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Decomposition of homogeneous vector fields of degree one and representation of the flow*

by

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ABSTRACT. – We first give a characterization for the set of real analytic diffeomorphisms which transform homogeneous vector fields of certain degree into homogeneous fields of the same degree with respect to an arbitrary dilation δ_ε^r . Such a set is constituted by the invertible analytic maps that are homogeneous of degree one with respect to δ_ε^r and can be endowed with the structure of a Lie Group whose Lie algebra is the space $H^{1,r}(\mathbb{R}^n)$ of the homogeneous fields of degree one with respect to δ_ε^r . Then we prove a decomposition theorem for the elements of the non semisimple Lie algebra $H^{1,r}(\mathbb{R}^n)$. This result is a non linear analog of the Jordan decomposition of a linear field, *i.e.* for $X \in H^{1,r}(\mathbb{R}^n)$, we can write $X = S + N$, with S linear semisimple and $[S, N] = 0$. We also give an explicit representation formula for the flow generated by a field in $H^{1,r}(\mathbb{R}^n)$. Finally we apply this result to obtain a simple representation for the trajectories of a class of affine control systems $\dot{x} = X_0(x) + Bu$, with $X_0 \in H^{1,r}(\mathbb{R}^n)$ and B a constant field, that constitute a natural extension of the linear control systems.

Key words: Homogeneous vector fields, decomposition of vector fields, nonlinear systems, representation of solutions, Lie algebra of a Lie group.

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RÉSUMÉ. – Tout d’abord, nous donnons une caractérisation de l’ensemble des difféomorphismes analytiques réels qui transforment des champs de vecteurs homogènes d’un certain degré en champs de vecteurs homogènes de même degré relativement à une dilatation arbitraire δ_ε^r . Un tel ensemble est constitué par des applications analytiques inversibles homogènes de degré 1 relativement à δ_ε^r , et il peut être doté d’une structure de groupe de Lie. L’espace $H^{1,r}(\mathbb{R}^n)$ des champs de vecteurs homogènes de degré 1 relativement à δ_ε^r est l’algèbre de Lie de cet ensemble. Ensuite, nous démontrons un théorème de décomposition pour les éléments de l’algèbre de Lie non-semisimple $H^{1,r}(\mathbb{R}^n)$. Ce résultat est l’analogie non linéaire de la décomposition de Jordan d’un champ linéaire, *i.e.*, pour $X \in H^{1,r}(\mathbb{R}^n)$, nous pouvons écrire $X = S + N$, où S est un champ linéaire semisimple et $[S, N] = 0$. Nous donnons aussi une formule explicite de représentation pour le flux d’un champ de $H^{1,r}(\mathbb{R}^n)$. Finalement, nous utilisons ce résultat pour obtenir une représentation simple des trajectoires d’une classe du système affines de contrôle $\dot{x} = X_0(x) + Bu$, où $X_0 \in H^{1,r}(\mathbb{R}^n)$ et B est un champ constant, qui constitue une extension naturelle des systèmes linéaires de contrôle.

1. INTRODUCTION

Consider an affine nonlinear control system

$$\dot{x} = X_0(x) + \sum_{j=1}^m u_j X_j(x), \quad (1.1)$$

where X_0, X_1, \dots, X_m are real analytic vector fields on \mathbb{R}^n and $u = (u_1, \dots, u_m)$ is the control. A well known technique for the local study of such a system consists in locally approximating the vector fields $X_j, j = 0, 1, \dots, m$, by fields $Y_j, j = 0, 1, \dots, m$, for which the analysis is easier and such that they “preserve” the property being studied. This technique has been the key in obtaining high order local controllability results and in the construction of asymptotically stabilizing feedback controls, *e.g.*, *see* [5], [19], [10], [11], [12], [13], [16]. Homogeneous vector fields with respect to a dilation δ_ε^r have often provided such “correct” approximations in basically non linear problems for which the usual linear approximations fail to yield sufficient information. Results in this direction can be found in [4], [9], [20].

In this paper we study some basic properties of the homogeneous vector fields of degree one with respect to an arbitrary dilation δ_ε^r and develop a general method of solution for every autonomous system of differential equations

$$\dot{x} = X(x), \tag{1.2}$$

where X is a real analytic vector field on \mathbb{R}^n (or on an n -dimensional manifold M^n) homogeneous of degree one with respect to δ_ε^r .

Specifically, let $\delta_\varepsilon^r : \mathbb{R}^n \rightarrow \mathbb{R}^n$ be a dilation on \mathbb{R}^n defined by $\delta_\varepsilon^r(x) = (\varepsilon^{r_1}x_1, \dots, \varepsilon^{r_n}x_n)$, where $\varepsilon > 0$, and $r_1 \leq \dots \leq r_n$ are positive integers. A polynomial $p : \mathbb{R}^n \rightarrow \mathbb{R}$ is *homogeneous of degree* $m \in \mathbb{N}$ with respect to δ_ε^r if $p(\delta_\varepsilon^r(x)) = \varepsilon^m p(x)$. (Throughout \mathbb{N} will denote the set of non negative integers $\{0, 1, \dots\}$.) The set of polynomials homogeneous of degree m with respect to δ_ε^r will be denoted $P^{m,r}(\mathbb{R}^n)$. We define a real analytic *vector field* $X(x) = \sum_{i=1}^n a_i(x)\partial/\partial x_i$ on \mathbb{R}^n , given in local coordinates $x = (x_1, \dots, x_n)$, to be *homogeneous of degree* $m \in \mathbb{Z}$ with respect to δ_ε^r if $a_i \in P^{r_i+m-1,r}(\mathbb{R}^n)$, $i = 1, \dots, n$. We denote by $H^{m,r}(\mathbb{R}^n)$ the family of such vector fields. This definition (although not universally used) agrees with the classical definition of homogeneity (i.e., $a_i(\varepsilon x) = \varepsilon^m a_i(x)$, $i = 1, \dots, n$) in the case X is homogeneous with respect to the *standard dilation* δ_ε^1 , having $r_1 = \dots = r_n = 1$. In particular, a field $X(x) = Ax$ that is linear in the local coordinates will be homogeneous of degree one w.r.t. δ_ε^1 . One can thus regard the concept of homogeneity of degree one w.r.t. an arbitrary dilation as a natural extension of the concept of linearity. In fact, classical results valid for linear vector fields have been obtained for such fields in nonlinear problems, where the homogeneous fields play the role of the linear approximations in classical theory, e.g., see [10].

The paper is organized as follows: in Section 2 we characterize the set of all real analytic diffeomorphisms ϕ on \mathbb{R}^n that transform any homogeneous vector field X of degree m with respect to a given dilation δ_ε^r into a homogeneous field $T_\phi X$ of the same degree with respect to the same dilation. If δ_ε^r is the standard dilation δ_ε^1 , such a set is clearly the set of all invertible linear transformations on \mathbb{R}^n , which in our notation coincides with the set of all invertible elements of $H^{1,1}(\mathbb{R}^n)$. We show that also in the case of an arbitrary dilation δ_ε^r , $r = (r_1, \dots, r_n)$, with $r_1 = 1$, the set of all real analytic changes of coordinates that transform homogeneous vector fields of certain degree into homogeneous fields of the same degree is precisely the set of all invertible elements of $H^{1,r}(\mathbb{R}^n)$, which will be denoted by $GH^{1,r}(\mathbb{R}^n)$. This set is a subgroup of the group of all real analytic diffeomorphisms on \mathbb{R}^n and can be endowed

with the structure of a finite dimensional Lie group. Moreover, if X is a field in $H^{1,r}(\mathbb{R}^n)$ and $(\exp tX)(p)$ denotes the solution, at time t , of the Cauchy problem $\dot{x} = X(x), x(0) = p$, then, for t fixed, the diffeomorphism $p \rightarrow (\exp tX)(p)$ lies in $GH^{1,r}(\mathbb{R}^n)$. It follows as an easy consequence that the space $H^{1,r}(\mathbb{R}^n)$ of the homogeneous fields of degree one with respect to a dilation δ_ε^r is the Lie algebra of the Lie group $GH^{1,r}(\mathbb{R}^n)$.

