

Normal Forms of Vector Fields

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July 5, 2017

1 Introduction and notations

As a well-known fact, a vector field $X \in \Gamma(TM)$ non-vanishing at a point x_0 can be represented as $\frac{\partial}{\partial s_1}$ for some coordinate chart $(U, \phi = (s_1, s_2, \dots, s_n))$ near that point, which means that X can be sent on a locally constant vector field through some diffeomorphism. However, there is not such a good local representation around singular points. The aim of this presentation is to show a relatively good local result formulated in [2], and to show its concrete applications.

Since the problem is local, we may well consider vector fields defined on a neighbourhood of $x_0 = 0 \in \mathbb{R}^n$ and vanishing at 0. We denote this set by \mathcal{X} . Note that \mathcal{X} is preserved by the Lie bracket, so $(\mathcal{X}, [\])$ is a Lie algebra. We also define \mathcal{G} as the group of local diffeomorphisms of \mathbb{R}^n that fix 0. We shall later study the action of \mathcal{G} on \mathcal{X} .

We will frequently use the natural identification $T\mathbb{R}^n \approx \mathbb{R}^{n+n}$ without mentioning it. Under this identification, we can also see $X \in \mathcal{X}$ as a function $X : U \ni 0 \rightarrow \mathbb{R}^n$. For any function $\phi : U \ni 0 \rightarrow \mathbb{R}^n$, we denote $[\phi]^k$ the up to k -th Taylor development of ϕ at 0 and $[\phi]_l$ the l -th homogenous term of its Taylor development at 0. If \mathcal{Y} is a subalgebra of \mathcal{X} , we denote $\mathcal{Y}^k := \{[Y]^k, Y \in \mathcal{Y}\}$ and $\mathcal{Y}_k := \{[Y]_k, Y \in \mathcal{Y}\}$. In the following text, we always consider subalgebras \mathcal{Y} that satisfy an additional condition:

$$\forall k \in \mathbb{N}, \mathcal{Y}_k \subseteq \mathcal{Y},$$

and we call them graded subalgebras as a reminder.

Finally, for $X \in \mathcal{X}$, as $A = [X]^1 = [X]_1$ is the linear part of X at 0, we see it as a matrix of $\mathcal{M}(n, \mathbb{R})$.

2 Statement of the main Theorem

Recall the Jordan-Dunford decomposition of matrix: for any $A \in \mathcal{M}(n, \mathbb{R})$, $\exists! A_s \in \mathcal{M}(n, \mathbb{R}), A_n \in \mathcal{M}(n, \mathbb{R})$ such that $A = A_s + A_n$ where A_s is diagonalisable, A_n is nilpotent, and $[A_s, A_n] = 0$.

Theorem 1. *Let \mathcal{Y} be a graded subalgebra of \mathcal{X} , $X \in \mathcal{Y}$ a vector field with $A = [X]_1$. We consider a matrix B which satisfies one of the two conditions:*

- $B = A_s$.
- $B \in \mathcal{Y}_1$ and B is adjoint to A under some inner product over \mathbb{R}^n .

*Then there exists a sequence $Y_i \in \mathcal{Y}_i$, $i \geq 2$, such that for any $k \geq 2$, $[(\varphi_{Y_2+\dots+Y_k}^1)_*X - A]^k$ commutes with B , where $\varphi_{Y_2+\dots+Y_k}^t$ is the flow of $Y_2 + \dots + Y_k$ at time t .*

Remark 1. *We shall see that the field $(\phi_{Y_2+\dots+Y_k}^1)_*X$ is a member of \mathcal{Y} .*

The vector field $(\varphi_{Y_2+\dots+Y_k}^1)_*X$ is called the *normal form* of X . We will see that its linear part is still A , but the next terms in its Taylor development have an additional property given by the commutation with B . In the best cases, all the monomial terms of order 2 to k are zero, which means that X , once pushed forward by the diffeomorphism $\varphi_{Y_2+\dots+Y_k}^1$, is very close to its linear part. Then the dynamics of the flow of X can be approached by the one of A . In the worst cases, e.g. if $B = 0$, the theorem gives no information at all.

3 Proof of the main theorem

We grade the ring $\mathbb{R}[x_1, x_2, \dots, x_n]$ by degree, *i.e.*

$$\mathbb{R}[x_1, x_2, \dots, x_n] = \bigoplus_{k \geq 0} E_k,$$

where E_k is the space of homogenous polynomials of degree k . We define $E^k := E_1 \oplus \dots \oplus E_k$ the vector space of polynomials P of degree lower than k such that $P(0) = 0$.

The main idea in this proof is to associate both vector fields $X \in \mathcal{X}$ and diffeomorphisms $\phi \in \mathcal{G}$ with endomorphisms of E^k , \widehat{X} and $\widetilde{\phi}$, which express their behaviour around 0 at order k . We therefore define two morphisms of vector spaces:

$$\begin{aligned} \widehat{\cdot} : \mathcal{X} &\longrightarrow \text{End}(E^k) \\ X &\longmapsto (P \in E^k \mapsto [X \cdot P]^k), \end{aligned}$$

and

$$\begin{aligned} \widetilde{\cdot} : \mathcal{G} &\longrightarrow \text{End}(E^k) \\ \phi &\longmapsto (P \in E^k \mapsto [P \circ \phi]^k). \end{aligned}$$

3.1 Analysis of the kernels

If $[X]^k = 0$, $X \in \mathcal{X}$, then

$$[XP]^k = \left[\sum_{i=1}^n X^i \partial_i P \right]^k = \sum_{i=1}^n [X^i \partial_i P]^k = \sum_{i=1}^n [[X^i]^k \partial_i P]^k = [[X]^k P]^k = 0.$$

