Phi$  denote the set of all such functions  $\varphi$ . Since a stable manifold can have no self-intersections, it follows that if  $\varphi_1, \varphi_2 \in \Phi$ , then either  $\Gamma_{\varphi_1} = \Gamma_{\varphi_2}$  or  $\Gamma_{\varphi_1} \cap \Gamma_{\varphi_2} = \emptyset$ . Now let  $\hat{\Phi}$  denote the set of all normal limits (uniform on compact subsets of  $\Delta_\eta$ ) of elements of  $\Phi$ . We note that by Hurwitz’s Theorem, the graphs  $\Gamma_\varphi, \varphi \in \hat{\Phi}$  have the same pairwise disjointness property. Finally, by [4],  $W^s(q_0)$  is dense in  $J^+$ , so the graphs  $\Gamma_\varphi, \varphi \in \hat{\Phi}$  give the local Riemann surface lamination.

If  $q_1$  is another saddle point, we may follow the same procedure and obtain a Riemann surface lamination whose graphs are given locally by  $\varphi \in \hat{\Phi}_1$ . However, we have seen that the tangent space to the foliation at a point  $q$  is given by  $H_q$ . Since these two foliations have the same tangent spaces everywhere, they must coincide.

We have seen that all the graphs are contained in  $J^+$ , so if  $J^+$  has boundary, then the boundary must coincide locally with one of the graphs. □

We will use the observation that  $K^+ \subset \{(x, y) \in \mathbb{C}^2 : |y| > \max(|x|, R)\}$ . Further, we will use the Green function  $G^-$  which has many properties, including

- (i)  $G^-$  is pluriharmonic on  $\{G^- > 0\}$ ,
- (ii)  $\{G^- = 0\} = K^-$ , and
- (iii)  $G^- \circ f = d^{-1}G^-$ .

Further, the restriction of  $G^-$  to  $\{|y| \leq \max(|x|, R)\}$  is a proper exhaustion.

**Lemma 2.2** *Suppose that  $J^+$  is a  $C^1$  smooth manifold-with-boundary, and  $M$  is a component of the boundary of  $J^+$ . Then  $M$  is a closed Riemann surface, and  $M \cap K \neq \emptyset$ .*

*Proof* We consider the restriction  $g := G^-|_M$ . If  $M \cap K = \emptyset$ , then  $g$  is harmonic on  $M$ . On the other hand,  $g$  is a proper exhaustion of  $M$ , which means that  $g(z) \rightarrow \infty$  as

$z \in M$  leaves every compact subset of  $M$ . This means that  $g$  must assume a minimum value at some point of  $M$ , which would violate the minimum principle for harmonic functions.  $\square$

**Lemma 2.3** *Suppose that  $J^+$  is a  $C^1$  smooth manifold-with-boundary. Then the boundary is empty.*

*Proof* Let  $M$  be a component of the boundary of  $J^+$ . By Lemma 2.2,  $M$  must intersect  $\Delta_R^2$ . Since  $J^+$  is  $C^1$ , there can be only finitely many boundary components of  $J^+ \cap \Delta_R^2$ . Thus there can be only finitely many components  $M$ , which must be permuted by  $f$ . If we take a sufficiently high iterate  $f^N$ , we may assume that  $M$  is invariant. Now let  $h := f^N|_M$  denote the restriction to  $M$ . We see that  $h$  is an automorphism of the Riemann surface  $M$ , and the iterates of all points of  $M$  approach  $K \cap M$  in forward time. It follows that  $M$  must have a fixed point  $q \in M$ , and  $|h'(q)| < 1$ . The other multiplier of  $Df$  at  $q$  is  $\delta/h'(q)$ .

We consider three cases. First, if  $|\delta/h'(q)| > 1$ , then  $q$  is a saddle point, and  $M = W^s(q)$ . On the other hand, by [4], the stable manifold of a saddle points is dense in  $J^+$ , which makes it impossible for  $M$  to be the boundary of  $J^+$ . This contradiction means that there can be no boundary component  $M$ .

The second case is  $|\delta/h'(q)| < 1$ . This case cannot occur because the multipliers are less than 1, so  $q$  is a sink, which means that  $q$  is contained in the interior of  $K^+$  and not in  $J^+$ .

The last case is where  $|\delta/h'(q)| = 1$ . In this case, we know that  $f$  preserves  $J^+$ , so  $Df$  must preserve  $T_q(J^+)$ . This means that the outward normal to  $M$  inside  $J^+$  is preserved, and thus the second multiplier must be  $+1$ . It follows that  $q$  is a semi-parabolic/semi-attracting fixed point. It follows that  $J^+$  must have a cusp at  $q$  and cannot be  $C^1$  (see Ueda [25] and Hakim [13]).  $\square$

### 3 Maps that Do Not Decrease Volume

We note the following topological result (see Samelson [21] for an elegant proof): If  $M$  is a smooth 3-manifold (without boundary) of class  $C^1$  in  $\mathbb{R}^4$ , then it is orientable. This gives:

**Proposition 3.1** *For any  $q \in M$ , there is a neighborhood  $U$  about  $q$  so that  $U - M$  consists of two components  $\mathcal{O}_1$  and  $\mathcal{O}_2$ , which belong to different components of  $\mathbb{R}^4 - M$ .*