In Section 3 we prove a decomposition theorem for the vector fields of the non semisimple Lie algebra $H^{1,r}(\mathbb{R}^n)$, providing a non linear analog to the Jordan decomposition of a linear vector field into a semisimple and nilpotent part.

THEOREM 3.1. – *Let $X(x) = \sum_{i=1}^n a_i(x)\partial/\partial x_i$ be a real analytic vector field on \mathbb{R}^n homogeneous of degree one with respect to a given dilation δ_ε^r . Then there is a polynomial change of coordinates $x = \phi(y), \phi \in GH^{1,r}(\mathbb{R}^n)$, such that*

$$T_\phi X = S + N, \tag{1.3}$$

($T_\phi X(y)$ denoting the transformed field after performing the coordinate change $x = \phi(y)$) where S and N are real analytic homogeneous vector fields of degree one with respect to δ_ε^r , satisfying

$$[S, N] = 0. \tag{1.4}$$

Moreover S is a linear semisimple vector field (i.e. the complexification of S is diagonalizable) and N is the sum of a linear nilpotent field and of a strictly non linear homogeneous vector field of degree one with respect to δ_ε^r .

Finally, in Section 4 we obtain a simple representation of the solutions of system (1.2), in terms of their Picard approximations, for a class of fields that can be regarded as a generalization of the linear nilpotent vector fields. This result, together with the decomposition theorem given in Section 3, yields a representation formula for the solutions of (1.2), for any $X \in H^{1,r}(\mathbb{R}^n)$. We also derive an explicit representation for the trajectories $x(\cdot, u)$ of an n -dimensional, single input, affine control system

$$\dot{x} = X_0(x) + B u, \tag{1.5}$$

where X_0 is an element of $H^{1,r}(\mathbb{R}^n)$, $r = (r_1, \dots, r_n)$, and B is a constant field whose local coordinate expression is given by an $n \times 1$ matrix having nonzero entries b_i only for those i such that $r_i = r_n$:

$$x(t, u) = (\exp tX_0) \left(x(0) + \int_0^t e^{-sA} B u(s) ds \right), \quad t \in \mathbb{R},$$

denoting by A the linear part of the field X_0 .

2. HOMOGENEOUS DIFFEOMORPHISMS OF DEGREE ONE WITH RESPECT TO AN ARBITRARY DILATION

2.1. Notations and Definitions

The general setting for the theory which follows is the Lie algebra of all real analytic vector fields on a real analytic n -dimensional manifold M^n . However, we deal only with local problems in which we study the vector fields on some neighborhood \mathcal{U} of a point $p \in M^n$. Therefore, instead of constantly referring to local coordinate charts $\chi : \mathcal{U} \rightarrow \mathbb{R}^n$, with $\chi(p) = 0$, we will identify any point q in \mathcal{U} with its local coordinate expression $x = \chi(q)$ in $\chi(\mathcal{U}) = \mathcal{U}$, and take the local viewpoint that our vector fields are defined on an open neighborhood U of $0 \in \mathbb{R}^n$.

We will denote a real analytic vector field X , in local coordinates $x = (x_1, \dots, x_n)$, equivalently by

$$X(x) = \sum_{i=1}^n a_i(x) \frac{\partial}{\partial x_i}, \quad X(x) = \begin{pmatrix} a_1(x) \\ \vdots \\ a_n(x) \end{pmatrix},$$

where each $a_i(x)$ is a real analytic function of x .

Given two real analytic vector fields X, Y , we let $[X, Y]$ denote their Lie product which, in local coordinates, is expressed by $(\partial X / \partial x)(x)Y(x) - (\partial Y / \partial x)(x)X(x)$, denoting by $(\partial X / \partial x)(x)$, $(\partial Y / \partial x)(x)$ respectively the Jacobians of X and Y . In order to simplify the notation, if X and Y are smooth maps from \mathbb{R}^n to \mathbb{R}^n , we still denote by $[X, Y]$ the map defined by the above expression even in cases where we do not interpret X, Y as local coordinate expressions of two vector fields. We also use $adX(Y) \doteq (adX, Y) \doteq [X, Y]$ and, inductively, $(ad^{k+1}X, Y) \doteq [X, (ad^kX, Y)]$.

If $X(x)$ denotes a vector field given in the x -coordinates and we perform a coordinate change $x = \phi(y)$, we denote by $T_\phi X(y)$ the transformed field expressed in the y -coordinates, which is given by $T_\phi X(y) = ((\partial \phi / \partial y)(y))^{-1} X(\phi(y))$.

For a given vector field X and $p \in \mathbb{R}^n$, we denote by $(\text{exp } tX)(p)$ the solution, at time t , of the Cauchy problem $\dot{x} = X(x)$, $x(0) = p$; thus the map $p \mapsto (\text{exp } tX)(p)$, represents the flow generated by the field X .

Throughout it will be used the term *strictly non linear* to denote an analytic map whose Taylor expansion starts with a homogeneous term (with respect to the standard dilation) of degree greater than one.

The definitions of dilation and homogeneity w.r.t. a dilation have already been recalled in the introduction. In particular we will use the

following notation. Let $\delta_\varepsilon^r : \mathbb{R}^n \rightarrow \mathbb{R}^n$ be a dilation on \mathbb{R}^n , $\delta_\varepsilon^r(x) = (\varepsilon^{r_1} x_1, \dots, \varepsilon^{r_n} x_n)$, $\varepsilon > 0$. The nondecreasing n -tuple of positive integers $r = (r_1, \dots, r_n)$ will be said to be of the type $(i_1, \dots, i_m; j_1, \dots, j_m)$ if the following equalities hold:

$$\begin{aligned} r_1 &= r_2 = \dots = r_{i_1} = j_1, \\ r_{i_1+1} &= r_{i_1+2} = \dots = r_{i_2} = j_2, \\ &\vdots \qquad \qquad \qquad \vdots \qquad \qquad \qquad \vdots \\ r_{i_{(m-1)}+1} &= r_{i_{(m-1)}+2} = \dots = r_{i_m} = j_m, \end{aligned} \tag{2.1}$$

with $i_m = n$ and $j_1 < j_2 < \dots < j_m$. (Throughout is used $i_0 = 0$.) For certain results the additional assumption $j_1 = 1$ will be required. Given a map $f : \mathbb{R}^n \mapsto \mathbb{R}^n$ with $f_i \in P^{r_i+m-1, r}(\mathbb{R}^n)$, $i = 1, \dots, n$, we write $f \in H^{m, r}(\mathbb{R}^n)$ even in cases where f does not denote the local coordinate expression of some vector field.