In the same way, if $[\phi]^k = 0$, $\phi \in \mathcal{G}$, then

$$\begin{aligned} [P \circ \phi]^k(x) &= [P \circ ([\phi]^k(x) + O(|x|^{k+1}))]^k = [P([\phi]^k(x) + O(|x|^{k+1}))]^k \\ &= [P \circ [\phi]^k]^k(x) = 0, \end{aligned}$$

hence $[P \circ \phi]^k = 0$. Moreover, when $P = \text{id}$, $\widehat{X}(\text{id}) = [X]^k$ and $\widetilde{\phi}(\text{id}) = [\phi]^k$. Thus we have $\text{Ker } \widehat{\cdot} = \{X \in \mathcal{X}, [X]^k = 0\}$ and $\text{Ker } \widetilde{\cdot} = \{\phi \in \mathcal{G}, [\phi]^k = 0\}$. We then conclude that the image of $\widehat{\cdot}$ and $\widetilde{\cdot}$ are determined by the k -th order development of their preimages.

3.2 Algebraic properties of $\widehat{\cdot}$ and $\widetilde{\cdot}$

Proposition 2. *We have:*

1. $\forall \phi, \psi \in \mathcal{G} \quad (\phi \circ \psi)^\sim = \widetilde{\psi} \circ \widetilde{\phi}$
2. $\forall \phi \in \mathcal{G}, \forall X \in \mathcal{X}, \phi_* X \in \mathcal{X}$ and $(\phi_* X)^\widehat{=} \widetilde{\phi}^{-1} \circ \widehat{X} \circ \widetilde{\phi}$
3. $[X, Y]^\widehat{=} \widehat{X} \circ \widehat{Y} - \widehat{Y} \circ \widehat{X}$

Proof. 1. In fact,

$$\begin{aligned} (\phi \circ \psi)(x) &= [\phi]^k(\psi(x)) + O(|\psi(x)|^{k+1}) \\ &= [\phi]^k([\psi]^k(x) + O(|x|^{k+1})) + O(|x|^{k+1}) \\ &= [\phi]^k \circ [\psi]^k(x) + O(|x|^{k+1}), \end{aligned}$$

hence for all P in E_k ,

$$(\phi \circ \psi)\tilde{\phi}(P) = [P \circ [\phi \circ \psi]^k]^k = [P \circ [\phi]^k \circ [\psi]^k]^k = [[P \circ [\phi]^k]^k \circ [\psi]^k]^k = \tilde{\psi}(\tilde{\phi}(P)).$$

2. Take $P \in E_k$, we have $(\phi_*X)P = X(P \circ \phi) \circ \phi^{-1}$, so

$$\begin{aligned} [(\phi_*X)P]^k &= \widetilde{\phi^{-1}} \left([X(P \circ \phi)]^k \right) = \widetilde{\phi^{-1}} \left(\sum_{i=1}^n [X^i \partial_i (P \circ \phi)]^k \right) \\ &= \widetilde{\phi^{-1}} \left(\sum_{i=1}^n [X^i L_i ([P \circ \phi]^{k+1})]^k \right), \end{aligned}$$

where $L_i : x_1^{j_1} \dots x_i^{j_i} \dots x_n^{j_n} \mapsto j_i x_1^{j_1} \dots x_i^{j_i-1} \dots x_n^{j_n}$. Since $\forall i, X^i(0) = 0$,

$$\begin{aligned} [(\phi_*X)P]^k &= \widetilde{\phi^{-1}} \left(\sum_{i=1}^n \left([X^i \partial_i [P \circ \phi]^k]^k + [X^i \partial_i [P \circ \phi]_{k+1}]^k \right) \right) \\ &= \widetilde{\phi^{-1}} \left([X[P \circ \phi]^k]^k \right) = \widetilde{\phi^{-1}} \circ \widehat{X} \circ \tilde{\phi}(P). \end{aligned}$$

3. By the same argument as in 2, as $X(0) = 0$,

$$\begin{aligned} \widehat{XY}(P) &= [X(YP)]^k = [X[YP]^{k+1}]^k = [X[YP]^k]^k + [X[YP]_{k+1}]^k \\ &= [X[YP]^k]^k = \widehat{X} \circ \widehat{Y}(P). \end{aligned}$$

□

3.3 Image of the flow of a vector field

For $X \in \mathcal{X}$, we now consider its flow $\varphi_X : \mathcal{D} \subseteq \mathbb{R} \times \mathbb{R}^n \rightarrow \mathbb{R}^n$. Since $\varphi_X(t, 0) = 0$ is defined for all $t \in \mathbb{R}$, we have $\varphi_X^1 = \varphi_X(1, \cdot)$ is defined on a neighbourhood of 0, hence $\varphi_X^1 \in \mathcal{G}$.

Proposition 3. We have $\widetilde{\varphi}_X^1 = \exp(\widehat{X})$.

Proof. Note that $\text{End}(E^k)$ is of finite dimension (with any norm on $\text{End}(E^k)$), thus the ODE $\alpha'(t) = \alpha(t) \circ \widehat{X}$ with initial condition $\alpha(0) = \mathbb{1}_{E^k}$ has a unique solution. The solution is exactly $\alpha(t) = \exp(t\widehat{X})$.

