THE SADDLE-NODE SEPARATRIX-LOOP BIFURCATION*

STEPHEN SCHECTER†

Abstract. We study vector fields $\dot{x} = f(x)$, $x \in \mathbb{R}^2$, having at some point an equilibrium of saddle-node type with a separatrix loop. Such vector fields fill a codimension two submanifold Σ of an appropriate Banach space. We give analytic conditions that determine whether a two-parameter perturbation of $\dot{x} = f(x)$ is transverse to Σ . The new condition is a version of Melnikov's integral around the separatrix loop. If it is nonzero, then as one perturbs away from $\dot{x} = f(x)$ in the direction in which an equilibrium of saddle-node type persists, the separatrix loop breaks in a nondegenerate manner. This integral is shown to be nonzero for the two-parameter pendulum equation $\beta \ddot{\phi} + \dot{\phi} + \sin \phi = \rho$ at its organizing center.

Key words. saddle-node separatrix-loop bifurcation, Melnikov integral, pendulum, Josephson junction

AMS(MOS) subject classification. 58F14

1. Introduction. We shall be concerned with vector fields

(1)
$$\dot{x} = f(x), \qquad x \in \mathbb{R}^2$$

having at some $p \in \mathbb{R}^2$ an equilibrium of saddle-node type with a separatrix loop Γ (see Fig. 1). We assume that the saddle-node has one negative eigenvalue and, of course, one zero eigenvalue. Such vector fields fill a codimension two submanifold Σ of an appropriate Banach space of planar vector fields. Consider a two-parameter unfolding of (1),

(2)
$$\dot{x} = \tilde{f}(x, \nu_1, \nu_2), \qquad x \in \mathbb{R}^2, \quad \nu_1, \nu_2 \in \mathbb{R}$$

where $\tilde{f}(x, 0, 0) = f(x)$. We shall give a computable condition that determines whether the family (2) is transverse to Σ at $(\nu_1, \nu_2) = (0, 0)$.

If the transversality condition is satisfied, there is a smooth nonsingular change of coordinates in parameter space,

$$(\nu_1, \nu_2) \leftrightarrow (\mu_1, \mu_2), \qquad (0, 0) \leftrightarrow (0, 0),$$

such that

$$\dot{x} = \tilde{f}(x, \nu_1(\mu_1, \mu_2), \nu_2(\mu_1, \mu_2)) = f(x, \mu_1, \mu_2)$$

has the bifurcation diagram of Fig. 2 in a neighborhood of $(\mu_1, \mu_2) = (0, 0)$. The curve C lies in $\{(\mu_1, \mu_2): \mu_1 \le 0, \mu_2 \ge 0\}$. It has a quadratic tangency with the μ_2 -axis at

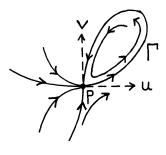


Fig. 1

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[†] Mathematics Department, North Carolina State University, Raleigh, North Carolina 27695-8205.

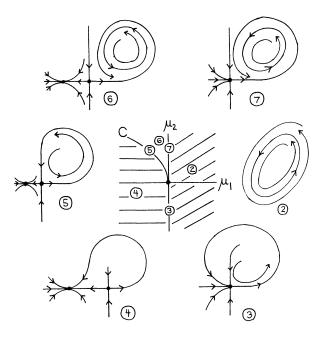


Fig. 2

- (0,0). The phase portrait of $\dot{x} = f(x, \mu_1, \mu_2)$ in a fixed neighborhood of Γ that is positively invariant for each vector field $\dot{x} = f(x, \mu_1, \mu_2)$ is as follows (see Fig. 2):
 - 1) $\mu_1 = 0$, $\mu_2 = 0$; a saddle-node and a separatrix loop.
 - 2) $\mu_1 > 0$; no equilibria, a unique stable closed orbit near Γ .
 - 3) $\mu_1 = 0$, $\mu_2 < 0$; a saddle-node.
 - 4) $\mu_1 < 0$, (μ_1, μ_2) below C; a saddle and a node.
 - 5) $\mu_1 < 0$, (μ_1, μ_2) on C; a saddle and a node; the saddle has a separatrix loop.
- 6) $\mu_1 < 0$, (μ_1, μ_2) above C; a saddle and a node; there is a unique stable closed orbit near Γ .
 - 7) $\mu_1 = 0$, $\mu_2 > 0$; a saddle-node and a unique stable closed orbit near Γ .

This bifurcation diagram is developed in [4], except that it is mistakenly stated there that the curve C is transverse to the μ_2 -axis.

Perhaps the best known example of this bifurcation diagram occurs in the study of the differential equation for a pendulum with linear damping and constant applied torque, which are the two parameters (see [3], in which the same equation arises in the study of the DC current-driven point Josephson junction). In § 5 we show that the pendulum equation satisfies our transversality condition at its organizing center. Thus the pendulum equation is a generic two-parameter unfolding of the saddle-node separatrix-loop bifurcation.

The heart of this paper is the study, in § 3, of Melnikov's integral (see [2]) around a saddle-node separatrix loop. The same method allows one to study time-periodic perturbations of a saddle-node separatrix loop. This subject is treated in the companion paper [6]. There the motivating example is the pendulum equation with, in addition, sinusoidal applied torque (or, equivalently, the AC-DC current-driven point Josephson junction).

- **2. Statement of results.** We shall consider vector fields $\dot{x} = f(x)$, $x \in \mathbb{R}^2$, satisfying the following conditions at some $p \in \mathbb{R}^2$:
 - (i) f(p) = 0.

(ii) Df(p) has eigenvalues 0 and $-\lambda$, where $\lambda > 0$.

Let u be a right eigenvector and w a left eigenvector of the eigenvalue 0, with w chosen so that wu > 0.

- (iii) $wD^2f(p)(u, u) > 0$.
- (iv) $\dot{x} = f(x)$ has a separatrix loop Γ at p.

Assumptions (i)-(iii) say that $\dot{x} = f(x)$ has a saddle-node at p with one negative eigenvalue (see [7]). Moreover, the assumptions imply that u (not -u) is one tangent vector to Γ at p (see Fig. 1). Let v be a right eigenvector of Df(p) for the eigenvalue $-\lambda$, chosen so that v is also tangent to Γ at p as in Fig. 1.

Let $\dot{x} = \tilde{f}(x, \nu_1, \nu_2)$ be a two-parameter family of vector fields on \mathbb{R}^2 such that $\dot{x} = \tilde{f}(x, 0, 0)$ satisfies (i)-(iv), and

- (v) $\tilde{f}(x, \nu_1, \nu_2)$ is $C^{k+1}, k \ge 5$.
- (vi) $wD_{\nu}, \tilde{f}(p, 0, 0) > 0$.