2.2. Statements of the main results

We here summarize the results presented in this section. Given a dilation δ_ε^r with $r = (r_1, \dots, r_n)$ of the type $(i_1, \dots, i_m; j_1, \dots, j_m)$ as defined in (2.1), let Z denote the vector field expressed, in x -local coordinates, by

$$Z(x) = \sum_{i=1}^n r_i x_i \frac{\partial}{\partial x_i}. \tag{2.2}$$

THEOREM 2.1. – *Let ϕ be a real analytic diffeomorphism on \mathbb{R}^n such that $\phi(0) = 0$. If $j_1 = 1$ in (2.1), then the following statements are equivalent:*

(i) *There exists an integer $m \geq 0$ such that*

$$X \in H^{m, r}(\mathbb{R}^n) \implies T_\phi X \in H^{m, r}(\mathbb{R}^n);$$

(ii)

$$T_\phi Z = Z;$$

(iii)

$$\phi \in H^{1, r}(\mathbb{R}^n).$$

If $j_1 > 1$ in (2.1), then (ii), (iii) are equivalent and they imply (i).

THEOREM 2.2. – *A map $\phi : \mathbb{R}^n \rightarrow \mathbb{R}^n$ is a real analytic diffeomorphism, homogeneous of degree one with respect to the dilation δ_ε^r , if and only if*

it has the form

$$\phi(x) = Ax + g(x), \tag{2.3}$$

where A is a matrix of the form

$$A = \text{diag}\{A_1, \dots, A_m\}, \tag{2.4}$$

each A_k being an invertible $(i_k - i_{k-1}) \times (i_k - i_{k-1})$ matrix with real entries, and g is a strictly non linear function in $H^{1,r}(\mathbb{R}^n)$ of the form

$$g(x) = \sum_{j=2}^{j_m} g^j(x), \tag{2.5}$$

each g^j being homogeneous of degree j with respect to the standard dilation δ_ε^1 .

Let $GL(\mathbb{R}^n)$ denote the set of all invertible linear transformations on \mathbb{R}^n . We will sometimes identify the elements of $GL(\mathbb{R}^n)$ with their matrix representation, denoting by A the linear transformation $x \mapsto Ax$. Consider the sets of transformations

$$\begin{aligned} G_1 &= \{A \in GL(\mathbb{R}^n) : A \text{ has the form (2.4)}\}, \\ G_2 &= \{\phi \in H^{1,r}(\mathbb{R}^n) : \phi = I + g\}, \\ GH^{1,r}(\mathbb{R}^n) &= \{\phi \in H^{1,r}(\mathbb{R}^n) : \phi = A + g, A \text{ is invertible}\}, \end{aligned}$$

(here I denotes as usual the identity matrix, and g the strictly non linear part of ϕ). From Theorems 2.1 and 2.2 it follows that, in the case δ_ε^r , $r = (r_1, \dots, r_n)$, is a dilation with $r_1 = 1$, $GH^{1,r}(\mathbb{R}^n)$ coincides with the set of all real analytic diffeomorphisms that transform elements of $H^{m,r}(\mathbb{R}^n)$ into elements of the same space.

THEOREM 2.3. – *The set $GH^{1,r}(\mathbb{R}^n)$ is a finite dimensional Lie group with the dimension depending on $r = (r_1, \dots, r_n)$. Moreover, the sets G_1, G_2 defined above are subgroups of $GH^{1,r}(\mathbb{R}^n)$, and for any $\phi \in GH^{1,r}(\mathbb{R}^n)$ there exist $\phi_1 \in G_1$ and $\phi_2, \psi_2 \in G_2$, such that*

$$\phi = \phi_1 \circ \phi_2 = \psi_2 \circ \phi_1. \tag{2.6}$$

Remark 2.4. – The group $GH^{1,r}(\mathbb{R}^n)$ is the natural generalization of the group of the invertible linear transformations $GL(\mathbb{R}^n)$, with which it coincides in the case δ_ε^r is the standard dilation δ_ε^1 .

THEOREM 2.5. – Let $X \in H^{1,r}(\mathbb{R}^n)$. Then, for each fixed $t \in \mathbb{R}$, the map

$$p \rightarrow (\exp tX)(p), \quad p \in \mathbb{R}^n$$

is a homogeneous analytic diffeomorphism of degree one with respect to δ_ε^r . Moreover, $H^{1,r}(\mathbb{R}^n)$ is the Lie algebra of the Lie group $GH^{1,r}(\mathbb{R}^n)$.

Example 2.6. – On \mathbb{R}^2 let $X(x) = 2x_1 \partial/\partial x_1 + (5x_1^3 - x_2) \partial/\partial x_2$. Note that X is a field in $H^{1,r}(\mathbb{R}^2)$, with $r = (1, 3)$. It can be easily computed that the flow-map of X is given by $(\exp tX)(p_1, p_2) = (e^{2t}p_1, (5/7)(e^{6t} - e^{-t})p_1^3 + e^{-t}p_2)$, which is clearly an element of $GH^{1,r}(\mathbb{R}^2)$ for any fixed $t \in \mathbb{R}$.

2.3. Preliminary lemmas

We collect here several preliminary results that will enable us to prove Theorem 2.1.

LEMMA 2.7. – Let Z be as in (2.2). A real analytic vector field $X(x) = \sum_{i=1}^n a_i(x) \partial/\partial x_i$ is homogeneous of degree m with respect to δ_ε^r if and only if

$$[X, Z] = (m - 1)X. \quad (2.7)$$

Proof. – The i -th component of the field $[X, Z]$ is given by

$$\sum_{j=1}^n \frac{\partial a_i}{\partial x_j}(x) r_j x_j - r_i a_i(x). \quad (2.8)$$

Since X is analytic, its components a_i can be expanded in terms of homogeneous polynomials; let \bar{a}_i denote a monomial of a_i of the form $\bar{a}_i(x) = ax_1^{\nu_1} \cdots x_n^{\nu_n}$. Then we have

$$\sum_{j=1}^n \frac{\partial \bar{a}_i}{\partial x_j}(x) r_j x_j = \sum_{j=1}^n \frac{\bar{a}_i(x)}{x_j} \nu_j r_j x_j = \left(\sum_{j=1}^n \nu_j r_j \right) \bar{a}_i(x).$$

Substituting the above in (2.8), for any monomial \bar{a}_i of a_i , we deduce that the relation (2.7) holds if and only if the exponents ν_1, \dots, ν_n of any monomial $\bar{a}_i(x) = ax_1^{\nu_1} \cdots x_n^{\nu_n}$ of a_i , satisfy the relation

$$\sum_{j=1}^n \nu_j r_j - r_i = m - 1,$$

for each $i = 1, \dots, n$. This means that $a_i \in P^{r_i+m-1,r}(\mathbb{R}^n)$, for each $i = 1, \dots, n$ which is equivalent to say that $X \in H^{m,r}(\mathbb{R}^n)$. \square

LEMMA 2.8. – Let X be a real analytic vector field on \mathbb{R}^n , homogeneous of degree m with respect to δ_ε^r . If ϕ is a real analytic diffeomorphism such that

$$T_\phi Z = Z, \tag{2.9}$$

then also $T_\phi X$ is an homogeneous field of degree m with respect to δ_ε^r .