We can also note that $\partial_t(P \circ \varphi_X) = (XP) \circ \varphi_X$. We consider the Taylor developments of $\partial_t(P \circ \varphi_X)$ and $(P \circ \varphi_X)$ at 0 with respect to x for each t :

$$\partial_t(P \circ \varphi_X)(t, x) = \sum_{|\nu| \leq k} \frac{1}{\nu!} \theta^\nu(t) x^\nu + O(|x|^{k+1})$$

and

$$(P \circ \varphi_X)(t, x) = \sum_{|\nu| \leq k} \frac{1}{\nu!} \beta^\nu(t) x^\nu + O(|x|^{k+1}),$$

where $\nu \in \mathbb{N}^n$ are multi-indices, and $\nu! = \prod_{i=1}^n \nu_i!$.

Then $(\beta^\nu)'(t) = \theta^\nu(t)$, for $\beta^\nu(t) = \partial_\nu(P \circ \varphi_X)(t, 0)$ and

$$\theta^\nu(t) = \partial_\nu(\partial_t(P \circ \varphi_X))(t, 0) = \partial_t(\partial_\nu(P \circ \varphi_X))(t, 0) = \frac{d}{dt} \beta^\nu(t)$$

Hence we have $\partial_t \widetilde{\varphi}_X^t(P) = \widetilde{\varphi}_X^t \circ \widehat{X}(P)$, and $\varphi_X^0(P) = P$, so $\widetilde{\varphi}_X^t = \exp(t\widehat{X})$. \square

3.4 Conjugated operators

We define two more operators. For $V \in \text{End}(E^k)$, if V is invertible, we define

$$\begin{aligned} Ad_V &: \text{End}(E^k) \longrightarrow \text{End}(E^k) \\ M &\longmapsto V^{-1} \circ M \circ V. \end{aligned}$$

and for any $L \in \text{End}(E^k)$,

$$\begin{aligned} ad_L &: \text{End}(E^k) \longrightarrow \text{End}(E^k) \\ M &\longmapsto [M, L] = M \circ L - L \circ M. \end{aligned}$$

Note that $Ad_V \in GL(\text{End}(E^k)) \subseteq \text{End}(\text{End}(E^k))$, and that $ad_L \in \text{End}(\text{End}(E^k))$.

Proposition 4. $\forall L \in \text{End}(E^k)$, $Ad_{\exp L} = \exp(ad_L)$.

Proof. Let $\alpha(t) = Ad_{\exp(tL)}$, then

$$\alpha'(t)M = -L \circ (\alpha(t)M) + (\alpha(t)M) \circ L = (ad_L \circ \alpha(t))M.$$

Thus $\alpha(t) = \exp(ad_{tL})$, so $Ad_{\exp L} = \exp(ad_L)$. \square

3.5 Proof of the theorem

The commutative statement of the k -th development of vector fields with the field B is now well characterised by the commutativity of their image under $\widehat{\cdot}$ by the third assumption of proposition 2. We'll now prove the theorem by induction, with the help of $\widehat{\cdot}$ and $\widetilde{\cdot}$. Suppose we already have $Y_i \in \mathcal{Y}_i$, $2 \leq i \leq k-1$, such that $[(\phi_Z^1)_*X - A]^{k-1}$ commutes with B and $(\phi_Z^1)_*X \in \mathcal{Y}$, where $Z = \sum_{2 \leq i \leq k-1} Y_i$.

By the second assumption of proposition 2 and propositions 3 and 4, we have

$$\begin{aligned} ((\phi_Z^1)_*X)^\widehat{} &= (\widetilde{(\phi_Z^1)^{-1}}) \circ \widehat{X} \circ \widetilde{\varphi_X^1} = \exp(-\widehat{Z}) \circ \widehat{X} \circ \exp(\widehat{Z}) \\ &= Ad_{\exp \widehat{Z}}(\widehat{X}) = \exp(ad_{\widehat{Z}})\widehat{X} \end{aligned}$$

Indeed, $\forall Y \in \mathcal{Y}^k$, $[Z, Y] \in \mathcal{Y}^k$ so $ad_{\widehat{Z}}(\widehat{Y}) = [\widehat{Y}, \widehat{Z}] = [Y, Z]^\widehat{}$, and from the kernel of $\widehat{\cdot}$, we know that $(\mathcal{Y}^k)^\widehat{} = \widehat{\mathcal{Y}} \subseteq End(E^k)$. Then $ad_{\widehat{Z}}$ preserves $\widehat{\mathcal{Y}} \subseteq End(E^k)$, and so does $\exp(ad_{\widehat{Z}}) = Ad_{\exp \widehat{Z}}$.

By the induction hypotheses and the preservation argument for dimension $k-1$, we have $((\phi_Z^1)_*X)^\widehat{} = Ad_{\exp \widehat{Z}}(\widehat{X}) = \widehat{A} + \widehat{N} + \widehat{R}_k$, where $N \in \mathcal{Y}^{k-1}$ has no linear term and commutes with B , and $R_k \in \mathcal{Y}_k$. For all $Y_k \in \mathcal{Y}_k$, note that Z consists of monomials with degree at least 2, which leads to $\widehat{Z} \circ \widehat{Y}_k = \widehat{Y}_k \circ \widehat{Z} = 0$. Hence we have $\exp(\widehat{Z} + \widehat{Y}_k) = \exp(\widehat{Z}) \circ \exp(\widehat{Y}_k)$. We then have equations:

$$\begin{aligned} \exp(-\widehat{Z} - \widehat{Y}_k) \circ \widehat{X} \circ \exp(\widehat{Z} + \widehat{Y}_k) &= \exp(-\widehat{Y}_k) \circ (\widehat{A} + \widehat{N} + \widehat{R}_k) \circ \exp(\widehat{Y}_k) \\ &= (\text{id} - \widehat{Y}_k) \circ (\widehat{A} + \widehat{N} + \widehat{R}_k) \circ (\text{id} + \widehat{Y}_k), \quad \text{since } \widehat{Y}_k \circ \widehat{Y}_k = 0 \\ &= \widehat{A} + \widehat{N} + \widehat{R}_k + \widehat{A} \circ \widehat{Y}_k - \widehat{Y}_k \circ \widehat{A}. \end{aligned}$$