Assumptions (iii) and (vi) imply that perturbation in the positive ν_1 direction eliminates the equilibrium p, while perturbation in the negative ν_1 direction splits the equilibrium in two (see [7]).

According to [7] there is a C^k function $\alpha(\nu_2)$, with $\alpha(0) = 0$, such that $\dot{x} = \tilde{f}(x, \nu_1, \nu_2)$ has an equilibrium of saddle-node type near p if and only if $\nu_1 = \alpha(\nu_2)$. Let $f(x, \mu_1, \mu_2) = \tilde{f}(x, \nu_1, \nu_2)$, where (μ_1, μ_2) and (ν_1, ν_2) are related by

$$\mu_1 = \nu_1 - \alpha(\nu_2), \qquad \mu_2 = \nu_2.$$

Then

$$\dot{x} = f(x, \mu_1, \mu_2)$$

is C^k , and has an equilibrium of saddle-node type near p if and only if $\mu_1 = 0$. Let $p(\mu_2)$ denote the saddle-node equilibrium near p of $\dot{x} = f(x, 0, \mu_2)$; $p(\mu_2)$ is C^k . If $\mu_1 < 0$, there are a saddle and a sink of (3) near p; if $\mu_1 > 0$, there are no equilibria of (3) near p.

If w and z are vectors in \mathbb{R}^2 , let $w \wedge z = w_1 z_2 - w_2 z_1$. Let q(t) be a solution of $\dot{x} = f(x, 0, 0)$ with $q(0) \in \Gamma$. Consider the expression

(4)
$$I = \frac{dp}{d\mu_2}(0) \wedge \lim_{t_1 \to \infty} f(q(t_1), 0, 0) \exp\left[-\int_0^{t_1} \operatorname{div} f(q(s), 0, 0) \, ds\right] + \int_{-\infty}^{\infty} \exp\left[-\int_0^t \operatorname{div} f(q(s), 0, 0) \, ds\right] f(q(t), 0, 0) \wedge \frac{\partial f}{\partial \mu_2}(q(t), 0, 0) \, dt.$$

THEOREM 1. The limit in (4) exists and is a negative multiple of v. The improper integral in (4) converges. If $I \neq 0$, then there is a C^{k-2} curve $\mu_1 = \psi(\mu_2) = -\ell^2 \mu_2^2 + \varepsilon(\mu_2^2)$, $\ell \neq 0$, such that for (μ_1, μ_2) sufficiently small, (3) has a separatrix loop near Γ if and only if $\mu_1 = \psi(\mu_2)$ and $I \cdot (u \wedge v) \cdot \mu_2 \geq 0$.

The integral in (4) is just the usual Melnikov integral used to study perturbations of a separatrix loop at a hyperbolic saddle (see [2], [5]). The limit in (4) is zero in the case of a hyperbolic saddle, but must be retained in the case of a saddle-node.

The bifurcation diagram presented in § 1 holds if $I \cdot (u \wedge v) > 0$, in which case separatrix loops occur for $\mu_1 = \psi(\mu_2)$ and $\mu_2 \ge 0$. If $I \cdot (u \wedge v) < 0$, this bifurcation diagram holds after the further change of parameter $\mu_2 \to -\mu_2$.

We remark that in order to compute I in applications, the only knowledge of the function $\alpha(\nu_2)$ that is needed is $\alpha'(0)$.

We shall now give a precise interpretation of the condition $I \neq 0$ as a transversality condition. Let E denote the space of C^{k+1} vector fields, $k \ge 5$, on a closed disk $D \subseteq \mathbb{R}^2$,

with the C^{k+1} topology. Let $\Sigma' = \{f \in E : f \text{ satisfies, at a unique } p \in \text{Int } D$, conditions (i)-(iii); and all other equilibria of f in D are hyperbolic}. Σ' is a C^k codimension one submanifold of E [7]. Let $\Sigma = \{f \in \Sigma' : f \text{ satisfies condition (iv), and } \Gamma \subset \text{Int } D\}$. We shall show in § 3 that Σ is a C^k codimension two submanifold of E, in fact a codimension one submanifold of Σ' . Let $\dot{x} = \tilde{f}(x, \nu_1, \nu_2)$ be a two-parameter family of vector fields in E with $\dot{x} = \tilde{f}(x, 0, 0) \in \Sigma$. Assume in addition that $\dot{x} = \tilde{f}(x, \nu_1, \nu_2)$ satisfies conditions (v) and (vi). Make the change of parameters $\mu_1 = \nu_1 - \alpha(\nu_2)$, $\mu_2 = \nu_2$ described earlier.

THEOREM 2. The family $\dot{x} = \tilde{f}(x, \nu_1, \nu_2)$ is transverse to Σ at $(\nu_1, \nu_2) = (0, 0)$ if and only if $I \neq 0$.

3. Proof of Theorem 1. The equilibrium (p, 0, 0) of

(5)
$$\dot{x} = f(x, \mu_1, \mu_2), \quad \dot{\mu}_1 = 0, \quad \dot{\mu}_2 = 0$$

has a 3-dimensional neutral subspace. The center manifold theorem [1, § 9.2] yields a 3-dimensional C^k local center manifold N_{loc} of (5), tangent at (p, 0, 0) to this subspace. N_{loc} meets each plane $\mathbb{R}^2 \times \{(\mu_1, \mu_2)\}$, (μ_1, μ_2) small, in a curve. $N_{loc} \cap \mathbb{R}^2 \times \{(0, 0)\}$ contains a portion of $\Gamma \times \{(0, 0)\}$ that is tangent at (p, 0, 0) to (u, 0, 0).

Let N denote the global center manifold that contains N_{loc} , i.e., the union of all integral curves of (5) that meet N_{loc} . N meets each plane $\mathbb{R}^2 \times \{(\mu_1, \mu_2)\}$ in a curve, which we denote $N(\mu_1, \mu_2) \times \{(\mu_1, \mu_2)\}$. Thus $N(\mu_1, \mu_2)$ is a curve in \mathbb{R}^2 . Let L be a line segment in \mathbb{R}^2 perpendicular to Γ at q(0). Then for (μ_1, μ_2) small, $N(\mu_1, \mu_2)$ meets L transversally near q(0). Therefore for (μ_1, μ_2) small there is a C^k function $x(\mu_1, \mu_2)$ such that x(0, 0) = q(0) and $x(\mu_1, \mu_2) \in N(\mu_1, \mu_2) \cap L$. Since a C^k vector field has a C^k flow, there is a C^k family of solutions of (3)

$$q^{c}(\mu_{1}, \mu_{2}, t), (\mu_{1}, \mu_{2})$$
 small,

such that $q^c(\mu_1, \mu_2, 0) = x(\mu_1, \mu_2)$. Then $q^c(0, 0, t) = q(t)$, and each curve $q^c(\mu_1, \mu_2, t)$ lies in $N(\mu_1, \mu_2)$. For $\mu_1 < 0$, $q^c(\mu_1, \mu_2, t)$ is a branch of the unstable manifold of the saddle of (3) near p. Similarly, $q^c(0, \mu_2, t)$ is the unstable separatrix of the saddle-node of $\dot{x} = f(x, 0, \mu_2)$ near p (see Fig. 3).