*Proof* Suppose that  $\mathcal{O}_1$  and  $\mathcal{O}_2$  belong to the same component of  $\mathbb{R}^4 - M$ . Then we can construct a simple closed curve  $\gamma \subset \mathbb{R}^4$  which crosses  $M$  transversally at  $q$  and has no other intersection with  $M$ . It follows that the (oriented) intersection is  $\gamma \cdot M = 1$  (modulo 2). But the oriented intersection modulo 2 is a homotopy invariant (see [18]), and  $\gamma$  is contractible in  $\mathbb{R}^4$ , so we must have  $\gamma \cdot M = 0$  (modulo 2).  $\square$

**Corollary 3.2** *If  $J^+$  is  $C^1$  smooth, then  $f$  is an orientation preserving map of  $J^+$ .*

*Proof*  $U^+ := \mathbb{C}^2 - K^+$  is a connected (see [15]) and thus it is a component of  $\mathbb{C}^2 - J^+$ . Since  $f$  preserves  $U^+$ , it also preserves the orientation of  $J^+$ , which is  $\pm\partial U^+$ .  $\square$

We recall the following result of Friedland and Milnor:

**Theorem** ([12]) *If  $|\delta| > 1$ , then  $K^+$  has zero Lebesgue volume, and thus  $J^+ = K^+$ . If  $|\delta| = 1$ , then  $\text{int}(K^+) = \text{int}(K^-) = \text{int}(K)$ . In particular, there exists  $R$  such that  $J^+ = K^+$  outside  $\Delta_R^2$ .*

*Proof of Theorem in the case  $|\delta| \geq 1$ .* Let  $q \in J^+$  be a point outside  $\Delta_R^2$ , as in the Theorem above. Then near  $q$  there must be a component  $\mathcal{O}$ , which is distinct from  $U^+ = \mathbb{C}^2 - K^+$ . Thus  $\mathcal{O}$  must belong to the interior of  $K^+$ . But by the Theorem above, the interior of  $K^+$  is not near  $q$ . □

### 4 Volume Decreasing Maps

Throughout this section, we continue to suppose that  $J^+$  is  $C^1$  smooth, and in addition we suppose that  $|\delta| < 1$ . For a point  $q \in J^+$ , we let  $T_q := T_q(J^+)$  denote the real tangent space to  $J^+$ . We let  $H_q := T_q \cap iT_q$  denote the unique (one-dimensional) complex subspace inside  $T_q$ . Since  $J^+$  is invariant under  $f$ , so is  $H_q$ , and we let  $\alpha_q$  denote the multiplier of  $D_q f|_{H_q}$ .

**Lemma 4.1** *Let  $q \in J^+$  be a fixed point. There is a  $D_q f$ -invariant subspace  $E_q \subset T_q(\mathbb{C}^2)$  such that  $H_q$  and  $E_q$  generate  $T_q$ . We denote the multiplier of  $D_q f|_{E_q}$  by  $\beta_q$ . Thus  $D_q f$  is linearly conjugate to the diagonal matrix with diagonal elements  $\alpha_q$  and  $\beta_q$ . Further,  $\beta_q \in \mathbb{R}$  and  $\beta_q > 0$ .*

*Proof* We have identified an eigenvalue  $\alpha_q$  of  $D_q f$ . If  $D_q f$  is not diagonalizable, then it must have a Jordan canonical form  $\begin{pmatrix} \alpha_q & 1 \\ 0 & \alpha_q \end{pmatrix}$ . The determinant is  $\alpha_q^2 = \delta$ , which has modulus less than 1. Thus  $|\alpha_q| < 1$ , which means that  $q$  is an attracting fixed point and thus in the interior of  $K^+$ , not in  $J^+$ . Thus  $D_q f$  must be diagonalizable, which means that  $H_q$  has a complementary invariant subspace  $E_q$ . Since  $E_q$  and  $T_q$  are invariant under  $D_q f$ , the real subspace  $E_q \cap T_q \subset E_q$  is invariant, too. Thus  $\beta_q \in \mathbb{R}$ . By Corollary 3.2,  $D_q f$  will preserve the orientation of  $T_q$ , and so  $\beta_q > 0$ . □

Let us recall the Riemann surface foliation of  $J^+$  which was obtained in Lemma 2.1. For  $q \in J^+$ , we let  $R_q$  denote the leaf of  $\mathcal{R}$  containing  $q$ . If  $q$  is a fixed point, then  $f$  defines an automorphism  $g := f|_{R_q}$  of the Riemann surface  $R_q$ . Since  $R_q \subset K^+$ , we know that the iterates of  $g^n$  are bounded in a complex disk  $q \in \Delta_q \subset R_q$ . Thus the derivatives  $(Dg)^n = D(g^n)$  are bounded at  $q$ . We conclude that  $|\alpha_q| = |D_q(g)| \leq 1$ . If  $|\alpha_q| = 1$ , then  $\alpha_q$  is not a root of unity. Otherwise  $g$  is an automorphism of  $R_q$  fixing  $q$ , and  $Dg^n(q) = 1$  for some  $n$ . It follows that  $g^n$  must be the identity on  $R_q$ . This means that  $R_q$  would be a curve of fixed points for  $f^n$ , but by [FM] all periodic points of  $f$  are isolated, so this cannot happen.