We already proved that $(\phi_{Z+Y_k}^1)_*X$ is a normal form of order k if and only if $\left[((\phi_{Z+Y_k}^1)_*X)^\widehat{} - \widehat{A}, \widehat{B} \right] = 0$. By the computation above, it's equivalent to

$$ad_{\widehat{B}} \left(\widehat{R}_k - ad_{\widehat{A}}(\widehat{Y}_k) \right) = 0.$$

The following proposition implies $ad_{\widehat{B}} \circ ad_{\widehat{A}}(\widehat{\mathcal{Y}}_k) = ad_{\widehat{B}}(\widehat{\mathcal{Y}}_k)$, hence the equation has a solution $Y_k \in \mathcal{Y}_k$, which finishes the proof.

Proposition 5. $\widehat{\mathcal{Y}}_k = (\widehat{\mathcal{Y}}_k \cap \ker(ad_{\widehat{B}})) + ad_{\widehat{A}}(\widehat{\mathcal{Y}}_k)$.

Proof. We give a proof only in the case where A is diagonalisable and $B = A$; the general case is quite more difficult. For this case, note that $\widehat{A} = \widehat{B}$ is diagonalisable. For simplicity, we omit the hats in the rest of the proof.

In an eigenbasis of A (of eigenvalues (λ_i)), we take $M = (m_{ij}) \in \text{End}(V)$ such that $[A, M] \in \ker(ad_B) = \ker(ad_A)$. This condition, written in coordinates, gives

$$\forall i, j, 0 = [[M, A], A]_{ij} = (\lambda_j - \lambda_i)^2 m_{ij},$$

which implies $(\lambda_j - \lambda_i) m_{ij} = 0 = [M, A]_{ij}$. Hence, we have

$$\ker(ad_A) \cap \text{im}(ad_A) = 0.$$

Moreover, $\dim(\ker(ad_A)) + \dim(\text{im}(ad_A)) = \dim(\text{End}(V))$. Thus we arrived to the equation. \square

Remark 2. When B is adjoint of A for the inner product \langle, \rangle , the same proof applies : if $[[M, A], B] = 0$,

$$\begin{aligned} 0 &= \langle [[M, A], B], M \rangle = \langle [M, A]B \rangle - \langle B[M, A], M \rangle \\ &= \langle [M, A], MA \rangle - \langle [M, A], AM \rangle = \|[M, A]\|^2 \end{aligned}$$

so $[M, A] = 0$ and $\ker(ad_A) \cap \text{im}(ad_A) = 0$.

4 Link with the Poincaré-Dulac theorem

In this section, we show that the main theorem includes the following one, exposed in [1] :

Theorem 6. (*Poincaré-Dulac*) Let $\mathcal{X}' = \{X \in \Gamma(T\mathbb{C}^n), \text{ holomorphic, such that } X(0) = 0\}$, \mathcal{G}' the group of local diffeomorphisms of \mathbb{C}^n that fix 0, $X \in \mathcal{X}'$, $A = [X]_1$ its linear part, $\lambda_1, \dots, \lambda_n$ the complex eigenvalues of A , (e_1, \dots, e_n) the eigenbasis of A in \mathbb{C}^n , and $k \geq 2$. Then

$$\exists \phi \in \mathcal{G}', \exists (\alpha_{m,s}) \in \mathbb{C}^{\mathcal{R}}, \phi_* X(x) = Ax + \sum_{(m,s) \in \mathcal{R}} \alpha_{m,s} x^m e_s + O(|x|^{k+1}),$$

where $\mathcal{R} = \{(m, s) \in \mathbb{N}^n \times \llbracket 1, n \rrbracket, \lambda^m = \lambda_s \text{ and } 2 \leq |m| \leq k\}$ and $x^m = \prod_{i=1}^n x_i^{m_i}$.

The last expression is the normal form of X . The monomials of the form $x^m e_s$, with (m, s) in \mathcal{R} , are called resonant terms. The dynamics of the flow strongly depends on \mathcal{R} , as we shall see in further examples. Thus this theorem tells which terms can be removed without changing the behaviour of φ_X .

Remark 3. The main theorem applies in \mathbb{R}^n , but this one is only valid for holomorphic vector fields in \mathbb{C}^n . Thus, in the proof, we identify \mathbb{C}^n with \mathbb{R}^{2n} through an isomorphism Φ .

In fact, the main theorem shows the similarity between Poincaré-Dulac and Poincaré-Birkhoff normal form theorems by reformulating both of them in a common context.

Proof. We just apply the main theorem in the case $B = A_s$.

We set $\mathcal{X} = \{X \in \Gamma(T\mathbb{R}^{2n}), X(0) = 0\}$ and $\mathcal{Y} = \Phi_*(\mathcal{X}')$. As the Lie bracket of holomorphic vector fields is an holomorphic vector field, \mathcal{Y} is a sub-algebra of \mathcal{X} . We also have $[\Phi^*X]_1 = \Phi^*[X]_1 = \Phi^*A$, and Φ^*B is the diagonalisable part of Φ^*A , so we can apply the main theorem.