We shall now define a μ -dependent change of coordinates on \mathbb{R}^2 that will make possible our computations. According to [1, § 9.2] there is a C^k change of coordinates

(6)
$$y(x, \mu_1, \mu_2) = (y_1(x, \mu_1, \mu_2), y_2(x, \mu_1, \mu_2)),$$

defined for (x, μ_1, μ_2) near (p, 0, 0), such that (1) y(p, 0, 0) = 0; $(2) N_{loc} \cap \mathbb{R}^2 \times \{(\mu_1, \mu_2)\}$ is transformed into the line $y_2 = 0$, which is therefore invariant; (3) the lines $y_1 = \text{constant}$ are mapped into each other by the flow. In other words, in the new coordinates we have a C^k differential equation of the form

$$\dot{y}_1 = a(y_1, \mu_1, \mu_2), \qquad \dot{y}_2 = y_2 b(y_1, y_2, \mu_1, \mu_2).$$

Since $p(\mu_2)$, defined in § 2, is C^k , we may assume that $p(\mu_2)$ is transformed to (0,0) for all μ_2 . In other words, $a(0,0,\mu_2) \equiv 0$. Since the stable manifold of $\dot{x} = f(x,0,0)$ at p is necessarily transformed into the line $y_1 = 0$, it is easy to arrange that

(7)
$$D_x y(p, 0, 0) u = (1, 0), \quad D_x y(p, 0, 0) v = (0, 1).$$

Taking into account assumptions (i)-(iii) and (vi), we have

(8)
$$\dot{y}_1 = \eta(\mu_2)y_1^2(1 + y_1g(y_1, \mu_2)) + \mu_1h(y_1, \mu_1, \mu_2), \\ \dot{y}_2 = -\lambda(y_1, \mu_1, \mu_2)y_2(1 + y_2k(y_1, y_2, \mu_1, \mu_2)),$$

with $\eta > 0$, h(0, 0, 0) > 0, $\lambda(0, 0, 0) = \lambda$.

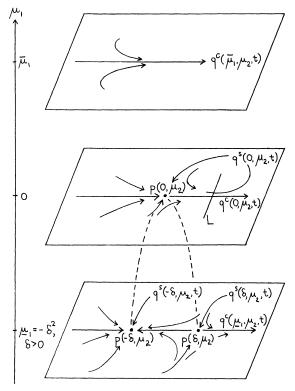


Fig. 3

LEMMA 1. There is a C^{k-2} mapping $p(\delta, \mu_2)$, defined for (δ, μ_2) near (0, 0), with values in \mathbb{R}^2 , such that p(0, 0) = p, and

$$p(\delta, \mu_2) \text{ is } \begin{cases} a \text{ saddle-node of } \dot{x} = f(x, 0, \mu_2) & \text{if } \delta = 0, \\ a \text{ saddle of } \dot{x} = f(x, -\delta^2, \mu_2) & \text{if } \delta > 0, \\ a \text{ sink of } \dot{x} = f(x, -\delta^2, \mu_2) & \text{if } \delta < 0. \end{cases}$$

Moreover, the mapping $(\delta, \mu_2) \rightarrow (p(\delta, \mu_2), -\delta^2, \mu_2)$ is a C^{k-2} diffeomorphism of a neighborhood of (0,0) in \mathbb{R}^2 onto a neighborhood of (p,0,0) in the set of equilibria of (5) near (0,0,0) (see Fig. 3).

We remark that $p(0, \mu_2)$ equals $p(\mu_2)$ defined in § 2.

Proof. The equilibria of the system

$$\begin{split} \dot{y}_1 &= \eta(\mu_2) y_1^2 (1 + y_1 g(y_1, \mu_2)) + \mu_1 h(y_1, \mu_1, \mu_2), \\ \dot{y}_2 &= -\lambda(y_1, \mu_1, \mu_2) y_2 (1 + y_2 k(y_1, y_2, \mu_1, \mu_2)), \\ \dot{\mu}_1 &= 0, \\ \dot{\mu}_2 &= 0 \end{split}$$

near (0, 0, 0, 0) comprise a set of the form

$$\{(y_1, 0, \mu_1, \mu_2): \mu_1 = \mu_1(y_1, \mu_2)\},\$$

where

$$\mu_1(0, \mu_2) = \frac{\partial \mu_1}{\partial y_1}(0, \mu_2) = 0$$
 and $\frac{\partial^2 \mu_1}{\partial y_1^2}(0, \mu_2) < 0$.

Therefore,

$$\mu_1 = -y_1^2(A(\mu_2) + y_1B(y_1, \mu_2)),$$

where $A(\mu_2) > 0$ and $A(\mu_2) + y_1 B(y_1, \mu_2)$ is C^{k-2} . For $\mu_1 \le 0$, let $\mu_1 = -\delta^2$. Then

(9)
$$\delta = y_1 (A(\mu_2) + y_1 B(y_1, \mu_2))^{1/2}.$$

(Taking into account the negative square root gives no additional information.) By the implicit function theorem, we can solve (9) for y_1 near $y_1 = 0$, $\delta = 0$, $\mu_2 = 0$. We obtain

(10)
$$y_1 = \hat{p}(\delta, \mu_2) \text{ with } \hat{p}(0, \mu_2) = 0, \frac{\partial \hat{p}}{\partial \delta}(0, \mu_2) > 0.$$

Then $\eta > 0$ and (10) imply that the equilibrium $(\hat{p}(\delta, \mu_2), 0)$ of (8) with $\mu_1 = -\delta^2$ is a saddle-node if $\delta = 0$, a saddle if $\delta > 0$, and a sink if $\delta < 0$.

Let

(11)
$$x = x(y, \mu_1, \mu_2)$$

be the change of coordinates inverse to (6). Define

(12)
$$p(\delta, \mu_2) = x((\hat{p}(\delta, \mu_2), 0), -\delta^2, \mu_2).$$

Then $p(\delta, \mu_2)$ satisfies the assertions of the lemma.

For future use, we note that

(13)
$$\frac{\partial p}{\partial \delta}(0,0) \text{ is a positive multiple of } u.$$

To see this, we compute from (12)

(14)
$$\frac{\partial p}{\partial \delta}(0,0) = D_y x((0,0),0,0) \left(\frac{\partial \hat{p}}{\partial \delta}(0,0),0\right).$$

Since $(\partial \hat{p}/\partial \delta)(0,0) > 0$ by (10), (14) is a positive multiple of u by (7).