Proof. – Recall that the Lie product is a coordinate-free operation; thus, transforming both sides of (2.7), we obtain

$$T_\phi[X, Z] = [T_\phi X, T_\phi Z] = (m - 1)T_\phi X. \tag{2.10}$$

Hence, substituting (2.9) into (2.10), we have $[T_\phi X, Z] = (m - 1)T_\phi X$ which, using lemma 2.7, enables us to conclude. \square

LEMMA 2.9. – Let $Y(x) = \sum_{i=1}^n a_i(x) \partial/\partial x_i$ be a non-zero real analytic homogeneous vector field of degree $j \in \mathbb{Z}$ with respect to $\delta_\varepsilon^r, r = (r_1, \dots, r_n)$, with $r_1 = 1$. Then

$$[X, Y] = 0 \text{ for all } X \in H^{m,r}(\mathbb{R}^n), \tag{2.11}$$

for some $m \geq 0$, if and only if one of the following two conditions holds

- (i) $\delta_\varepsilon^r = \delta_\varepsilon^1, \quad m = 0, \quad Y(x) = \sum_{i=1}^n a_i \partial/\partial x_i, \quad a_i \in \mathbb{R};$
- (ii) $m = 1, \quad Y(x) = aZ(x), \quad a \in \mathbb{R}.$

Proof. – That (i) implies (2.11) is immediate since an homogeneous vector field of degree zero with respect to the standard dilation is a constant vector field and the Lie bracket of two constant fields is clearly zero.

Next suppose (ii); from lemma 2.7, using the linearity of the Lie product, it follows

$$[X, aZ] = a(m - 1)X, \text{ for any } X \in H^{m,r}(\mathbb{R}^n).$$

Therefore, since $m = 1$, (2.11) is satisfied.

Now suppose that (2.11) holds for some $m \geq 0$. Note first that, since $m \geq 0$ implies $r_k + m - 1 \geq 0, \quad k = 1, \dots, n$, then $\{x_1^{r_k+m-1} \partial/\partial x_k, \quad k = 1, \dots, n\}$ is a set of non zero fields in $H^{m,r}(\mathbb{R}^n)$. Therefore (2.11) is in

particular satisfied by the elements of this set:

$$0 = [x_1^{r_k+m-1} \partial / \partial x_k, Y]$$

$$= \begin{pmatrix} -\frac{\partial a_1}{\partial x_k}(x) x_1^{r_k+m-1} \\ \vdots \\ (r_k + m - 1)x_1^{r_k+m-2} a_1(x) - \frac{\partial a_k}{\partial x_k}(x) x_1^{r_k+m-1} \\ \vdots \\ -\frac{\partial a_n}{\partial x_k}(x) x_1^{r_k+m-1} \end{pmatrix}$$

$$k = 1, \dots, n. \quad (2.12)$$

Thus it follows

$$\frac{\partial a_i}{\partial x_k}(x) \equiv 0 \quad i, k = 1, \dots, n, \quad i \neq k, \quad (2.13)$$

which implies, using $Y \in H^{j,r}(\mathbb{R}^n)$ and therefore $a_i \in P^{r_i+j-1,r}(\mathbb{R}^n)$,

$$a_i(x) = a_i x_i^{\nu_i}, \quad a_i \in \mathbb{R}, \quad r_i \nu_i = r_i + j - 1, \quad i = 1, \dots, n. \quad (2.14)$$

We now consider three cases:

Case 1. – Suppose $a_1 = 0$. From (2.12) it follows $(\partial a_k / \partial x_k)(x) \equiv 0$, for each $k = 1, \dots, n$ and so $\nu_i = 0$ for each $i = 1, \dots, n$ in (2.14), which implies $r_i = 1$, for each $i = 1, \dots, n$, and $j = 0$. Thus condition (i) is satisfied.

Case 2. – Suppose $a_1 \neq 0$, $\nu_k = 0$ for some k , $1 \leq k \leq n$. Then $m \geq 0$ and (2.12) imply $r_k = 1$, $m = 0$. Hence $r_i = 1$ for each $i = 1, \dots, k$ which, using again (2.12), implies $(\partial a_i / \partial x_i)(x) \equiv 0$, for each $i = 1, \dots, k$. Thus $a_i(x) \equiv a_i \in \mathbb{R}$, $\nu_i = 0$, for each $i = 1, \dots, k$ and $j = 0$. If $\nu_i > 0$ for some $i > k$, then, using (2.14), we would have $j = r_i(\nu_i - 1) + 1 \geq 1$ which gives a contradiction. Hence $\nu_i = 0$ for each $i = 1, \dots, n$ and condition (i) is satisfied.

Case 3. – Suppose $a_1 \neq 0$, $\nu_i > 0$ for each $i = 1, \dots, n$. From (2.12), using (2.14), it follows

$$0 = (r_k + m - 1)x_1^{r_k+m-2} a_1 x_1^{\nu_1} - a_k \nu_k x_k^{\nu_k-1} x_1^{r_k+m-1},$$

$$k = 1, \dots, n, \quad (2.15)$$

which implies $\nu_k = 1$ for each $k = 2, \dots, n$, and so by (2.14) $j = 1$. Moreover, since $r_1 = 1$, (2.14) implies also $\nu_1 = j$. Thus $\nu_i = 1$ for each $i = 1, \dots, n$. Finally (2.15), with $k = 1$, implies $m = 1$, and, with $k = 2, \dots, n$, implies $a_k = r_k a_1$ for each $k = 2, \dots, n$. Hence condition (ii) is satisfied with $a = a_1$. \square

LEMMA 2.10. – *Let $Y(x) = \sum_{i=1}^n a_i(x) \partial/\partial x_i$ be a non-zero real analytic vector field on \mathbb{R}^n . Then the same conclusions of lemma 2.9 hold.*

Proof. – First expand Y in homogeneous vector fields with respect to δ_ε^r , i.e.,

$$Y(x) = \sum_{j=1-r_n}^{\infty} Y^j(x), \quad Y^j \in H^{j,r}(\mathbb{R}^n).$$

Next observe that, since $[X, Y^j], j \geq 1 - r_n$, are homogeneous vector fields of different degrees, the equation

$$0 = [X, Y] = \sum_{j=1-r_n}^{\infty} [X, Y^j], \text{ for any } X \in H^{m,r}(\mathbb{R}^n),$$

is satisfied if and only if the equations

$$[X, Y^j] = 0, \text{ for any } j \geq 1 - r_n, X \in H^{m,r}(\mathbb{R}^n),$$

are satisfied simultaneously. Hence we can apply lemma 2.9 and obtain the same conclusion. \square

LEMMA 2.11. – *Let ϕ be a real analytic diffeomorphism on \mathbb{R}^n such that $\phi(0) = 0$. Assume that $j_1 = 1$ in (2.1). If*

$$X \in H^{m,r}(\mathbb{R}^n) \implies T_\phi X \in H^{m,r}(\mathbb{R}^n), \tag{2.16}$$

for some $m \geq 0$, then

$$T_\phi Z = Z. \tag{2.17}$$

Proof. – By lemma 2.7 we know that (2.16) implies

$$[T_\phi X, Z] = (m - 1)T_\phi X \text{ for any } X \in H^{m,r}(\mathbb{R}^n). \tag{2.18}$$

Transforming both sides of (2.18) under the action of the map $T_{\phi^{-1}}$, which is the inverse map of T_ϕ , we obtain

$$[X, T_{\phi^{-1}} Z] = (m - 1)X \text{ for any } X \in H^{m,r}(\mathbb{R}^n),$$

which, together with (2.7), implies

$$[X, T_{\phi^{-1}}Z - Z] = 0, \text{ for any } X \in H^{m,r}(\mathbb{R}^n). \quad (2.19)$$

We now distinguish three cases:

Case 1. – Suppose $\delta_\varepsilon^r = \delta_\varepsilon^1$, $m = 0$. By lemma 2.10 it follows

$$T_{\phi^{-1}}Z = \sum_{i=1}^n a_i \frac{\partial}{\partial x_i} + Z, \quad a_i \in \mathbb{R}. \quad (2.20)$$

that can be explicitly written as

$$\left(\frac{\partial \phi^{-1}}{\partial x}(x) \right)^{-1} \begin{pmatrix} \phi_1^{-1}(x) \\ \vdots \\ \phi_n^{-1}(x) \end{pmatrix} = \begin{pmatrix} a_1 \\ \vdots \\ a_n \end{pmatrix} + \begin{pmatrix} x_1 \\ \vdots \\ x_n \end{pmatrix}.$$

If we evaluate the above at $x = 0$ and observe that $\phi(0) = 0$ implies $\phi^{-1}(0) = 0$, we obtain $a_i = 0, i = 1, \dots, n$. Hence $T_{\phi^{-1}}Z = Z$ which implies (2.17).