We get $\phi = \Phi^*(\varphi_{Y_2+\dots+Y_k}^1) \in \mathcal{G}'$ such that $[\Phi_*(\phi_*X - A)]^k$ commutes with $\Phi_*(B)$, which is equivalent to the commutation of $[\phi_*X - A]^k$ and B .

We write the Taylor development of ϕ_*X :

$$\exists(\alpha_{m,s}) \in \mathbb{C}^{\mathbb{N}^n \times \llbracket 1, n \rrbracket}, \phi_*X(x) = Ax + \sum_{\substack{(m,s) \\ 2 \leq |m| \leq k}} \alpha_{m,s} x^m e_s + O(|x|^{k+1}).$$

As $B(x) = \sum_{s=1}^n \lambda_s x_s e_s$, we compute :

$$\begin{aligned} 0 &= [[\phi_*X - A]^k, B](x) \\ &= \sum_{\substack{(m,s) \\ 2 \leq |m| \leq k}} \alpha_{m,s} B(x)^m e_s - B \left(\sum_{\substack{(m,s) \\ 2 \leq |m| \leq k}} \alpha_{m,s} x^m e_s \right) + O(|x|^{k+1}) \\ &= \sum_{\substack{(m,s) \\ 2 \leq |m| \leq k}} \alpha_{m,s} \lambda^m x^m e_s - \sum_{\substack{(m,s) \\ 2 \leq |m| \leq k}} \alpha_{m,s} x^m \lambda_s e_s + O(|x|^{k+1}) \\ &= \sum_{\substack{(m,s) \\ 2 \leq |m| \leq k}} \alpha_{m,s} (\lambda^m - \lambda_s) x^m e_s + O(|x|^{k+1}). \end{aligned}$$

The family $(x^m e_s)_{1 \leq |m| \leq k, 1 \leq s \leq n}$ is a base of E^k , so all the terms of the sum are zero. If $(m, s) \in \mathcal{R}$, as $\lambda^m - \lambda_s \neq 0$, we conclude that $\alpha_{m,s} = 0$, which gives us the expected result. \square

5 A direct proof of Poincaré-Dulac theorem

As the main theorem is quite abstract and more complex than Poincaré-Dulac theorem, we give an elementary proof of the latest in the case of a diagonalisable linear part $A = B$. This proof shows how to choose the morphism ϕ . We once again use an induction on k . The case $k = 1$ is trivial. Let $k \geq 2$ such that the theorem holds for $k - 1$. We know that

$$\exists \phi \in \mathcal{G}', \exists (\alpha_{m,s}) \in \mathbb{C}^{\mathcal{R}}, \phi_* X(x) = Ax + \sum_{(m,s) \in \mathcal{R}} \alpha_{m,s} x^m e_s + v(x) + O(|x|^{k+1}),$$

with v a homogeneous polynomial of degree k containing no resonant term. We want to find another a homogeneous polynomial h of degree k such that the change of variables $y = x - h(x)$ leads to

$$\phi_* X(y) = Ay + \sum_{(m,s) \in \mathcal{R}} \alpha_{m,s} y^m e_s + O(|y|^{k+1}).$$

We compute

$$\begin{aligned} \phi_* X(y) &= \phi_* X(x - h(x)) = \phi_* X(x) - d_x(\phi_* X)(h(x)) + O(|x|^{2k}) \\ &= Ax + \sum_{(m,s) \in \mathcal{R}} \alpha_{m,s} x^m e_s + v(x) - \partial_x h Ax + O(|x|^{k+1}) \\ &= Ay + Ah(x) + \sum_{(m,s) \in \mathcal{R}} \alpha_{m,s} y^m e_s + v(x) - \partial_x h Ax + O(|x|^{k+1}). \end{aligned}$$

Thus we get the good result if and only if $v(x) = \partial_x h Ax - Ah(x) =: L_A h(x)$. The operator L_A is linear, so we can decompose

$$v(x) = \sum_{\substack{(m,s) \notin \mathcal{R} \\ |m|=k}} \alpha_{m,s} x^m e_s$$

and look for a solution of each equation $L_A h(x) = x^m e_s$. But

$$\begin{aligned} L_A(x^m e_s) &= \sum_{i=1}^n \partial_i x^m e_s A x_i - A x^m e_s \\ &= \sum_{i=1}^n m_i \frac{x^m}{x_i} e_s x_i - \lambda_s x^m e_s = (\lambda^m - \lambda_s) x^m e_s. \end{aligned}$$

The $\lambda^m - \lambda_s$ are different from zero for $(m, s) \notin \mathcal{R}$, so we just have to take

$$h(x) = \sum_{\substack{(m,s) \notin \mathcal{R} \\ |m|=k}} \frac{\alpha_{m,s}}{\lambda^m - \lambda_s} x^m e_s.$$

Finally, the change of variables $\psi(x) = y = x - h(x)$ is a local diffeomorphism, and $\psi_* \phi_* X$ is the normal form of X at order k , which concludes the proof.

Remark 4. Note that the last computations are the same in the two proofs : indeed, $L_A(h)$ is equal to the bracket of h with B in the case $B = A_s$.

6 Examples

In this section, we just study a few examples to understand the link between the previous theorems and the dynamics of the flow.