System (8) with $\mu_1 = -\delta^2$ has at the equilibrium ($\hat{p}(\delta, \mu_2)$, 0) the invariant manifold $\{(y_1, y_2): y_1 = \hat{p}(\delta, \mu_2)\}$, a line. For $\delta = 0$ this line is the stable manifold of the saddle-node (0,0); for $\delta > 0$ it is the stable manifold of the saddle ($\hat{p}(\delta, \mu_2)$, 0); and for $\delta < 0$ it is the strong stable manifold of the sink ($\hat{p}(\delta, \mu_2)$, 0). These lines correspond to invariant manifolds of $\dot{x} = f(x, -\delta^2, \mu_2)$ at $p(\delta, \mu_2)$. Let $v(\delta, \mu_2) = D_y x((\hat{p}(\delta, \mu_2), 0), -\delta^2, \mu_2)(0, 1)$. Then $\dot{x} = f(x, -\delta^2, \mu_2)$ has at $p(\delta, \mu_2)$ an invariant curve tangent to $v(\delta, \mu_2)$ and these invariant curves vary in a C^{k-2} manner with (δ, μ_2). For $(\delta, \mu_2) = (0, 0)$, this invariant curve contains Γ . Now a construction similar to that of $q^c(\mu_1, \mu_2, t)$ yields a C^{k-2} family

$$q^s(\delta, \mu_2, t), \quad (\delta, \mu_2) \text{ small},$$

each a solution of $\dot{x} = f(x, -\delta^2, \mu_2)$, such that $q^s(\delta, \mu_2, t) \rightarrow p(\delta, \mu_2)$ as $t \rightarrow \infty$ along the negative $v(\delta, \mu_2)$ direction. Again we require $q^s(\delta, \mu_2, 0) \in L$; thus $q^s(0, 0, t) = q(t)$. Note that $q^s(0, \mu_2, t)$ is a branch of the stable manifold of the saddle-node $p(0, \mu_2)$ of $\dot{x} = f(x, 0, \mu_2)$; and if $\delta > 0$, $q^s(\delta, \mu_2, t)$ is a branch of the stable manifold of the saddle $p(\delta, \mu_2)$ of $\dot{x} = f(x, -\delta^2, \mu_2)$ (see Fig. 3).

For any vector $w = (w_1, w_2) \in \mathbb{R}^2$, let $w^{\perp} = (-w_2, w_1)$. Define $d^c(\mu_1, \mu_2)$ and $d^s(\delta, \mu_2)$ by

$$q^{c}(\mu_{1}, \mu_{2}, 0) = q(0) + [d^{c}(\mu_{1}, \mu_{2}) / ||f(q(0), 0, 0)||^{2}]f^{\perp}(q(0), 0, 0),$$

$$q^{s}(\delta, \mu_{2}, 0) = q(0) + [d^{s}(\delta, \mu_{2}) / ||f(q(0), 0, 0)||^{2}]f^{\perp}(q(0), 0, 0).$$

Then d^c is C^k and d^s is C^{k-2} . The number $d^c(\mu_1, \mu_2)$ (resp. $d^s(\delta, \mu_2)$) determines where on L the curve $q^c(\mu_1, \mu_2, t)$ (resp. $q^s(\delta, \mu_2, t)$) starts. We have

$$d^{c}(\mu_{1}, \mu_{2}) = f^{\perp}(q(0), 0, 0) \cdot [q^{c}(\mu_{1}, \mu_{2}, 0) - q(0)]$$
$$= f(q(0), 0, 0) \wedge [q^{c}(\mu_{1}, \mu_{2}, 0) - q(0)].$$

Similarly,

$$d^{s}(\delta, \mu_{2}) = f(q(0), 0, 0) \wedge [q^{s}(\delta, \mu_{2}, 0) - q(0)].$$

There is a separatrix loop of $\dot{x} = f(x, -\delta^2, \mu_2)$ through $p(\delta, \mu_2)$ if and only if $\delta \ge 0$ and

(15)
$$d(\delta, \mu_2) = d^c(-\delta^2, \mu_2) - d^s(\delta, \mu_2) = 0.$$

Here $d(\delta, \mu_2)$ is C^{k-2} .

We shall show that $(\partial d/\partial \delta)(0,0)$ is a negative multiple of $u \wedge v$ (hence is nonzero), and $(\partial d/\partial \mu_2)(0,0) = I$. Given these facts, the proof of Theorem 1 is completed as follows: if $I \neq 0$, then $\{(\delta, \mu_2): d(\delta, \mu_2) = 0\}$ is a C^{k-2} curve through (0,0) of the form

(16)
$$\delta = \ell \mu_2 + o(\mu_2), \qquad \ell = -I / \frac{\partial d}{\partial \delta}(0, 0).$$

Squaring both sides of (16) yields

$$\mu_1 = -\ell^2 \mu_2^2 + o(\mu_2^2).$$

The condition $\delta \ge 0$, applied to (16), shows that $\mu_2 = 0$ or ℓ and μ_2 have the same sign. But ℓ has the sign of $I \cdot (u \wedge v)$.

We now turn to the computation of $(\partial d/\partial \delta)(0,0)$ and $(\partial d/\partial \mu_2)(0,0)$. We shall need the following variational equations for $q^c(-\delta^2, \mu_2, t)$ and $q^s(\delta, \mu_2, t)$:

$$\frac{d}{dt} \frac{\partial q^c}{\partial \mu_2}(0,0,t) = D_x f(q(t),0,0) \frac{\partial q^c}{\partial \mu_2}(0,0,t) + \frac{\partial f}{\partial \mu_2}(q(t),0,0),$$

$$\frac{d}{dt} \frac{\partial q^s}{\partial \delta}(0,0,t) = D_x f(q(t),0,0) \frac{\partial q^s}{\partial \delta}(0,0,t),$$

$$\frac{d}{dt} \frac{\partial q^s}{\partial \mu_2}(0,0,t) = D_x f(q(t),0,0) \frac{\partial q^s}{\partial \mu_2}(0,0,t) + \frac{\partial f}{\partial \mu_2}(q(t),0,0).$$

As in [2], we define

$$\Delta_{\mu_2}^c(t) = f(q(t), 0, 0) \wedge \frac{\partial q^c}{\partial \mu_2}(0, 0, t)$$

and define $\Delta^s_{\delta}(t)$ and $\Delta^s_{\mu_2}(t)$ analogously.