Case 2. – Suppose $m = 1$. From (2.19), using lemma 2.10, it follows

$$T_{\phi^{-1}}Z = (a + 1)Z, \quad (2.21)$$

for some $a \in \mathbb{R} \setminus \{-1\}$ ($a \neq -1$ since $T_{\phi^{-1}}Z = 0$ implies $Z = 0$, by applying T_ϕ to (2.21)). Set $b = (a + 1)^{-1}$. Then, using the linearity of T_ϕ , (2.21) implies $T_\phi Z = bZ$, $b \in \mathbb{R} \setminus \{0\}$, from which it follows

$$\begin{pmatrix} r_1 \phi_1(x) \\ \vdots \\ r_n \phi_n(x) \end{pmatrix} = \left(\frac{\partial \phi}{\partial x}(x) \right) \begin{pmatrix} b r_1 x_1 \\ \vdots \\ b r_n x_n \end{pmatrix}. \quad (2.22)$$

Since ϕ is analytic, its components ϕ_i can be expanded in terms of homogeneous polynomials; let $\bar{\phi}_i$ denote a monomial of ϕ_i of the form $\bar{\phi}_i(x) = \alpha x_1^{\nu_1} \cdot \dots \cdot x_n^{\nu_n}$. Then (2.22) implies

$$r_i \bar{\phi}_i(x) = \sum_{k=1}^n \frac{\bar{\phi}_i(x)}{x_k} \nu_k b r_k x_k = b \left(\sum_{k=1}^n \nu_k r_k \right) \bar{\phi}_i(x). \quad (2.23)$$

Thus (2.22) is equivalent to the condition that the exponents ν_1, \dots, ν_n of all the monomials $\bar{\phi}_i(x) = \alpha x_1^{\nu_1} \cdot \dots \cdot x_n^{\nu_n}$ of the same i -th component ϕ_i of ϕ , must satisfy the relation

$$r_i = b \left(\sum_{k=1}^n \nu_k r_k \right),$$

for each $i = 1, \dots, n$. This relation implies $b > 0$ and

$$\phi_i \in P^{r_i+l_i-1,r}(\mathbb{R}^n), \tag{2.24}$$

for each $i = 1, \dots, n$, with l_i integers satisfying

$$r_i = b(r_i + l_i - 1), \quad i = 1, \dots, n. \tag{2.25}$$

In particular (2.25), evaluated for $i = 1$, gives $l_1 = b^{-1}$ and thus, by substituting it back in (2.25), we have

$$r_i + l_i - 1 = r_i l_1, \quad i = 1, \dots, n. \tag{2.26}$$

Note that from the above it follows that l_1 is a positive integer since $b > 0$. Moreover, if $l_1 > 1$, from (2.24) and (2.26) we have

$$\phi_n \in P^{r_n l_1, r}(\mathbb{R}^n), \quad r_n l_1 > r_n,$$

which implies that the n -th component of ϕ is the sum of homogeneous terms of degree greater than one with respect to the standard dilation. But this would imply $\det((\partial\phi/\partial x)(0)) = 0$ which cannot be since ϕ is a diffeomorphism. Thus $l_1 = 1$ which implies $b = 1$ and $a = 0$ in (2.21), from which (2.17) follows.

Case 3. – Suppose $\delta_\varepsilon^r = \delta_\varepsilon^1, m > 1$, or $\delta_\varepsilon^r \neq \delta_\varepsilon^1, m \neq 1$. From (2.19), using lemma 2.10, it follows that $T_{\phi^{-1}}Z - Z = 0$ and thus (2.17). \square

COROLLARY 2.12. – *Let ϕ be a real analytic diffeomorphism on \mathbb{R}^n such that $\phi(0) = 0$. Then the following statements are equivalent:*

- (i) $T_\phi Z = Z$;
- (ii) $\phi \in H^{1,r}(\mathbb{R}^n)$.

Proof. – First suppose that (i) is satisfied. We have shown, in the proof of case 2 of lemma 2.11, that this condition, which is equivalent to (2.22) with $b = 1$, implies (2.24). Then, using (2.25) with $b = 1$, it follows that $\phi_i \in P^{r_i,r}(\mathbb{R}^n)$, for each $i = 1, \dots, n$ and so $\phi \in H^{1,r}(\mathbb{R}^n)$.

Next suppose that (ii) is satisfied. This condition implies that (2.23), in the proof of case 2 of lemma 2.11, is satisfied with $b = 1$ by any monomial $\bar{\phi}_i$ of each component ϕ_i of ϕ . It follows (2.22) with $b = 1$, which is equivalent to condition (i). \square

2.4. Proofs of the main results

Proof of Theorem 2.1. – We need only to observe that:

if $j_1 = 1$, condition (i) implies (ii) by lemma 2.11;
 if $j_1 \geq 1$, condition (ii) implies (i) by lemma 2.8;
 if $j_1 \geq 1$, condition (ii) is equivalent to (iii) by Corollary 2.12. \square

Proof of Theorem 2.2. – From the definition of homogeneity with respect to a dilation, it follows immediately that ϕ is a real analytic homogeneous function of degree one with respect to δ_ε^r if and only if it has the form (2.3) - (2.5), with A_k generic $(i_k - i_{k-1}) \times (i_k - i_{k-1})$ matrices with real entries. Thus we need only to show that, under the hypothesis $\phi \in H^{1,r}(\mathbb{R}^n)$, the matrices A_k in (2.4) are invertible if and only if ϕ is a diffeomorphism on \mathbb{R}^n .

One implication is obvious since if ϕ is a diffeomorphism then the Jacobian $(\partial\phi/\partial x)(0)$ is invertible and, from (2.3), it is clear that $(\partial\phi/\partial x)(0) = A$ which, having the form (2.4), is invertible if and only if each block A_k is invertible, $k = 1, \dots, m$.