**Lemma 4.2** *If  $q \in J^+$  is a fixed point, then  $q$  is a saddle point, and  $\alpha_q = \delta/d$ , and  $\beta_q = d$ .*

*Proof* First we claim that  $|\alpha_p| < 1$ . Otherwise, we have  $|\alpha_q| = 1$ , and by the discussion above, this means that  $\alpha_q$  is not a root of unity. Thus the restriction  $g = f|_{R_q}$  is

an irrational rotation. Let  $\Delta \subset R_q$  denote a  $g$ -invariant disk containing  $q$ . Since  $|\delta| = |\alpha_q \beta_q| = |\beta_q|$  has modulus less than 1, we conclude that  $f$  is normally attracting to  $\Delta$ , and thus  $q$  must be in the interior of  $K^+$ , which contradicts the assumption that  $q \in J^+$ .

Now we have  $|\alpha_q| < 1$ , so if  $|\beta_q| = 1$ , we have  $\beta_q = 1$ , since  $\beta_q$  is real and positive. This means that  $q$  is a semi-parabolic, semi-attracting fixed point for  $f$ . We conclude by Ueda [25] and Hakim [13] that  $J^+$  has a cusp at  $q$  and thus is not smooth. Thus we conclude that  $|\beta_q| > 1$ , which means that  $q$  is a saddle point.

Now since  $E_q$  is transverse to  $H_q$ , it follows that  $W^u(q)$  intersects  $J^+$  transversally, and thus  $J^+ \cap W^u(q)$  is  $C^1$  smooth. Let us consider the uniformization

$$\phi : \mathbb{C} \rightarrow W^u(q) \subset \mathbb{C}^2, \quad \phi(0) = q, \quad f \circ \phi(\zeta) = \phi(\lambda^u \zeta).$$

The pre-image  $\tau := \phi^{-1}(W^u(q) \cap J^+) \subset \mathbb{C}$  is a  $C^1$  curve passing through the origin and invariant under  $\zeta \mapsto \lambda^u \zeta$ . It follows that  $\lambda^u \in \mathbb{R}$ , and  $\tau$  is a straight line containing the origin. Further,  $g^+ := G^+ \circ \phi$  is harmonic on  $\mathbb{C} - \tau$ , vanishing on  $\tau$ , and satisfying  $g^+(\lambda^u \zeta) = d \cdot g^+(\zeta)$ . Since  $\tau$  is a line, it follows that  $g^+$  is piecewise linear, so we must have  $\lambda^u = \pm d$ . Finally, since  $f$  preserves orientation, we have  $\lambda^u = d$ .

**Lemma 4.3** *There can be at most one fixed point in the interior of  $K^+$ . There are at least  $d - 1$  fixed points contained in  $J^+$ , and at each of these fixed points, the differential  $Df$  has multiplier of  $d$ .*

*Proof* Suppose that  $q$  is a fixed point in the interior of  $K^+$ . Then  $q$  is contained in a recurrent Fatou domain  $\Omega$ , and by [4],  $\partial\Omega = J^+$ . If there is more than one fixed point in the interior of  $K^+$ , we would have  $J^+$  simultaneously being the boundary of more than one domain, in addition to being the boundary of  $U^+ = \mathbb{C}^2 - K^+$ . This is not possible if  $J^+$  is a topological submanifold of  $\mathbb{C}^2$ .

By [FM] there are exactly  $d$  fixed points, counted with multiplicity. By Lemma 4.3, the fixed points in  $J^+$  are of saddle type, so they have multiplicity 1. Thus there are at least  $d - 1$  of them. □

### 5 Fixed Points with Given Multipliers

If  $q = (x, y)$  is a fixed point for  $f = f_n \circ \dots \circ f_1$ , then we may represent it as a finite sequence  $(x_j, y_j)$  with  $j \in \mathbb{Z}/n\mathbb{Z}$ , subject to the conditions  $(x, y) = (x_1, y_1) = (x_{n+1}, y_{n+1})$  and  $f_j(x_j, y_j) = (x_{j+1}, y_{j+1})$ . Given the form of  $f_j$ , we have  $x_{j+1} = y_j$ , so we may drop the  $x_j$ 's from our notation and write  $q = (y_n, y_1)$ . We identify this point with the sequence  $\hat{q} = (y_1, \dots, y_n) \in \mathbb{C}^n$ , and we define the polynomials

$$\begin{aligned} \varphi_1 &:= p_1(y_1) - \delta_1 y_n - y_2 \\ \varphi_2 &:= p_2(y_2) - \delta_2 y_1 - y_3 \\ &\dots\dots\dots \\ \varphi_n &:= p_n(y_n) - \delta_n y_{n-1} - y_1. \end{aligned}$$

The condition to be a fixed point is that  $\hat{q} = (y_1, \dots, y_n)$  belongs to the zero locus  $Z(\varphi_1, \dots, \varphi_n)$  of the  $\varphi_i$ 's. We define  $q_i(y_i) := p_i(y_i) - y_i^{d_i}$  and  $Q_i := q_i(y_i) - y_{i+1} - \delta_i y_{i-1}$ , so

$$\varphi_i = y_i^{d_i} + q_i(y_i) - y_{i+1} - \delta_i y_{i-1} = y_i^{d_i} + Q_i \tag{*}$$

Since  $p_j$  is monic, the degrees of  $q_i$  and  $Q_i$  are  $\leq d_i - 1$ .

By the Chain Rule, the differential of  $f$  at  $q = (y_n, y_1)$  is given by

$$Df(q) = \begin{pmatrix} 0 & 1 \\ -\delta_n & p'_n(y_n) \end{pmatrix} \cdots \begin{pmatrix} 0 & 1 \\ -\delta_1 & p'_n(y_1) \end{pmatrix}$$

We will denote this by  $M_n = M_n(y_1, \dots, y_n) := \begin{pmatrix} m_{11}^{(n)} & m_{12}^{(n)} \\ m_{21}^{(n)} & m_{22}^{(n)} \end{pmatrix}$ .