6.1 A system with only one resonance

Let $p, q \geq 2$ integers, $c \in \mathbb{R} \setminus \{0\}$, and $X : (x, y) \mapsto (x, py + cx^q)$ a holomorphic vector field of \mathbb{C}^2 . Its linear part is $A = \begin{pmatrix} 1 & 0 \\ 0 & p \end{pmatrix}$. The dynamical systems associated to A and X are respectively

$$\begin{cases} \frac{d}{dt}x = x \\ \frac{d}{dt}y = py \end{cases} \quad \text{and} \quad \begin{cases} \frac{d}{dt}x = x \\ \frac{d}{dt}y = py + cx^q \end{cases}$$

These systems can be solved explicitly. The first one has solutions of the form $x(t) = x_0 e^t$, $y(t) = y_0 e^{pt}$, which gives a trajectory $y(x) = y_0 \left(\frac{x}{x_0}\right)^p$.

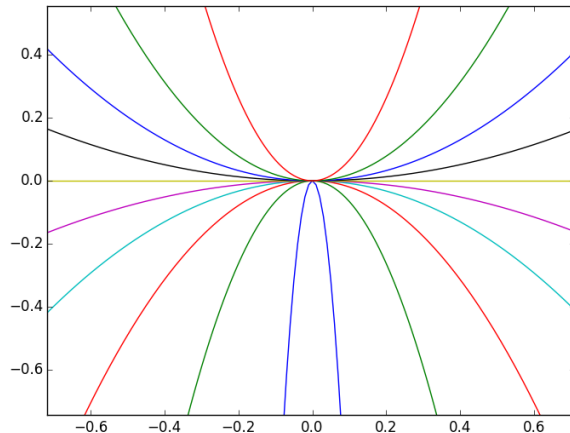


Figure 1: The trajectories of φ_A in the case $p = 2$, with $c = 1$. They all converge at 0, and three trajectories (including the one staying at 0) can be joined to form the C^∞ colored curves.

For the second one, the solution depends on the values of p and q .

- If $p = q$, we have $x(t) = x_0 e^t$, and $y(t) = (y_0 + c x_0^p t) e^{pt}$,
so $y(x) = y_0 \left(\frac{x}{x_0}\right)^p + c x^p \log\left(\frac{x}{x_0}\right)$.
- If $p \neq q$, we have $x(t) = x_0 e^t$, and $y(t) = \left(y_0 - \frac{c x_0^q}{q-p}\right) e^{pt} + \frac{c x_0^q}{q-p} e^{qt}$,
so $y(x) = \left(y_0 - \frac{c x_0^q}{q-p}\right) \left(\frac{x}{x_0}\right)^p + \frac{c}{q-p} x^q$.

Here, as $\lambda_1 = 1$ and $\lambda_2 = p$, the only resonance is $\mathcal{R} = \{(m, s)\} = \{((p, 0), 2)\}$. The resonant term $x^p e_y$ is non-zero only in the case $p = q$. Thus Poincaré-Dulac theorem states that X is locally isomorphic to A at order k for all $k \geq 2$ whenever $p \neq q$. Indeed, the holomorphic application

$$\begin{aligned} \phi : \quad \mathbb{C}^2 &\longrightarrow \mathbb{C}^2 \\ (x, y) &\longmapsto \left(x, y + \frac{c}{q-p} x^q\right) \end{aligned}$$

is a local diffeomorphism, and $\phi_* A = X$.

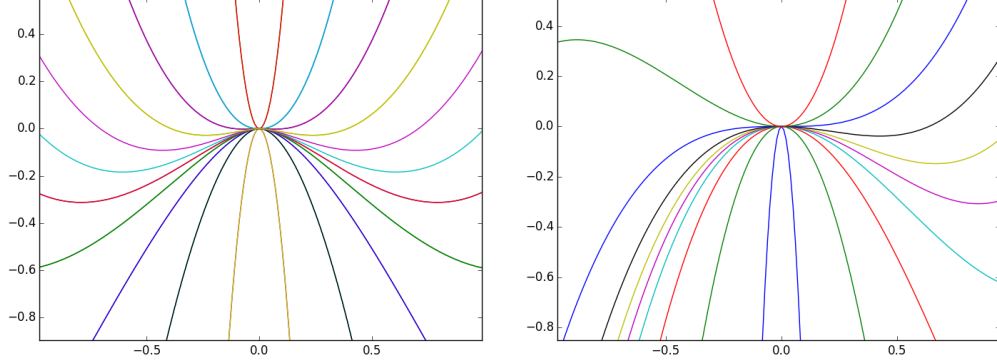


Figure 2: The trajectories of φ_X for $p = q = 2$. They are not C^2 at 0. Figure 3: The trajectories of φ_X for $p = 2$ and $q = 3$. They form C^∞ curbs.

At contrary, in the case $p = q$, we can find a C^{q-1} -diffeomorphism

$$\begin{aligned} \phi : \mathbb{C}^2 &\longrightarrow \mathbb{C}^2 \\ (x, y) &\longmapsto (x, y + c x^p \log x) \end{aligned}$$

but there is no C^q -diffeomorphism sending φ_A on φ_X . Indeed, we can join 3 trajectories of φ_A (the ones passing at (x_0, y_0) , $(0, 0)$ and $(-x_0, (-1)^p y_0)$ respectively) to make the C^∞ -curb $\left\{ \left(x, y_0 \left(\frac{x}{x_0} \right)^p \right) \right\}$, but that's not possible with φ_X because the function $y(x)$ written above is not C^p at 0. We then use the proposition below :

Proposition 7. *With the same notations as before, let $\phi \in \mathcal{G}_Y$ and $X, Y \in \mathcal{Y}$ such that $\phi_* X = Y$. If the trajectories of φ_X can be joined to form curbs C^k at 0, then the trajectories of φ_Y can be joined to form curbs C^k at 0 too.*