Suppose now that A_k in (2.4) are invertible matrices for each $k = 1, \dots, m$. We will show that ϕ has a continuously differentiable inverse map on \mathbb{R}^n . Define the map

$$\psi : \mathbb{R}^n \rightarrow \mathbb{R}^n, \quad \psi(x) = A^{-1}x + h(x), \tag{2.27}$$

where h is a function whose i -th component is recursively defined by

$$h_i(x) = 0, \quad i = 1, \dots, i_1;$$

$$\begin{pmatrix} h_{i_{k-1}+1}(x) \\ \vdots \\ h_{i_k}(x) \end{pmatrix} = -A_k^{-1} \begin{pmatrix} g_{i_{k-1}+1} \left(\begin{pmatrix} A_1^{-1} & & \\ & \ddots & \\ & & A_{k-1}^{-1} \end{pmatrix} \begin{pmatrix} x_1 \\ \vdots \\ x_{i_{k-1}} \end{pmatrix} + \begin{pmatrix} h_1(x) \\ \vdots \\ h_{i_{k-1}}(x) \end{pmatrix} \right) \\ \vdots \\ g_{i_k} \left(\begin{pmatrix} A_1^{-1} & & \\ & \ddots & \\ & & A_{k-1}^{-1} \end{pmatrix} \begin{pmatrix} x_1 \\ \vdots \\ x_{i_{k-1}} \end{pmatrix} + \begin{pmatrix} h_1(x) \\ \vdots \\ h_{i_{k-1}}(x) \end{pmatrix} \right) \end{pmatrix} \tag{2.28}$$

$k = 2, \dots, m.$

Note that this definition makes sense since the fact that g is a strictly non linear map in $H^{1,r}(\mathbb{R}^n)$ implies that the components $g_i, i = i_{k-1} + 1, \dots, i_k$ are functions depending only on the previous i_{k-1} variables $(x_1, \dots, x_{i_{k-1}})$.

Thus it can be easily checked, by recursion on $k = 1, \dots, m$, that

$$\begin{pmatrix} (\phi \circ \psi)_{i_{k-1}+1}(x) \\ \vdots \\ (\phi \circ \psi)_{i_k}(x) \end{pmatrix} = \begin{pmatrix} (\psi \circ \phi)_{i_{k-1}+1}(x) \\ \vdots \\ (\psi \circ \phi)_{i_k}(x) \end{pmatrix} = \begin{pmatrix} x_{i_{k-1}+1} \\ \vdots \\ x_{i_k} \end{pmatrix}.$$

Hence $\psi = \phi^{-1}$ and, from definition (2.27), (2.28), it is clear that ψ is continuously differentiable, thus concluding the proof. \square

Proof of Theorem 2.3. – By Theorem 2.2 and Corollary 2.12 we know that $\phi \in GH^{1,r}(\mathbb{R}^n)$ if and only if ϕ is a diffeomorphism that fixes zero such that $T_\phi Z = Z$. Thus for $\phi \in GH^{1,r}(\mathbb{R}^n)$, from

$$Z = T_{id}Z = T_{\phi^{-1} \circ \phi}Z = T_{\phi^{-1}}T_\phi Z = T_{\phi^{-1}}Z,$$

it follows that the inverse ϕ^{-1} is an element of $GH^{1,r}(\mathbb{R}^n)$. Also we have, using the definition of homogeneity of degree one with respect to a dilation,

$$\phi \circ \psi(\delta_\varepsilon^r(x)) = \phi(\psi(\delta_\varepsilon^r(x))) = \phi(\delta_\varepsilon^r(\psi(x))) = \delta_\varepsilon^r(\phi(\psi(x))),$$

for any $\phi, \psi \in H^{1,r}(\mathbb{R}^n)$. This shows that $GH^{1,r}(\mathbb{R}^n)$ is a group. The fact that G_1, G_2 are subgroups of $GH^{1,r}(\mathbb{R}^n)$ is an immediate consequence of their definitions and of the form (2.27), (2.28) of the inverse map of an element in $GH^{1,r}(\mathbb{R}^n)$. Let now ϕ be an element of $GH^{1,r}(\mathbb{R}^n)$ of the form (2.3) - (2.5), $\phi(x) = Ax + g(x)$, and consider the following maps

$$\begin{aligned} \phi_1(x) &= Ax, \\ \phi_2(x) &= x + A^{-1}g(x), \\ \psi_2(x) &= x + g(A^{-1}x). \end{aligned}$$

Note that $A^{-1} \circ g$ and $g \circ A^{-1}$ are compositions of two elements of $H^{1,r}(\mathbb{R}^n)$ and therefore still elements of $H^{1,r}(\mathbb{R}^n)$ by the previous argument; also they are strictly non linear maps since g is. Hence $\phi_1 \in G_1$ and $\phi_2, \psi_2 \in G_2$ and the equality (2.6) is clearly satisfied.

Finally, to prove that $GH^{1,r}(\mathbb{R}^n)$ is a Lie group, observe that $H^{1,r}(\mathbb{R}^n)$ is a finite dimensional analytic manifold isomorphic to \mathbb{R}^{p_r} , with p_r a constant depending on $r = (r_1, \dots, r_n)$ (more precisely on the $2m$ -tuple $(i_1, \dots, i_m; j_1, \dots, j_m)$ associated to r in (2.1)). Note that $GH^{1,r}(\mathbb{R}^n)$ coincides with the open subset $\{ \phi \in H^{1,r}(\mathbb{R}^n) : \phi = A + g, \det A \neq 0 \}$ of $H^{1,r}(\mathbb{R}^n)$. Thus also $GH^{1,r}(\mathbb{R}^n)$ is a p_r dimensional analytic manifold and the group operation (the composition of maps) and the inversion ($i : \phi \rightarrow \phi^{-1}$) are clearly analytic. \square

Proof of Theorem 2.5. – We first observe that, by Theorem 4.3 in Section 4, the flow $(\exp tX)(x)$, $x \in \mathbb{R}^n$, generated by the field X is defined for all $t \in \mathbb{R}$. Let Z be the vector field defined in (2.2). It can be easily verified that

$$\delta_\varepsilon^r(x) = (\exp(\ln \varepsilon)Z)(x), \quad x \in \mathbb{R}^n, \quad \varepsilon > 0.$$

Since, by Lemma 2.7, $X \in H^{1,r}(\mathbb{R}^n)$ implies $[X, Z] = 0$, it follows that, for any fixed $t \in \mathbb{R}$,

$$\begin{aligned} \delta_\varepsilon^r((\exp tX)(x)) &= (\exp(\ln \varepsilon)Z) \circ (\exp tX)(x) \\ &= (\exp tX) \circ (\exp(\ln \varepsilon)Z)(x) \\ &= (\exp tX)(\delta_\varepsilon^r(x)), \quad x \in \mathbb{R}^n, \end{aligned}$$

which shows the homogeneity of the map $x \rightarrow (\exp tX)(x)$, t fixed. Moreover, it is well known from the theory of differential equations that the map $x \rightarrow (\exp tX)(x)$ is analytic and therefore we can conclude that it is an element of $H^{1,r}(\mathbb{R}^n)$. To show that $H^{1,r}(\mathbb{R}^n)$ is the Lie algebra of the Lie group $GH^{1,r}(\mathbb{R}^n)$, we observe that, by Theorem 2.2, $GH^{1,r}(\mathbb{R}^n)$ is the group of the invertible elements of the associative algebra $H^{1,r}(\mathbb{R}^n)$ (with the composition of functions as multiplication). Thus, the conclusion follows from a general result in the theory of Lie algebras (see [24, Section 2.3]). \square

3. DECOMPOSITION OF HOMOGENEOUS VECTOR FIELDS OF DEGREE ONE WITH RESPECT TO AN ARBITRARY DILATION

3.1. Proof of the main result

Before giving the proof of Theorem 3.1, we want to observe that the decomposition of vector fields given in this theorem would be an immediate consequence of a classical result in the theory of Lie algebras (the Jordan decomposition for elements of finite dimensional, semisimple Lie algebras: see [24, Thm. 3.10.6]), if $H^{1,r}(\mathbb{R}^n)$ or, at least $H^{1,r}(\mathbb{R}^n)/C(H^{1,r}(\mathbb{R}^n))$ ($C(H^{1,r}(\mathbb{R}^n))$ denoting the center of $H^{1,r}(\mathbb{R}^n)$), would have been semisimple. Nevertheless it can be easily seen that this is not the case for any dilation δ_ε^r different from the standard dilation δ_ε^1 .