We consider special monomials in  $p'_j = p'_j(y_j)$  which have the form  $(p')^L := p'_{\ell_1} \cdots p'_{\ell_s}$ , with  $L = \{\ell_1, \dots, \ell_s\} \subset \{1, \dots, n\}$ . Note that the factors  $p'_{\ell_i}$  in  $(p')^L$  are distinct. Let us use the notation  $|L|$  for the number of elements in  $L$  and  $H_{\mathbf{m}}$  for the linear span of  $\{(p')^L : |L| = m - 2k, 0 \leq k \leq n/2\}$ . With this notation,  $\mathbf{m}$  indicates the maximum number of factors of  $p'_j$  in any monomial, and in every case the number of factors differs from  $\mathbf{m}$  by an even number.

**Lemma 5.1** *The entries of  $M_n$ :*

- (1)  $m_{11}^{(n)}$  and  $m_{22}^{(n)} - p'_1(y_1) \cdots p'_n(y_n)$  both belong to  $H_{\mathbf{n}-2}$ .
- (2)  $m_{12}^{(n)}, m_{21}^{(n)} \in H_{\mathbf{n}-1}$ .

*Proof* We proceed by induction. The case  $n = 1$  is clear. If  $n = 2$ ,

$$M_2 = \begin{pmatrix} 0 & 1 \\ -\delta_2 & p'_2 \end{pmatrix} \begin{pmatrix} 0 & 1 \\ -\delta_1 & p'_1 \end{pmatrix} = \begin{pmatrix} -\delta_1 & p'_1 \\ -\delta_1 p'_2 & p'_1 p'_2 - \delta_2 \end{pmatrix}$$

which satisfies (1) and (2). For  $n > 2$ , we have

$$M_n = \begin{pmatrix} 0 & 1 \\ -\delta_n & p'_n \end{pmatrix} M_{n-1} = \begin{pmatrix} m_{21}^{(n-1)} & m_{22}^{(n-1)} \\ -\delta_n m_{11}^{(n-1)} + m_{21}^{(n-1)} p'_n & -\delta_n m_{12}^{(n-1)} + p'_n m_{22}^{(n-1)} \end{pmatrix}$$

which gives (1) and (2) for all  $n$ . □

The condition for  $Df$  to have a multiplier  $\lambda$  at  $q$  is  $\Phi(\hat{q}) = 0$ , where

$$\Phi = \det \left( M_n - \begin{pmatrix} \lambda & 0 \\ 0 & \lambda \end{pmatrix} \right)$$

**Lemma 5.2**  $\Phi - p'_1(y_1) \cdots p'_n(y_n) \in H_{\mathbf{n}-2}$ .

*Proof* The formula for the determinant gives

$$\Phi = \lambda^2 - \lambda \text{Tr}(M_n) + \det(M_n) = \lambda^2 - \lambda(m_{11}^{(n)} + m_{22}^{(n)}) + \delta$$

since  $\delta$  is the Jacobian determinant of  $Df$ . The Lemma now follows from Lemma 5.1. □

The degree of the monomial  $y^a := y_1^{a_1} \cdots y_n^{a_n}$  is  $\deg(y^a) = a_1 + \cdots + a_n$ . We will use the *graded lexicographical order* on the monomials in  $\{y_1, \dots, y_n\}$ . That is,  $y^a > y^b$  if either  $\deg(y^a) > \deg(y^b)$ , or if  $\deg(y^a) = \deg(y^b)$  and  $a_i > b_i$ , where  $i = \min\{1 \leq j \leq n : a_j \neq b_j\}$ . If  $f \in \mathbb{C}[y_1, \dots, y_n]$ , we denote  $\text{LT}(f)$  for the leading term of  $f$ ,  $\text{LC}(f)$  for the leading coefficient, and  $\text{LM}(f)$  for the leading monomial.

**Lemma 5.3** *With the graded lexicographical order,  $G := \{\varphi_1, \dots, \varphi_n\}$  is a Gröbner basis.*

*Proof* We will use Buchberger’s Algorithm (see [10, Chapter 2]). For each  $i = 1, \dots, n$ ,  $\text{LT}(\varphi_i) = \text{LM}(\varphi_i) = y_i^{d_i}$ , so for  $i \neq j$ , the least common multiple of the leading terms is  $\text{L.C.M.} = y_i^{d_i} y_j^{d_j}$ . The  $S$ -polynomial is

$$S(\varphi_i, \varphi_j) := \frac{\text{L.C.M.}}{\text{LM}(\varphi_j)} \varphi_i - \frac{\text{L.C.M.}}{\text{LM}(\varphi_i)} \varphi_j = y_j^{d_j} Q_i - y_i^{d_i} Q_j = \varphi_j Q_i - Q_j \varphi_i$$

where we use the  $Q_j$  from (4.1) and cancel terms. Now let  $\mu_i := \deg(Q_i)$ . Since  $\mu_i < d_i$  for all  $i$ , the monomials  $\text{LM}(\varphi_j Q_i) = y_j^{d_j} y_i^{\mu_i}$  and  $\text{LM}(\varphi_i Q_j) = y_i^{d_i} y_j^{\mu_j}$  are not equal in our monomial ordering. Thus  $\text{LM}(S(\varphi_i, \varphi_j)) \geq \max(\text{LM}(\varphi_j Q_i), \text{LM}(\varphi_i Q_j))$ . It follows from Buchberger’s Algorithm that  $\{\varphi_1, \dots, \varphi_n\}$  is a Gröbner basis. □