For $d^c(-\delta^2, \mu_2)$ and $d^s(\delta, \mu_2)$ we have the derivative formulas

(17)
$$\frac{\partial d^{c}}{\partial \delta}(0,0) = \frac{\partial d^{c}}{\partial \mu_{1}}(0,0) \cdot \frac{d\mu_{1}}{d\delta}(0) = 0, \qquad \frac{\partial d^{c}}{\partial \mu_{2}}(0,0) = \Delta^{c}_{\mu_{2}}(0),$$

$$\frac{\partial d^{s}}{\partial \delta}(0,0) = \Delta^{s}_{\delta}(0), \qquad \frac{\partial d^{s}}{\partial \mu_{2}}(0,0) = \Delta^{s}_{\mu_{2}}(0).$$

Using the variational equations for q^c and q^s , we compute as in [2]:

(18)
$$\frac{d}{dt} \Delta_{\mu_2}^c(t) = \operatorname{div} f(q(t), 0, 0) \Delta_{\mu_2}^c(t) + f(q(t), 0, 0) \wedge \frac{\partial f}{\partial \mu_2}(q(t), 0, 0),$$

(19)
$$\frac{d}{dt}\Delta_{\delta}^{s}(t) = \operatorname{div} f(q(t), 0, 0)\Delta_{\delta}^{s}(t),$$

(20)
$$\frac{d}{dt} \Delta_{\mu_2}^s(t) = \operatorname{div} f(q(t), 0, 0) \Delta_{\mu_2}^s(t) + f(q(t), 0, 0) \wedge \frac{\partial f}{\partial \mu_2}(q(t), 0, 0).$$

Solving these linear differential equations, we obtain, for any t_1 ,

(21)
$$\Delta_{\mu_{2}}^{c}(0) = \Delta_{\mu_{2}}^{c}(t_{1}) \exp \int_{t_{1}}^{0} \operatorname{div} f(q(t), 0, 0) dt + \int_{t_{1}}^{0} \exp \left[-\int_{0}^{t} \operatorname{div} f(q(s), 0, 0) ds \right] \cdot f(q(t), 0, 0) \wedge \frac{\partial f}{\partial \mu_{2}}(q(t), 0, 0) dt,$$

$$(22) \qquad -\Delta_{\delta}^{s}(0) = -\Delta_{\delta}^{s}(t_{1}) \exp \left[-\int_{0}^{t_{1}} \operatorname{div} f(q(t), 0, 0) dt \right],$$

$$-\Delta_{\mu_{2}}^{s}(0) = -\Delta_{\mu_{2}}^{s}(t_{1}) \exp \left[-\int_{0}^{t_{1}} \operatorname{div} f(q(t), 0, 0) dt \right]$$

$$+\int_{0}^{t_{1}} \exp \left[-\int_{0}^{t} \operatorname{div} f(q(s), 0, 0) ds \right] \cdot f(q(t), 0, 0) \wedge \frac{\partial f}{\partial \mu_{2}}(q(t), 0, 0) dt.$$

We shall first evaluate (22) in the limit $t_1 \to \infty$. Using the definition of Δ_{δ}^s , we write

(24)
$$-\Delta_{\delta}^{s}(0) = \frac{\partial q^{s}}{\partial \delta}(0, 0, t_{1}) \wedge f(q(t_{1}), 0, 0) \exp\left[-\int_{0}^{t_{1}} \operatorname{div} f(q(s), 0, 0) ds\right].$$

LEMMA 2.

$$\lim_{t\to\infty}\frac{\partial q^s}{\partial \delta}(0,0,t)=\frac{\partial p}{\partial \delta}(0,0),\qquad \lim_{t\to\infty}\frac{\partial q^s}{\partial \mu_2}(0,0,t)=\frac{\partial p}{\partial \mu_2}(0,0).$$

Proof. We shall use the coordinates (6). Define

(25)
$$\tilde{q}^{s}(\delta, \mu_{2}, t) = y(q^{s}(\delta, \mu_{2}, t), -\delta^{2}, \mu_{2}) \\ = (\hat{p}(\delta, \mu_{2}), y_{2}(\delta, \mu_{2}, t)),$$

where $\hat{p}(\delta, \mu_2)$ is given by (10) and $y_2(\delta, \mu_2, t)$ is defined by (25). Since $q^s(\delta, \mu_2, t) \rightarrow p(\delta, \mu_2)$ as $t \rightarrow \infty$, for each (δ, μ_2) near (0, 0), $\tilde{q}^s(\delta, \mu_2, t)$ and hence $y_2(\delta, \mu_2, t)$ are defined for sufficiently large t. It follows from (7) that $y_2(\delta, \mu_2, t) > 0$ for large t. From (8), $y_2(\delta, \mu_2, t)$ satisfies a differential equation of the form

(26)
$$\frac{dz}{dt} = -\lambda(\delta, \mu_2)z(1+zG(z, \delta, \mu_2)),$$

where $\lambda(0,0) = \lambda$. Here λ and zG are C^{k-3} .

In order to prove the lemma, we shall study the asymptotic behavior of solutions of (26) as $t \to \infty$ by solving (26) by separation of variables. Let

(27)
$$z^{-1}[1+zG(z,\delta,\mu_2)]^{-1}=z^{-1}+H(z,\delta,\mu_2).$$

Then H is C^{k-4} . Fix $z_0 > 0$. Let

$$J(z, \delta, \mu_2) = \int_{z_0}^z H(s, \delta, \mu_2) \ ds.$$

Then J is C^{k-4} . Solving (26) by separation of variables using (27) yields

$$\ln z + J(z, \delta, \mu_2) = -\lambda(\delta, \mu_2)t + A(\delta, \mu_2),$$

or

$$z \exp J(z, \delta, \mu_2) = B(\delta, \mu_2) \exp(-\lambda(\delta, \mu_2)t).$$

Here $B(\delta, \mu_2)$ is determined by the value of $y_2(\delta, \mu_2, t)$ at some $t = t_0$. Hence B is C^{k-4} and B > 0. Since

$$\frac{\partial}{\partial z}[z \exp J(z, \delta, \mu_2)](0, \delta, \mu_2) \neq 0,$$

by the implicit function theorem we can solve the equation

$$z \exp J(z, \delta, \mu_2) = v$$

for z when z and v are near 0. We obtain

$$z = K(v, \delta, \mu_2),$$

where K is C^{k-4} and

(28)
$$K(0, \delta, \mu_2) \equiv 0.$$

Putting $z = y_2$ and $v = B(\delta, \mu_2) \exp(-\lambda(\delta, \mu_2)t)$, we obtain

(29)
$$y_2 = K(B(\delta, \mu_2) \exp(-\lambda(\delta, \mu_2)t), \delta, \mu_2).$$