Recall that a Lie algebra is said to be semisimple if it does not possess any non zero solvable ideal. From lemma 2.9 it follows that $C(H^{1,r}(\mathbb{R}^n)) = \{aZ(x) : a \in \mathbb{R}\}$. Thus, since the center is a non zero abelian (thus solvable) ideal, it is clear that $H^{1,r}(\mathbb{R}^n)$ is not semisimple.

Regarding $H^{1,r}(\mathbb{R}^n)/C(H^{1,r}(\mathbb{R}^n))$, if $\delta_\varepsilon^r \neq \delta_\varepsilon^1$ we can consider the subspace \mathcal{G} that is generated by the following set

$$\left\{ x^\nu \frac{\partial}{\partial x_i} + C(H^{1,r}(\mathbb{R}^n)) : x^\nu \frac{\partial}{\partial x_i} \in H^{j_m,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n) \right\}.$$

Note that, if $x^\alpha \partial/\partial x_s$ is an element in $H^{j,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$, $j \in \{j_1, \dots, j_m\}$, we have

$$\left[x^\nu \frac{\partial}{\partial x_i}, x^\alpha \frac{\partial}{\partial x_s} \right] \in H^{j+j_m-1,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n).$$

But, since $H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n) = \{0\}$, for any $k > j_m$, it follows that the above product, if not zero, must be an element of $H^{j_m,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$. Thus it can be easily seen that \mathcal{G} is an ideal of $H^{1,r}(\mathbb{R}^n)/C(H^{1,r}(\mathbb{R}^n))$. Moreover, observing again that the Lie product of two elements in $H^{j_m,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$ is an element of $H^{2j_m-1,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$, we deduce that their product is zero since $j_m > 1$ implies $2j_m - 1 > j_m$. Therefore \mathcal{G} is an abelian ideal of $H^{1,r}(\mathbb{R}^n)/C(H^{1,r}(\mathbb{R}^n))$ that cannot be semisimple.

Proof of Theorem 3.1. – Since $X \in H^{1,r}(\mathbb{R}^n)$, its expansion in homogeneous fields with respect to the standard dilation has the form

$$X(x) = Ax + \sum_{j=2}^{j_m} X^j(x), \quad X^j \in H^{j,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n). \quad (3.1)$$

We know that, if S and \widehat{N} denote the semisimple and nilpotent part of the linear field A , we have $[S, \widehat{N}] = 0$. Thus, in order to obtain the decomposition (1.3) satisfying (1.4), we will look for coordinate transformations that leave unchanged the linear part of X and remove, from the non linear part, terms of increasing degree not commuting with the semisimple field S . More precisely we will show, using induction, that there exist a finite sequence of transformations $\phi_k \in GH^{1,r}(\mathbb{R}^n)$, $1 \leq k \leq j_m$, satisfying

$$T_{\phi_1} X = X, \quad (3.2)$$

$$T_{\phi_k} T_{\phi_1 \circ \dots \circ \phi_{k-1}} X = \sum_{j=1}^{k-1} (T_{\phi_1 \circ \dots \circ \phi_{k-1}} X)^j + \sum_{j=k}^{j_m} Y^j,$$

$$[S, \sum_{j=2}^{k-1} (T_{\phi_1 \circ \dots \circ \phi_{k-1}} X)^j + Y^k] = 0, \quad k > 1, \quad (3.3)$$

where $(T_{\phi_1 \circ \dots \circ \phi_{k-1}} X)^j$ denotes the homogeneous part of degree j (with respect to the standard dilation) of the field X after performing the transformation $x = \phi_1 \circ \dots \circ \phi_{k-1}(y)$, and $Y^j \in H^{j,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$.

Since we know, from Theorem 2.3, that the set $GH^{1,r}(\mathbb{R}^n)$ is a group, the composition of these transformations is also an element of $GH^{1,r}(\mathbb{R}^n)$ and therefore their product $\phi = \phi_1 \circ \dots \circ \phi_{j_m}$ will produce a change of coordinates in $GH^{1,r}(\mathbb{R}^n)$ that gives X the form (1.3), (1.4), thus proving the theorem.

It is clear that the map $\phi_1 = I$ is an element of $GH^{1,r}(\mathbb{R}^n)$ that satisfies (3.2). Next suppose $\phi_s \in GH^{1,r}(\mathbb{R}^n)$, $1 \leq s < k$, $k > 1$, satisfying (3.3) have already been constructed. Since $\phi_s \in GH^{1,r}(\mathbb{R}^n)$, it follows from Theorem 2.1 that $T_{\phi_1 \circ \dots \circ \phi_s} X \in H^{1,r}(\mathbb{R}^n)$, for all $1 \leq s < k$. This, together with the inductive assumption, implies that $T_{\phi_1 \circ \dots \circ \phi_{k-1}} X$ is a field of the form

$$A + \sum_{j=2}^{j_m} \tilde{X}^j, \quad [S, \sum_{j=2}^{k-1} \tilde{X}^j] = 0, \quad \tilde{X}^j \in H^{j,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n). \quad (3.4)$$

Suppose that the field $T_{\phi_1 \circ \dots \circ \phi_{k-1}} X$ is given in the x -coordinates and consider the coordinate transformation $x = \phi_k(y)$, with

$$\phi_k = I + g^k, \quad g^k \in H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n). \quad (3.5)$$

The field (3.4), expressed in the new y -coordinates, becomes

$$\left(I + \frac{\partial g^k}{\partial y}(y) \right)^{-1} \left(Ay + Ag^k(y) + \sum_{j=2}^{j_m} \tilde{X}^j(y + g^k(y)) \right). \quad (3.6)$$

Since g^k is a strictly non linear map in $H^{1,r}(\mathbb{R}^n)$, we have $(\partial g_i^k / \partial y_s)(y) \equiv 0$, for all $1 \leq i \leq s \leq n$ (g_i^k denoting the i -th components of the map g^k) which implies that the Jacobian $(\partial g^k / \partial y)(y)$ is a lower triangular matrix with zeros on the diagonal and hence a nilpotent matrix. Therefore we can write

$$\left(I + \frac{\partial g^k}{\partial y}(y) \right)^{-1} = \sum_{s=0}^N (-1)^s \left(\frac{\partial g^k}{\partial y}(y) \right)^s,$$

for some $N \in \mathbb{N}$. Substitution of this into (3.6) and expansion of the resulting expression (using the analyticity of the fields \tilde{X}^j), retaining only

homogeneous terms of degree k (with respect to the standard dilation) and lower, produces

$$Ay + Ag^k(y) + \sum_{j=2}^k \tilde{X}^j(y) - \left(\frac{\partial g^k}{\partial y}(y) \right) Ay,$$

which can be rewritten in the form

$$Ay + \sum_{j=2}^k \tilde{X}^j(y) + [A, g^k](y). \tag{3.7}$$

Observe now that $g^k \in H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$ and $A \in H^{1,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$, imply $[A, g^k] \in H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$. Thus the map

$$g^k \longmapsto ad_k A(g^k) = [A, g^k], \tag{3.8}$$

is a linear operator acting on the linear space $H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$.