We will use the Multivariable Division Algorithm, by which any polynomial  $g \in \mathbb{C}[y_1, \dots, y_n]$  may be written as  $g = A_1 \varphi_1 + \cdots + A_n \varphi_n + R$  where  $\text{LM}(g) \geq \text{LM}(A_j \varphi_j)$  for all  $1 \leq j \leq n$ , and  $R$  contains no terms divisible by any  $\text{LM}(\varphi_j)$ . An important property of a Gröbner basis is that  $g$  belongs to the ideal  $\langle \varphi_1, \dots, \varphi_n \rangle$  if and only if  $R = 0$  (see, for instance, [10] or [1]).

If all fixed points have the same value of  $\lambda$  as multiplier, then it follows that  $\Phi$  must vanish on the whole zero set  $Z(\varphi_1, \dots, \varphi_n)$ . Since we have a Gröbner basis, we easily determine the following:

**Corollary 5.4**  $\Phi \notin \langle \varphi_1, \dots, \varphi_n \rangle$ .

*Proof* The leading monomial of  $\Phi$  is  $y_1^{d_1-1} \cdots y_n^{d_n-1}$ , but this is not divisible by any of the leading monomials  $\text{LM}(\varphi_j) = y_j^{d_j}$ . Since  $\{\varphi_1, \dots, \varphi_n\}$  is a Gröbner basis, it follows that  $\Phi$  does not belong to the ideal  $\langle \varphi_1, \dots, \varphi_n \rangle$ . □

### 6 Proof of the Theorem

In this section we prove the Theorem, which will follow from 4.3, in combination with:



**Proposition 6.1** *Suppose  $F = f_n \circ \dots \circ f_1$ ,  $n \geq 3$ , is a composition of generalized Hénon maps with  $|\delta| < 1$ . Suppose that  $F$  has  $d = d_1 \cdot \dots \cdot d_n$  distinct fixed points. It is not possible that  $d - 1$  of these points have the same multipliers.*

*Proof that Proposition 6.1 implies the Theorem* To prove the Theorem, it remains to deal with the case  $|\delta| < 1$ . If  $f = f_1$  is a single generalized Hénon map, we consider  $F = f_1 \circ f_1 \circ f_1$  with  $n = 3$  and the same Julia set. Lemma 3.4 asserts that if  $J^+$  is  $C^1$ , there are  $d - 1$  saddle points with unstable multiplier  $\lambda = d$ . So by Proposition 4.1, we conclude that  $J^+$  cannot be  $C^1$  smooth.  $\square$

We give the proof of Proposition 6.1 at the end of this section. For  $J \subset \{1, \dots, n\}$ , we write

$$\Lambda_J := \{(p')^L : L \subset J, |L| = |J| - 2k, \text{ for some, } 1 \leq k \leq |J|/2\},$$

We let  $H_J$  denote the linear span of  $\Lambda_J$ . To compare with our earlier notation, we note that  $H_J \subset H_{|J|-2}$  and that  $(p')^J \notin H_J$ . The elements of  $H_J$  depend only on the variables  $y_j$  for  $j \in J$ . Now we formulate a result for dividing certain terms by  $\varphi_j$ :

**Lemma 6.2** *Suppose that  $J \subset \{1, \dots, n\}$  and  $h \in H_J$ . Then for each  $j \in J$  and  $\alpha \in \mathbb{C}$ , we have*

$$(y_j - \alpha) \left( (p')^J + h \right) = A(y)\varphi_j + B(y) \left( (p')^{J-\{j\}} + \rho_1 \right) + (y_j - \alpha) \cdot \rho_2, \quad (\dagger)$$

where  $\rho_1, \rho_2 \in H_{J-\{j\}}$ , and  $B = \eta_j(y_j) + d_j y_{j+1} + d_j \delta_j y_{j-1}$  with

$$\eta_j(y_j) = y_j q'_j(y_j) - \alpha p'_j(y_j) - d_j q_j(y_j). \quad (\ddagger)$$

The leading monomials satisfy

$$LM \left( (y_j - \alpha) \left( (p')^J + h \right) \right) = LM(A(y)\varphi_j)$$

*Proof* Let us start with the case  $J = \{1, \dots, m\}$ ,  $m \leq n$ , and  $j = 1$ , so  $J - \{j\} = J_1 = \{2, \dots, n\}$ . We divide by  $p'_1$  and remove any factor of  $p'_1$  in  $h$ . This gives

$$(p')^J + h = p'_1(y_1)\mu_1 + \rho_2$$

where  $\mu_1 = (p')^{J_1} + \rho_1$ ,  $\rho_1, \rho_2 \in H_{\{2, \dots, m\}}$ , and  $\mu_1, \rho_1, \rho_2$  are independent of the variable  $y_1$ . Thus

$$\begin{aligned} (y_1 - \alpha) \left( (p')^J + h \right) &= (y_1 - \alpha)(d_1 y_1^{d_1-1} + q'_1(y_1))\mu_1 + (y_1 - \alpha)\rho_2 \\ &= d_1 y_1^{d_1} \mu_1 + (y_1 q'_1(y_1) - \alpha p'_1(y_1))\mu_1 + (y_1 - \alpha)\rho_2 \\ &= (d_1 \mu_1)\varphi_1 + (\eta_1(y_1) + d_1 y_2 + d_1 \delta_1 y_n)\mu_1 + (y_1 - \alpha)\rho_2 \end{aligned}$$

where in the last line we substitute  $\eta_1$  defined by  $(\ddagger)$ . Using  $(*)$ , we see that this gives  $(\dagger)$ .