From (29) and (28) it follows that $(\partial y_2/\partial \delta)(\delta, \mu_2, t)$ and $(\partial y_2/\partial \mu_2)(\delta, \mu_2, t)$ approach 0 as $t \to \infty$. Therefore (25) implies that as $t \to \infty$,

$$(30) \qquad \frac{\partial \tilde{q}^{s}}{\partial \delta}(\delta, \mu_{2}, t) \rightarrow \left(\frac{\partial \hat{p}}{\partial \delta}(\delta, \mu_{2}), 0\right), \qquad \frac{\partial \tilde{q}^{s}}{\partial \mu_{2}}(\delta, \mu_{2}, t) \rightarrow \left(\frac{\partial \hat{p}}{\partial \mu_{2}}(\delta, \mu_{2}), 0\right).$$

Now

$$\frac{\partial q^s}{\partial \delta}(\delta, \mu_2, t) = \frac{\partial}{\partial \delta} x(\tilde{q}^s(\delta, \mu_2, t), -\delta^2, \mu_2),$$

where $x(y, \mu_1, \mu_2)$ is given by (11). Therefore

$$\frac{\partial q^s}{\partial \delta}(0,0,t) = D_y x(\tilde{q}^s(0,0,t),0,0) \frac{\partial \tilde{q}^s}{\partial \delta}(0,0,t).$$

By (30) and (14),

$$\lim_{t \to \infty} \frac{\partial q^{s}}{\partial \delta}(0, 0, t) = D_{y}x((0, 0), 0, 0) \left(\frac{\partial \hat{p}}{\partial \delta}(0, 0), 0\right)$$
$$= \frac{\partial p}{\partial \delta}(0, 0).$$

Similarly, the second formula of the lemma follows from (30) and the following formula derived from (12):

$$\frac{\partial p}{\partial \mu_2}(0,0) = D_y x((0,0),0,0) \left(\frac{\partial \hat{p}}{\partial \mu_2}(0,0),0 \right) + \frac{\partial x}{\partial \mu_2}((0,0),0,0). \quad \Box$$

We shall now use the computations we have done in proving Lemma 2 to study the other terms of (24). By (25),

(31)
$$\dot{q}^{s}(\delta, \mu_{2}, t) = D_{x}y(q^{s}(\delta, \mu_{2}, t), -\delta^{2}, \mu_{2})\dot{q}^{s}(\delta, \mu_{2}, t).$$

Let $\delta = \mu_2 = 0$ in (31). Since $\dot{q}^s(0, 0, t) = \dot{q}(t) = f(q(t), 0, 0)$, we obtain

$$f(q(t), 0, 0) = [D_x y(q(t), 0, 0)]^{-1} \ddot{q}^s(0, 0, t).$$

It follows easily from (25), (28) and (29) that

(32)
$$q(t) = p + \mathcal{O}(\exp(-\lambda t)),$$
$$\tilde{q}^{s}(t) = (0, -C \exp(-\lambda t) + \rho(\exp(-\lambda t))),$$

where C > 0. Therefore, setting $t = t_1$,

(33)
$$f(q(t_1), 0, 0) = \{ [D_x y(p, 0, 0)]^{-1} + \mathcal{O}(\exp(-\lambda t_1)) \} \cdot (0, -C \exp(-\lambda t_1) + \rho(\exp(-\lambda t_1))).$$

From (32) we also have

$$\operatorname{div} f(q(t), 0, 0) = -\lambda + \mathcal{O}(\exp(-\lambda t)).$$

Therefore

(34)
$$\exp\left[-\int_0^{t_1} \operatorname{div} f(q(s), 0, 0) \ ds\right] = \exp\left(\lambda t_1\right) \cdot \exp\left[\int_0^{t_1} \mathcal{O}(\exp\left(-\lambda s\right)) \ ds.$$

Then (33) and (34) give

(35)
$$\lim_{t_1 \to \infty} f(q(t_1), 0, 0) \exp \left[-\int_0^{t_1} \operatorname{div} f(q(s), 0, 0) \, ds \right] \\ = \left[D_x y(p, 0, 0) \right]^{-1} \cdot \left(0, -C \exp \int_0^{\infty} \mathcal{O}(\exp(-\lambda s)) \, ds \right),$$

where the integral clearly converges.

Now (24), Lemma 2 and (35) imply that

$$(36) \qquad -\Delta_{\delta}^{s}(0) = \frac{\partial p}{\partial \delta}(0,0) \wedge \lim_{t_1 \to \infty} f(q(t_1),0,0) \exp\left[-\int_{0}^{t_1} \operatorname{div} f(q(s),0,0) \, ds\right],$$

where the limit exists.

Notice that (7) implies that (35) is a negative multiple of v. Then (13) implies that $-\Delta_{\delta}^{s}(0)$ is a negative multiple of $u \wedge v$. By (15) and (17), $(\partial d/\partial \delta)(0,0)$ is also a negative multiple of $u \wedge v$.

We now turn to (23). We claim that

$$-\Delta_{\mu_{2}}^{s}(0) = \frac{\partial p}{\partial \mu_{2}}(0,0) \wedge \lim_{t_{1} \to \infty} f(q(t_{1}),0,0) \exp\left[-\int_{0}^{t_{1}} \operatorname{div} f(q(s),0,0) \, ds\right]$$

$$+ \int_{0}^{\infty} \exp\left[-\int_{0}^{t} \operatorname{div} f(q(s),0,0) \, ds\right] f(q(t),0,0)$$

$$\wedge \frac{\partial f}{\partial \mu_{2}}(q(t),0,0) \, dt.$$

The proof is modeled on that of (36). Using Lemma 2, we first show that the first summand of (23) approaches, as $t_1 \rightarrow \infty$, the first summand of (37), where the limit

exists. Then, since $-\Delta_{\mu_2}^s(0)$ is finite, the second summand of (23), the integral, must approach a limit as $t_1 \to \infty$.

Finally we turn to (21).

LEMMA 3.

$$\lim_{t\to-\infty}\frac{\partial q^c}{\partial \mu_2}(0,0,t)=0.$$

Proof. Again we shall use the coordinates (6). Define

$$\tilde{q}^c(\mu_2, t) = y(q^c(0, \mu_2, t), 0, \mu_2) = (y_1(\mu_2, t), 0).$$

Since $q^c(0, \mu_2, t) \to p(0, \mu_2)$ as $t \to -\infty$, for each μ_2 near 0, $\tilde{q}^c(\mu_2, t)$, and hence $y_1(\mu_2, t)$, is defined for sufficiently negative t. From (7), $y_1(\mu_2, t) > 0$ for sufficiently negative t. From its definition, $y_1(\mu_2, t)$ satisfies a differential equation of the form

(38)
$$\frac{dz}{dt} = \eta(\mu_2)z^2(1+zG(z,\mu_2)).$$

Here η and zG are C^{k-2} .