Remark 5. *As the trajectories are solutions to $\frac{d}{dt}x(t) = Y(x(t))$, the trajectories are C^∞ at any regular point.*

Proof. Let $c :] - \varepsilon, \varepsilon[\mapsto \mathbb{R}^n$ a curb following trajectories of φ_X , i.e. such that $\forall \theta \in] - \varepsilon, \varepsilon[$, $\frac{d}{d\theta}c(\theta)$ is colinear to $X(c(\theta))$. We assume that $c(0) = 0$ and c is C^k at 0. Then $\phi \circ c$ follows trajectories of Y because

$$\forall \theta \in] - \varepsilon, \varepsilon[, \quad \frac{d}{d\theta}\phi \circ c(\theta) = d\phi_{c(\theta)} \frac{d}{d\theta}c(\theta)$$

is colinear to $Y(c(\theta)) = \phi_* X(c(\theta)) = d\phi_{c(\theta)} X(c(\theta))$. □

Remark 6. The proposition remains true in \mathbb{C}^n with holomorphic vector fields, using the same argument as in Poincaré-Dulac theorem.

6.2 Harmonic oscillator

We now study a harmonic oscillator with a non-linear term, of the form

$$\begin{cases} \frac{d}{dt}x = y \\ \frac{d}{dt}y = -x + c x^a y^b \end{cases}$$

We know that the solutions of the linear system are of the form $x(t) = c_1 e^{it} + c_2 e^{-it}$, $y(t) = c_1 e^{it} - c_2 e^{-it}$, hence $|x|^2 + |y|^2 = 2(|c_1|^2 + |c_2|^2)$ is constant, so the trajectories are circles.

Here, $A = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$, so $\lambda_1 = i$ and $\lambda_2 = -i$. Hence

$$\mathcal{R} = \bigcup_{k \geq 0} \{((k+1, k), 1), ((k, k+1), 2)\}.$$

Consequently, the system above is resonant if and only if $b = a + 1$. When multiplying the second equation by y , we get

$$y \dot{y} = -x \dot{x} + c x^a y^{b+1},$$

so

$$\frac{1}{2} \frac{d}{dt} (|x|^2 + |y|^2) = x \dot{x} + y \dot{y} = c x^a y^{b+1}.$$

If a is even and $b = a + 1$, $x^a \dot{x}^{b+1} = (x \dot{x})^a \dot{x}^2 \geq 0$ and this term is different from zero (except at $x = 0$ or $y = 0$), so the radius $r = \sqrt{|x|^2 + |y|^2}$ depends monotonously on time at a finite order, increasing if $\Re(c) > 0$ and decreasing if $\Re(c) < 0$, and the trajectories are spirals.

On the other hand, when $b \neq a + 1$ the Poincaré-Dulac theorem implies that the trajectories remain close to cycles at any order k , as they are C^k diffeomorph to the circles of the linear system at order k .

To make it more clear, define $D(x, y) = |x(2\pi) - x(0)|^2 + |y(2\pi) - y(0)|^2$ to measure the evolution of the coordinates after a time 2π . In the linear case, φ_A is 2π -periodic, so $D(x, y) = 0$. With a non-linear term such that $b \neq a + 1$,

using Poincaré-Dulac theorem, $\forall k, D(x, y) = o(r^k)$. Finally, when $b = a + 1$, as the variations of r^2 are of order $a + b + 1 = 2a + 2$, $D(x, y) = \Theta(r^{a+1})$.

We now understand better the choice of the word *resonance*: a monomial term is resonant when he corresponds to a perturbation of the linear system that can affect the properties of the dynamics around the origin.

6.3 Coupled oscillators

We now study a system composed of two harmonic oscillators coupled through their position :

$$\begin{cases} \dot{x}_1 = \omega_1 y_1 + c x_2 \\ \dot{y}_1 = \omega_1 x_1 \\ \dot{x}_2 = \omega_2 y_2 + c x_1 \\ \dot{y}_2 = \omega_2 x_2 \end{cases}$$

with $\omega_1 > \omega_2 > 0$. Of course, each oscillator has its own resonances, as in the previous example. What we want to know is whether they are non-trivial resonances caused by the interaction between the oscillators. Let's compute the eigenvalues of A :

$$\begin{aligned} \det(A - uI) &= \det \begin{pmatrix} -u & \omega_1 & c & 0 \\ -\omega_1 & -u & 0 & 0 \\ c & 0 & -u & \omega_2 \\ 0 & 0 & -\omega_2 & -u \end{pmatrix} \\ &= (u^2 + \omega_1^2)(u^2 + \omega_2^2) - c^2 u^2 \\ &= u^4 + (\omega_1^2 + \omega_2^2 - c^2) u^2 + \omega_1^2 \omega_2^2. \end{aligned}$$

For c sufficiently small, the eigenvalues are

$$\lambda_{1,2} = i \left(\frac{\omega_1^2 + \omega_2^2 - c^2 \pm \sqrt{(\omega_1^2 + \omega_2^2 - c^2)^2 - 4\omega_1^2 \omega_2^2}}{2} \right)^{\frac{1}{2}}$$

and their complex conjugates. The non-trivial resonances are the ones of the form $a \lambda_1 + b \lambda_2 = 0$, with $(a, b) \in \mathbb{N}^2 \setminus \{(0, 0)\}$. The function $f(c) = \frac{\lambda_1}{\lambda_2}(c)$ is continuous, monotonic, defined on $[0, \sqrt{\omega_1^2 + \omega_2^2}]$, and reaches $\frac{\omega_1^2}{\omega_2^2}$ at 0 and 1 at $\sqrt{\omega_1^2 + \omega_2^2}$, so f reaches rational values on an infinite countable set of values of c .