It remains to look at the leading terms of  $T_1 := (y_1 - \alpha) ((p')^J + h)$  and  $T_2 := d_1 \mu_1 \varphi_1$ . We see that  $T_1$  and  $T_2$  both contain nonzero multiples of  $y_j \prod_{i=1}^m y_i^{d_i-1}$ , and all other monomials in  $T_1$  and  $T_2$  have lower degree. Thus we have  $LM(T_1) = LM(T_2)$  for the graded ordering, independent of any ordering on the variables  $y_1, \dots, y_n$ . The choices of  $J = \{1, \dots, m\}$  and  $j = 1$  just correspond to a permutation of variables, and this does not affect the conclusion that  $LM(T_1) = LM(T_2)$ .  $\square$

**Lemma 6.3** *For any  $\alpha \in \mathbb{C}$ ,  $(y_1 - \alpha)\Phi \notin \langle \varphi_1, \dots, \varphi_n \rangle$ .*

*Proof* By [FM], we may assume that  $p_j(y_j) = y_j^{d_j} + q_j(y_j)$  and  $\deg(q_j) \leq d_j - 2$ . We consider two cases. The first case is that there is at least one  $j$  such that  $\eta_j$  is not the zero polynomial. If we conjugate by  $f_{j-1} \circ \dots \circ f_1$ , we may “rotate” the maps in  $f$  so that the factor  $f_j$  becomes the first factor. If there exists a  $j$  for which  $\eta_j(y_j)$  is nonconstant, we choose this for  $f_1$ . Otherwise, if all the  $\eta_j$  are constant, we choose  $f_1$  to be any factor such that  $\eta_1 \neq 0$ .

We will apply the Multivariate Division Algorithm on  $(y_1 - \alpha)\Phi$  with respect to the set  $\{\varphi_1, \dots, \varphi_n\}$ . We will find that there is a nonzero remainder, and since  $\{\varphi_1, \dots, \varphi_n\}$  is a Gröbner basis, it will follow that  $(y_1 - \alpha)\Phi$  does not belong to the ideal  $\langle \varphi_1, \dots, \varphi_n \rangle$ .

We start with Lemma 5.2, according to which  $\Phi = p'_1 \cdots p'_n + h$ , where  $h \in H_{n-2} = H_{\{1, \dots, n\}}$ . The leading monomial of  $(y_1 - \alpha)\Phi$  is  $y_1^{d_1} \prod_{i=2}^n y_i^{d_i-1}$ , and  $\varphi_1$  is the only element of the basis whose leading monomial divides this. Thus we apply Lemma 6.2, with  $J = \{1, \dots, n\}$ ,  $j = 1$ , and  $J_{\hat{1}} := J - \{j\} = \{2, \dots, n\}$ . This gives

$$\begin{aligned} (y_1 - \alpha)\Phi &= A_1\varphi_1 + (\eta_1(y_1) + d_1y_2 + d_1\delta_1y_n) \left( \prod_{i=2}^n p'_i(y_i) + \rho_1 \right) + (y_1 - \alpha)\rho_2 \\ &= A_1\varphi_1 + \left[ d_1y_2 \left( (p')^{J_{\hat{1}}} + \rho_1 \right) \right] \\ &\quad + \left[ d_1\delta_1y_n \left( (p')^{J_{\hat{1}}} + \rho_1 \right) \right] + \left[ \eta_1 \left( (p')^{J_{\hat{1}}} + \rho_1 \right) \right] + \ell.o.t \\ &= A_1\varphi_1 + T_2 + T_n + R_1 + \ell.o.t \end{aligned}$$

where  $\rho_1, \rho_2 \in H_{\{2, \dots, n\}}$ . In particular,  $T_2$  and  $T_n$  depend on  $y_2, \dots, y_n$  but not on  $y_1$ . We note that  $T_2$  (respectively,  $T_n$ ) contains a term divisible by  $LM(\varphi_2)$  (respectively,  $LM(\varphi_n)$ ). We view  $R_1$  as a remainder term, and note that  $LM(R_1)$  is divisible by  $y_2^{d_2-1} \cdots y_n^{d_n-1}$ , as well as the largest power of  $y_1$  in  $\eta_1(y_1)$ . By “*ℓ.o.t.*,” we mean that none of its monomials is divisible by  $LM(R_1)$  or by any of the  $LM(\varphi_j)$ .