Let

(39)
$$z^{-2}(1+zG(z,\mu_2))^{-1}=z^{-2}+A(\mu_2)z^{-1}+H(z,\mu_2).$$

Here A is C^{k-3} and H is C^{k-4} . Fix $z_0 > 0$. Let

$$J(z, \mu_2) = \int_{z_0}^z H(s, \mu_2) ds.$$

J is C^{k-4} . Then solving (38) by separation of variables using (39) yields

$$-z^{-1} + A(\mu_2) \ln z + J(z, \mu_2) = \eta(\mu_2)t + B(\mu_2).$$

Here $B(\mu_2)$ is determined by the value of $y_1(\mu_2, t)$ at some $t = t_0$. Hence B is C^{k-4} . Rearranging yields

(40)
$$z[1-A(\mu_2)z \ln z - zJ(z,\mu_2)]^{-1} = -[\eta(\mu_2)t + B(\mu_2)]^{-1}.$$

Let $\Phi(z, \mu_2)$ equal the left-hand side of (40). $\Phi(z, \mu_2)$ is a C^1 function of z and μ_2 on a neighborhood of (0,0) in $\{(z, \mu_2): z \ge 0\}$; $\Phi(0, \mu_2) = 0$, and $(\partial/\partial z)\Phi(0, \mu_2) = 1$. By the implicit function theorem, we can solve the equation $\Phi(z, \mu_2) = v$ for z when z and v are near 0, $v \ge 0$ (in which case $z \ge 0$). We obtain

$$z=v+R(v,\mu_2),$$

where R is C^1 , $R(0, \mu_2) \equiv 0$ and R is o(v). Putting $z = y_1$ and $v = -[\eta(\mu_2)t + B(\mu_2)]^{-1}$, t large negative, we obtain

$$y_1 = -[\eta(\mu_2)t + B(\mu_2)]^{-1} + R(-[\eta(\mu_2)t + B(\mu_2)]^{-1}, \mu_2).$$

It follows that $(\partial y_1/\partial \mu_2)(\mu_2, t) \to 0$ as $t \to -\infty$, so

(41)
$$\frac{\partial \tilde{q}^c}{\partial \mu_2}(\mu_2, t) \to 0 \quad \text{as } t \to -\infty.$$

Lemma 3 follows from (41) the way Lemma 2 follows from (30). \Box

To evaluate (21) in the limit $t_1 \to -\infty$, we use the definition of $\Delta_{\mu_2}^c$ to write (21) as

(42)
$$\Delta_{\mu_{2}}^{c}(0) = -\frac{\partial q^{c}}{\partial \mu_{2}}(0, 0, t_{1}) \wedge f(q(t), 0, 0) \exp \int_{t_{1}}^{0} \operatorname{div} f(q(s), 0, 0) ds + \int_{t_{1}}^{0} \exp \left[-\int_{0}^{t} \operatorname{div} f(q(s), 0, 0) ds\right] f(q(t), 0, 0) \wedge \frac{\partial f}{\partial \mu_{2}}(q(t), 0, 0) dt.$$

Of course

(43)
$$f(q(t), 0, 0) \to 0 \text{ as } t_1 \to -\infty$$

Moreover,

(44)
$$\operatorname{div} f(q(t), 0, 0) = -\lambda + \mathcal{O}([\eta(0)t + c(0)]^{-1}) \text{ as } t \to -\infty.$$

By Lemma 3, (43) and (44), we have

$$\lim_{t_1\to-\infty} -\frac{\partial q^c}{\partial \mu_2}(0,0,t_1) \wedge f(q(t_1),0,0) \exp \int_{t_1}^0 \operatorname{div} f(q(s),0,0) \ ds = 0.$$

Therefore the second summand of (42), the integral, approaches a limit as $t_1 \to -\infty$, so

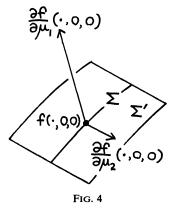
(45)
$$\Delta_{\mu_2}^c(0) = \int_{-\infty}^0 \exp\left[-\int_0^t \operatorname{div} f(q(s), 0, 0) \, ds\right] f(q(t), 0, 0) \wedge \frac{\partial f}{\partial \mu_2}(q(t), 0, 0) \, dt.$$

Finally, we complete the proof of Theorem 1 by calculating

$$\frac{\partial d}{\partial \mu_2}(0,0) = \frac{\partial d^c}{\partial \mu_2}(0,0) - \frac{\partial d^s}{\partial \mu_2}(0,0) = \Delta^c_{\mu_2}(0) - \Delta^s_{\mu_2}(0) = (45) + (37) = I.$$

4. Proof of Theorem 2. First we show that Σ is a C^k codimension one submanifold of Σ' . Let $f \in \Sigma$ with saddle-node p and let L be a line segment perpendicular to the separatrix loop Γ as in § 3. For g near f in Σ' , there is a unique saddle-node p_g near p; p_g is a C^k function of $g \in \Sigma'$. The stable and center manifolds of p_g also depend C^k on g [1, § 9.2]. Thus their intersections with L are C^k functions of g. Therefore the function $d(0, \mu_2)$ from § 2 extends to a C^k function d(g) defined for $g \in \Sigma'$ near f; d(g) measures the separation of these points of intersection. (We remark that $d(0, \mu_2)$ is C^k although $d(\delta, \mu_2)$ is only C^{k-2} .) Then d(g) = 0 if and only if $g \in \Sigma$. Since it is easy to find a perturbation $f + \varepsilon h$ in Σ' such that $d/d\varepsilon|_{\varepsilon=0}d(f+\varepsilon h) \neq 0$, Σ is a C^k codimension one submanifold of Σ' .

To prove Theorem 2, it suffices to show that $f(\cdot, \mu_1, \mu_2)$ is transverse to Σ at $(\mu_1, \mu_2) = (0, 0)$ if and only if $I \neq 0$. Since $f(\cdot, 0, \mu_2) \in \Sigma'$ for all small μ_2 , $(\partial f/\partial \mu_2)(\cdot, 0, 0)$ is tangent to Σ' (see Fig. 4). But $(\partial \tilde{f}/\partial \nu_1)(\cdot, 0, 0) = (\partial f/\partial \mu_1)(\cdot, 0, 0) - (\partial f/\partial \mu_2)(\cdot, 0, 0)\alpha'(0)$. Since $(\partial \tilde{f}/\partial \nu_1)(\cdot, 0, 0)$ is transverse to Σ' by



assumption (vi) (see [7]), and $(\partial f/\partial \mu_2)(\cdot, 0, 0)$ is tangent to Σ' , $(\partial f/\partial \mu_1)(\cdot, 0, 0)$ is transverse to Σ' . Thus we need to show that $(\partial f/\partial \mu_2)(\cdot, 0, 0)$ is transverse to Σ if and only if $I \neq 0$. But this follows from the formula $(\partial d/\partial \mu_2)(0, 0) = I$.