As a conclusion, non-trivial resonant terms exist only for specific values of c , ω_1 and ω_2 . We can see the eigenvalues as modified frequencies, because $\lambda_1 \xrightarrow{c \rightarrow 0} \omega_1$, and $\lambda_2 \xrightarrow{c \rightarrow 0} \omega_2$. The resonances appear whenever the ratio of these frequencies is rational.

7 Generalisation to a family of vector fields

It's very natural to consider the regularity of the normal forms of a family of vector fields. That's to say if we have a family of vector fields X_λ singular at 0 parameterised by $\lambda \in \Lambda$, and depends on λ with some regularity, then we want to know if there's some family (Y_λ^k) of normal forms of the X_λ with the same regularity on λ . We're going to formalise the problem and give it an affirmative answer.

Definition 1. *Assume Λ is a Banach space, we say the family of vector fields X_λ depends on λ C^r regularly, if the map*

$$\begin{aligned} \Lambda \times U &\longrightarrow \mathbb{R}^n \\ (\lambda, x) &\longmapsto X_\lambda(x) \end{aligned}$$

is C^r smooth as a map from an open subset of the product Banach space $\Lambda \times \mathbb{R}^n$ to \mathbb{R}^n . Here $U \ni 0$ is an open set of \mathbb{R}^n and we regard $X_\lambda \in \Gamma(TU)$ as a map $X_\lambda : U \rightarrow \mathbb{R}^n$ by the natural identification $T\mathbb{R}^n \approx \mathbb{R}^{n+n}$ mentioned in the beginning.

The answer to the problem is stated here.

Theorem 8. *Let Λ be a Banach space, X_λ a C^r family of vector fields, $A \in \mathcal{Y}_1$ the linear part of X_0 at 0, and we choose B as in the main theorem. Then, there exists a C^r family of vector fields $Y_\lambda^k \in \mathcal{Y}^k$, such that Y_0^k has no linear term, and there exists an open set $V \subseteq \Lambda$, such that $\forall \lambda \in V, [(\phi_{Y_\lambda^k}^1)_* X_\lambda - A]^k, B = 0$.*

Proof. Similarly to the proof of the main theorem, we prove it by induction on k . Suppose that we have built C^r vector fields Y_λ^{k-1} , $k \geq 2$. We then, by the same computation, get the equation:

$$ad_{\widehat{B}}(\widehat{R}_\lambda) - ad_{\widehat{B}} \circ ad_{\widehat{A}_\lambda}(\widehat{Y}_\lambda) = 0,$$

where R_λ and A_λ are the homogenous term of degree k and the linear term of $(\phi_{Y_\lambda}^1)_{*} X_\lambda$ respectively.

The induction hypotheses implies that R_λ and A_λ are C^r families. By proposition 5, we have $ad_{\widehat{B}} \circ ad_{\widehat{A}} : \widehat{\mathcal{Y}}_k \rightarrow ad_{\widehat{B}}(\widehat{\mathcal{Y}}_k)$ is surjective. We may choose a $K \subseteq \widehat{\mathcal{Y}}_k$ as the complementary of its kernel, thus

$$ad_{\widehat{B}} \circ ad_{\widehat{A}} : K \rightarrow ad_{\widehat{B}}(\widehat{\mathcal{Y}}_k)$$

is an isomorphism. By the regularity of A_λ , $\exists V \subseteq \Lambda$ open set, such that $\forall \lambda \in V, ad_{\widehat{B}} \circ ad_{\widehat{A}_\lambda} : K \rightarrow ad_{\widehat{B}}(\widehat{\mathcal{Y}}_k)$ is an isomorphism. We denote

$$\theta_\lambda = (ad_{\widehat{B}} \circ ad_{\widehat{A}_\lambda})^{-1} : ad_{\widehat{B}}(\widehat{\mathcal{Y}}_k) \rightarrow K,$$

then θ_λ is C^r on λ . Take $Y_\lambda = -\theta_\lambda \circ ad_{\widehat{B}}(R_\lambda)$, which is a C^r family of vector fields and solves the equation.

We finally need to set the initial condition, i.e. to show the existence of a C^r family of vector fields $L_\lambda \in \mathcal{Y}_1$, such that $L_0 = 0$ and

$$\forall \lambda \in V, [\exp(-\widehat{L}_\lambda) \widehat{A}_\lambda \exp(\widehat{L}_\lambda) - \widehat{A}, \widehat{B}] = 0.$$

Suppose $K \subseteq \widehat{\mathcal{Y}}_1$ a complementary of $\ker(ad_{\widehat{B}} \circ ad_{\widehat{A}})$. Consider the map

$$\begin{aligned} F : K \times \Lambda &\rightarrow ad_{\widehat{B}}(\widehat{\mathcal{Y}}_1) \\ (L, \lambda) &\mapsto ad_{\widehat{B}}(\exp(-\widehat{L}) \widehat{A}_\lambda \exp(\widehat{L}) - \widehat{A}). \end{aligned}$$

We have $\partial_L F(0, 0) = ad_{\widehat{B}} \circ ad_{\widehat{A}}$, which is an isomorphism from $K \approx T_0 K$ to $ad_{\widehat{B}}(\widehat{\mathcal{Y}}_1) \approx T_0(ad_{\widehat{B}}(\widehat{\mathcal{Y}}_1))$. Then by the implicit function theorem, we have a C^r family of linear maps \widehat{L}_λ which lifts back to a C^r family of vector fields L_λ .

□

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