Now we apply Lemma 6.2 to  $T_2$ , this time with  $J = \{2, \dots, n\}$  and  $j = 2$ , with  $J - \{2\} = J_{\hat{2}} = \{3, \dots, n\}$ . We have

$$\begin{aligned} T_2 &= A_2\varphi_2 + d_2y_3((p')^{J_{\hat{2}}} + \rho_1^{(2)}) + d_2\delta_2y_1 \left( (p')^{J_{\hat{2}}} + \rho_1^{(2)} \right) + \eta_2(y_2)(p')^{J_{\hat{2}}} + \ell.o.t. \\ &= A_2\varphi_2 + T_2^{(2)} + R_1^{(2)} + R_2^{(2)} + \ell.o.t. \end{aligned}$$

We see that  $T_2^{(2)}$  contains terms that are divisible by  $\text{LM}(\varphi_3)$ , but the monomials in  $R_1^{(2)}$  and  $R_2^{(2)}$  are not divisible by  $\text{LM}(\varphi_i)$  for any  $i$ . The remainder term here is  $R_1^{(2)} + R_2^{(2)}$ , and we observe that this cannot cancel the largest term in  $R_1$ . This is because  $\text{LM}(R_1^{(2)})$  lacks a factor of  $y_2$ , and  $\text{LM}(R_2^{(2)})$  is equal to  $y_3^{d_3-1} \dots y_n^{d_n-1}$  times the largest power of  $y_2$  in  $\eta_2(y_2)$ , and by  $(\ddagger)$ , this power is no bigger than  $d_2 - 1$ . If  $\eta_1$  is not constant, then we see that  $\text{LM}(R_1) > \text{LM}(R_2^{(2)})$ . If  $\eta_1$  is constant, then  $\eta_2$  must be constant, too, and again we have  $\text{LM}(R_1) > \text{LM}(R_2^{(2)})$ . Thus, with our earlier notation,  $R_1^{(2)} + R_2^{(2)} = \ell.o.t.$

We do a similar procedure with  $T_n, T_2^{(2)}$ , etc., and again find that the remainder term does not contain a multiple of the leading monomial of  $R_1$ . We see that each time we do this process, the size of the exponent  $L$  decreases in the term  $(p')^L$ . When we have  $L = \emptyset$ , there are no terms that can be divided by any  $\text{LM}(\varphi_j)$ . Thus we end up with

$$(y_1 - \alpha)\Phi = A_1\varphi_1 + \dots + A_n\varphi_n + R_1 + \ell.o.t.$$

and  $\text{LT}((y_1 - \alpha)\Phi) \geq \text{LT}(A_j\varphi_j)$  for all  $1 \leq j \leq n$ , and none of the remaining terms is divisible by any of the leading monomials of  $\varphi_j$ . Thus we have now finished the Multivariate Division Algorithm, and we have a nonzero remainder. Thus  $(y_1 - \alpha)\Phi$  does not belong to the ideal of the  $\varphi_j$ 's.

Now we turn to the second case, in which  $\eta_j = 0$  for all  $j$ . By [12], we may assume that  $\text{deg}(q_j) \leq d_j - 2$ . It follows that  $\alpha = 0$  and  $q_j = 0$ . Thus  $p_j = y_j^{d_j}$  for all  $1 \leq j \leq n$ , so  $p'_j = d_j y_j^{d_j-1}$ , and  $H_J$  consists of linear combinations of products  $(p')^I = y_{i_1}^{d_{i_1}-1} \dots y_{i_k}^{d_{i_k}-1}$  for  $I = \{i_1, \dots, i_k\} \subset J$ , for even  $k \leq |J| - 2$ . We will go through the Multivariate Division Algorithm again. The principle is the same as before, but the details are different; in the first case we needed  $n \geq 2$ , and now we will need  $n \geq 3$ .

Again, it is only  $\varphi_1$  which has a leading monomial which can divide some terms in  $(y_1 - \alpha)\Phi$ . As before, we apply Lemma 6.2 with  $J = \{1, \dots, n\}$ ,  $j = 1$ , and  $J - \{1\} = J_1 = \{2, \dots, n\}$ . The polynomial in  $(\ddagger)$  becomes  $B = d_j y_{j+1} + d_j \delta_j y_{j-1}$ , and we have

$$\begin{aligned} y_1\Phi &= A_1\varphi_1 + d_1 y_2 \left( (p')^{J_1} + \rho_1 \right) + d_1 \delta_1 y_n \left( (p')^{J_1} + \rho_1 \right) + y_1 \rho_2 \\ &= A_1\varphi_1 + T_2 + T_n + \ell.o.t. \end{aligned}$$

where  $\rho_1, \rho_2 \in H_{\{2, \dots, n\}}$ . Now we apply Lemma 6.2 to divide  $T_2$  (respectively,  $T_n$ ) by  $\varphi_2$  (respectively,  $\varphi_n$ ). This yields

$$y_1\Phi = A_1\varphi_1 + A_2\varphi_2 + A_n\varphi_n + T_3 + T_n + R + \ell.o.t.,$$

where

$$T_3 = d_1 d_2 y_3 \left( (p')^{J_{12}} + \tilde{\rho}_3 \right), \quad T_n = d_1 d_n \delta_1 \delta_n y_{n-1} \left( (p')^{J_{1n}} + \tilde{\rho}_n \right)$$

with  $\tilde{\rho}_3 \in H_{\{3, \dots, n\}}$  and  $\tilde{\rho}_n \in H_{\{2, \dots, n-1\}}$ , and

$$R = \left( d_1 d_2 \delta_2 y_1 y_n^{d_n-1} + d_1 d_n \delta_1 y_1 y_2^{d_2-1} \right) \prod_{i=3}^{n-1} y_i^{d_i-1}$$

Since  $n > 2$ ,  $R$  is not the zero polynomial. We will continue the Multivariate Division Algorithm by dividing  $T_3$  by  $\varphi_3$  and  $T_n$  by  $\varphi_n$ , but we see that any terms created cannot cancel  $R$ . Thus when we finish the Multivariable Division Algorithm, we will have a nonzero remainder. As in the previous case, we conclude that  $y_1 \Phi$  is not in the ideal  $\langle \varphi_1, \dots, \varphi_n \rangle$ .