5. The pendulum equation. We consider the differential equation for the damped pendulum with constant applied torque, in the dimensionless form studied in [5]:

$$\beta\ddot{\phi} + \dot{\phi} + \sin \phi = \rho.$$

Putting $y = \dot{\phi}$, we have

(46)
$$\dot{\phi} = y, \qquad \dot{y} = \frac{1}{\beta} \left(-y - \sin \phi + \rho \right).$$

We identitify ϕ and $\phi + 2\pi$, so that (46) defines a vector field on a cylinder; ρ and β are parameters, which we shall assume positive. We remark that Theorem 1 applies equally well to vector fields on the compact set $\{(\phi, y): |y| \le d\}$.

Let $x = (\phi, y)$,

$$f(x,\beta,\rho) = (f_1((\phi,y),\beta,\rho), f_2((\phi,y),\beta,\rho)) = \left(y,\frac{1}{\beta}(-y-\sin\phi+\rho)\right).$$

If $\rho < 1$, $\dot{x} = f(x, \rho, \beta)$ has two equilibria, one a saddle and one a sink; $\dot{x} = f(x, 1, \beta)$ has one equilibrium, at $(\pi/2, 0)$ independent of β ; if $\rho > 1$, $\dot{x} = f(x, \rho, \beta)$ has no equilibria. We note that

$$D_{x}f\left(\left(\frac{\pi}{2},0\right),1,\beta\right)=\begin{bmatrix}0&1\\0&-1/\beta\end{bmatrix},$$

which has eigenvalues $0, -1/\beta$. Corresponding right eigenvectors are u = (1, 0) and $v = (-\beta, 1)$; a left eigenvector for the eigenvalue 0 is $(1, \beta)$.

To show that $\dot{x} = f(x, 1, \beta)$ has a saddle-node at $(\pi/2, 0)$, we compute (assumption (iii)):

$$(1,\beta) \cdot D_x^2 f\left(\left(\frac{\pi}{2},0\right),1,\beta\right) \cdot ((1,0),(1,0)) = (1,\beta) \cdot \left(0,\frac{\partial^2 f}{\partial \phi^2}\left(\left(\frac{\pi}{2},0\right),1,\beta\right)\right)$$
$$= (1,\beta) \cdot \left(0,\frac{1}{\beta}\right) = 1.$$

We also note (assumption (iv)):

$$(1, \beta) \cdot D_{\rho} f\left(\left(\frac{\pi}{2}, 0\right), 1, \beta\right) = (1, \beta) \cdot \left(0, \frac{1}{\beta}\right) = 1.$$

It is shown in [3] that there is a unique positive β_0 such that $\dot{x} = f(x, 1, \beta_0)$ has a separatrix loop at the saddle-node $(\pi/2, 0)$. The separatrix loop, considered as a curve in ϕy -space with ϕ and $\phi + 2\pi$ not yet identified, can be expressed as $y = y(\phi)$, $\pi/2 < \phi < 5\pi/2$; $y(\phi) > 0$ for all ϕ . As $\phi \to \pi/2$, $y \to 0$ and $y/(\phi - \pi/2) \to 0$; as $\phi \to 5\pi/2$, $y \to 0$ and $y/(\phi - 5\pi/2) \to -1/\beta$) (see Fig. 5).

Thus assumptions (i)-(vi) are verified if we put $\tilde{f} = f$, $\nu_1 = \rho - 1$, $\nu_2 = \beta - \beta_0$. The change of variables from ν to μ is not necessary here, i.e., we may put $\mu_1 = \nu_1$, $\mu_2 = \nu_2$. For simplicity we shall continue to use the parameters ρ and β .

To compute I, we first note that the saddle-node $p(\beta)$ of $\dot{x} = f(x, 1, \beta)$ is identically $(\pi/2, 0)$, so that the first summand of I is 0. Now div $f(x, \rho, \beta) = -1/\beta$, and

$$f((\phi, y), \rho, \beta) \wedge \frac{\partial f}{\partial \beta}((\phi, y), \rho, \beta) = -\frac{1}{\beta^2}y(-y - \sin \phi + \beta) = -\frac{1}{\beta}y\dot{y}.$$

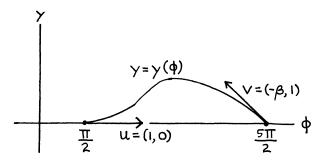


Fig. 5

Therefore,

$$\begin{split} I &= -\frac{1}{\beta_0} \int_{-\infty}^{\infty} e^{t/\beta_0} y \dot{y} \, dt \\ &= \lim_{S, T \to \infty} \left\{ -\frac{1}{\beta_0} e^{t/\beta_0} \cdot \frac{y^2}{2} \bigg|_{-S}^{T} + \frac{1}{\beta_0^2} \int_{-S}^{T} e^{t/\beta_0} \cdot \frac{y^2}{2} \, dt \right\}. \end{split}$$

From (29), as $t \to \infty$, $(\phi(t) - (5\pi/2), y(t)) = c e^{-t/\beta_0} (-\beta_0, 1) + o(e^{-t/\beta_0})$ for some positive constant c. Therefore $[y(t)]^2 = c^2 e^{-2t/\beta_0} + o(e^{-2t/\beta_0})$ as $t \to \infty$. Hence

$$\lim_{T\to\infty} -\frac{1}{\beta_0} e^{T/\beta_0} \cdot \frac{[y(T)]^2}{2} = 0.$$

Of course,

$$\lim_{S\to\infty} -\frac{1}{\beta_0} e^{-S/\beta_0} \cdot \frac{[y(-S)]^2}{2} = 0,$$

since $y(t) \to 0$ as $t \to -\infty$. Therefore

$$I = \frac{1}{\beta_0^2} \int_{-\infty}^{\infty} e^{t/\beta_0} \cdot \frac{y^2}{2} dt > 0.$$

Since $u \wedge v > 0$, Theorem 1 implies that for (ρ, β) near $(1, \beta_0)$, $\dot{x} = f(x, \rho, \beta)$ has a separatrix loop if and only if $\rho = 1 - \ell^2 (\beta - \beta_0)^2 + \cdots$, with $\ell \neq 0$, and $\beta - \beta_0 \ge 0$ (see Fig. 6). This result is in agreement with statements in [5].

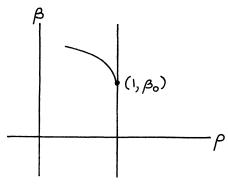


Fig. 6

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