*Proof of Proposition 6.1.* The fixed points of  $f$  coincide with the elements of  $Z(\varphi_1, \dots, \varphi_n)$ , which is a variety of pure dimension zero. Saddle points have multiplicity 1, and since there are  $d - 1$  of these, and since the total multiplicity is  $d$ , there must be one more fixed point, also of multiplicity 1. It follows that the ideal  $I := \langle \varphi_1, \dots, \varphi_n \rangle$  is equal to its radical (see [1]). Since the saddle points all have multiplier  $\lambda$ ,  $\Phi$  must vanish at all the saddle points. If  $(\alpha, \beta)$  is the other fixed point, we conclude that  $(y_1 - \alpha)\Phi$  vanishes at all the fixed points. Thus  $(y_1 - \alpha)\Phi$  belongs to the radical of  $I$ , and thus  $I$  itself. This contradicts Lemma 6.3, which completes the proof of Proposition 6.1. □

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### Appendix: Nonsmoothness of $J$ , $J^*$ , and $K$

Let us turn our attention to other dynamical sets for polynomial diffeomorphisms of positive entropy. These are  $J := J^+ \cap J^-$ ,  $K := K^+ \cap K^-$ , and the set  $J^*$ , which coincides with the closure of the set of periodic points of saddle type. (See [3, 5], and [2] for other characterizations of  $J^*$ .) We have  $J^* \subset J \subset K$ . We note that none of these sets can be a smooth 3-manifold: otherwise, for any saddle point  $p$ , it would be a bounded set containing  $W^s(p)$  or  $W^u(p)$ , which is the holomorphic image of  $\mathbb{C}$ . The following was suggested by Remark 5.9 of Cantat in [9]; we sketch his proof:

**Proposition 6.1** *If  $J = J^*$ , then it is not a smooth 2-manifold.*

*Proof* Let  $p$  be a saddle point, and let  $W^u(p)$  be the unstable manifold. The slice  $J \cap W^u(p)$  is smooth and invariant under multiplication by the multiplier of  $Df$ . This means that in fact, the multiplier must be real, and the restriction of  $G^+$  to the slice must be linear on each (half-space) component of  $W^u(p) - J$ .

The identity  $G^+ \circ f = d \cdot G^+$  means that the canonical metric (defined in [6]) is multiplied by  $d$ . Thus  $f$  is quasi-expanding on  $J^*$ . Now, applying this argument to  $f^{-1}$  we get that  $f$  is quasi-hyperbolic. Further,  $J^* = J$ , so it is quasi-hyperbolic on  $J$ . If  $f$  fails to be hyperbolic, then by [7] there will be a one-sided saddle point, which cannot happen since  $J$  is smooth.

Now that  $f$  is hyperbolic on  $J$ , there is a splitting  $E^s \oplus E^u$  of the tangent bundle, so we conclude that  $J$  is a 2-torus. The dynamical degree must be the spectral radius of an

invertible 2-by-2 integer matrix, but this means it is not an integer, which contradicts the fact the dynamical degree of a Hénon map is its algebraic degree.

**Proposition 6.2** *Suppose that the complex Jacobian is not equal to  $\pm 1$ . Then for each saddle (periodic) point  $p$  and each neighborhood  $U$  of  $p$ , neither  $J \cap U$  nor  $J^* \cap U$  nor  $K \cap U$  is a  $C^1$  smooth 2-manifold.*

*Proof* Let us write  $M := J \cap U$  and  $g := f|_M$ . (The following argument works, too, if we take  $M = J^* \cap U$  or  $M = K \cap U$ .) The tangent space  $T_p M$  is invariant under  $Df$ . The stable/unstable spaces  $E^{s/u} \subset T_p \mathbb{C}^2$  are invariant under  $D_p f$ . The space  $E^s$  (or  $E^u$ ) cannot coincide with  $T_p M$ , for otherwise the complex stable manifold  $W^s(p)$  (or  $W^u(p)$ ) would be locally contained in  $M$ , and thus globally contained in  $J$ . But the  $W^{s/u}$  are uniformized by  $\mathbb{C}$ , whereas  $J$  is bounded. We conclude that  $p$  is a saddle point for  $g$ , and thus the local stable manifold  $W_{\text{loc}}^s(p; g)$  is a  $C^1$ -curve inside the complex stable manifold  $W^s(p)$ . As in Lemma 4.3, we conclude that the multiplier for  $D_p f|_{E_p^u}$  is  $\pm d$  and the multiplier for  $D_p f|_{E_p^s}$  is  $\pm 1/d$ . Thus the complex Jacobian is  $\delta = \pm 1$ .

*Solenoids* The two results above concern smoothness, but no example is known where  $J$ ,  $J^*$ , or  $K$  is even a topological 2-manifold. In the cases where  $J^+$  has been shown to be a topological 3-manifold (see [8, 11, 16, 20]) it also happens that  $J$  is a (topological) real solenoid, and in these cases it is also the case that  $J = J^*$ . Further, for every saddle (periodic) point  $p$ , there is a real arc  $\gamma_p = W_{\text{loc}}^u(p) \cap J$ . If we apply the argument of Proposition 6.2 to this case, we conclude that  $\gamma_p$  is not  $C^1$  smooth.

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