At this point we need to state a technical lemma that is proved after the theorem using a standard argument in normal form theory (e.g., see [22, Coroll. 2.1] and [23, Thm. 2.5]).

LEMMA 3.2. – *Let A be a linear map in $H^{1,r}(\mathbb{R}^n)$ and S the semisimple part of A . Denote by $Ker(ad_k S)$ and $Im(ad_k A)$ respectively the kernel and the range of the linear operators $ad_k S, ad_k A$ defined as in (3.8), with $k \geq 1$. Then there exists a subspace V_k of $Ker(ad_k S)$ that is a complement to $Im(ad_k A)$ in $H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$, i.e. such that*

$$H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n) = Im(ad_k A) \oplus V_k. \tag{3.9}$$

Using this lemma we can find $g^k, h^k \in H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$, such that

$$-\tilde{X}^k = ad_k A(g^k) + h^k, \quad [S, h^k] = 0. \tag{3.10}$$

Consequently, after the transformation (3.5) with g^k chosen as in (3.10), the field (3.4) takes the form (3.3), which concludes the proof. \square

Proof of Lemma 3.2. – Denote by \hat{N} the nilpotent part of A . We first show that $ad_k A = ad_k S + ad_k \hat{N}$ is the semisimple-nilpotent decomposition of the linear operator $ad_k A$. Since $[S, \hat{N}] = 0$ it follows, using the Jacobi identity, that $ad_k S$ and $ad_k \hat{N}$ commute. To prove that $ad_k S$ is semisimple, by considering the natural isomorphism between the complexification $(H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n))^c$ of the space $H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$,

and the space $H^{k,1}(\mathbb{C}^n) \cap H^{1,r}(\mathbb{C}^n)$ (which is defined in analogous way to the corresponding real one), it will be sufficient to show that $H^{k,1}(\mathbb{C}^n) \cap H^{1,r}(\mathbb{C}^n)$ has a basis consisting of eigenvectors of $ad_k S^c$ (where S^c denotes the complexification of the map S). Moreover, let $B \in H^{1,1}(\mathbb{C}^n) \cap H^{1,r}(\mathbb{C}^n)$ be the linear transformation that puts S^c in diagonal form $T_B S^c = B^{-1} S^c B = D$. Then, by Theorem 2.1., T_B is an invertible linear operator on $H^{k,1}(\mathbb{C}^n) \cap H^{1,r}(\mathbb{C}^n)$. Hence, it will be equivalent to show that there exists a basis of $H^{k,1}(\mathbb{C}^n) \cap H^{1,r}(\mathbb{C}^n)$ consisting of eigenvectors of $T_B \circ (ad_k S^c) \circ T_{B^{-1}} = ad_k T_B S^c = ad_k D$, where $D = diag\{\lambda_1, \dots, \lambda_n\}$ is a diagonal linear map with $\lambda_1, \dots, \lambda_n$ denoting the eigenvalues of S (and therefore of A). If we denote by $\{e_1, \dots, e_n\}$ the canonical basis of \mathbb{C}^n , the set

$$\mathcal{B} = \left\{ x^\nu e_i : 1 \leq i \leq n, \nu = (\nu_1, \dots, \nu_n) \in \mathbb{N}^n, \sum_{s=1}^n \nu_s = k, \sum_{s=1}^n \nu_s r_s = r_i \right\} \tag{3.11}$$

constitutes a basis for $H^{k,1}(\mathbb{C}^n) \cap H^{1,r}(\mathbb{C}^n)$. A straightforward computation shows that, for any $x^\nu e_i \in \mathcal{B}$, we have

$$ad_k D(x^\nu e_i) = \left(\lambda_i - \sum_{s=1}^n \nu_s \lambda_s \right) x^\nu e_i, \tag{3.12}$$

thus proving that the elements of \mathcal{B} are eigenvectors of $ad_k D$ with corresponding eigenvalues of the form

$$\begin{aligned} \mu_{\nu,i} &= \lambda_i - \sum_{s=1}^n \nu_s \lambda_s, \quad 1 \leq i \leq n, \nu = (\nu_1, \dots, \nu_n) \in \mathbb{N}^n, \\ &\sum_{s=1}^n \nu_s = k, \sum_{s=1}^n \nu_s r_s = r_i. \end{aligned} \tag{3.13}$$

Regarding the nilpotency of $ad_k \widehat{N}$, observe that for any $Z \in H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$ and for any $j > k$, we have $D^j Z = 0$ (denoting by $D^j Z$ the j -th order differential of Z). Then, since $(ad^m \widehat{N}, Z)(x)$ is a linear combination of elements of the form

$$\widehat{N}^s D^t Z(x) [\widehat{N}^{i_1} x, \dots, \widehat{N}^{i_t} x], \quad 0 \leq s, t \leq m, i_t \leq \dots \leq i_1,$$

with $t = i_1 = 0$ when $s = m$, and $1 \leq i_1, t + i_1 = m + 1 - s$ when $s < m$, it follows that $(ad_k \widehat{N})^{3p} = 0$ for $p > \max\{k, \bar{n}\}$, where \bar{n} denotes the nilpotency order of \widehat{N} .

Thus $ad_k S$ constitutes the semisimple part in the Jordan decomposition of the operator $ad_k A$, which implies that the kernel of the adjoint operator $(ad_k A)^T$ is a subspace of the kernel of $ad_k S$. Moreover, by simple facts from linear algebra (the Fredholm alternative) we know that the kernel of $(ad_k A)^T$ is a complement space to $Im(ad_k A)$. Thus $Ker((ad_k A)^T)$ is a possible choice for a subspace of $Ker(ad_k S)$ that satisfies (3.9). \square

Remark 3.3. – Denote by $C(H^{1,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n))$ the center of the Lie algebra $H^{1,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$. It is not difficult to verify that $(H^{1,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n))/C(H^{1,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n))$ is a finite dimensional semisimple Lie algebra. Therefore, since the map $A \mapsto ad_k A$ is a representation of $H^{1,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$ in $H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$, the semisimple-nilpotent decomposition $ad_k A = ad_k S + ad_k \widehat{N}$ in the proof of lemma 3.2 can also be derived from a well known result in the theory of Lie algebras (see [15, Coroll. 6.4]).

Remark 3.4. – Adapting results analogous to lemma 3.2, of A. Vanderbauwhede [21], and C. Elphick *et al.* [7], it can be easily shown that we can introduce a particular inner product on $H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$ such that the adjoint operator satisfies $(ad_k A)^T = ad_k A^T$. Thus we may choose the kernel of $ad_k A^T$ as a subspace V_k of $H^{k,1}(\mathbb{R}^n) \cap H^{1,r}(\mathbb{R}^n)$ that satisfies the condition of the lemma.

Example 3.5. – (i) On \mathbb{R}^3 consider the field

$$X(x) = x_1 \frac{\partial}{\partial x_1} - x_2 \frac{\partial}{\partial x_2} + (3x_3 + 5x_1^3 - x_1^2 x_2 + 2x_1 x_2^2 + 7x_2^3) \frac{\partial}{\partial x_3}.$$

Note that $X \in H^{1,r}(\mathbb{R}^3)$, with $r = (1, 1, 3)$, and that its linear part A is represented by a symmetric matrix. Therefore, by the previous remark, we can choose a complementary space to $Im(ad_3 A)$ satisfying (3.9) to be the kernel of $ad_3 A$, that is equal to $span\{x_1^3 \partial/\partial x_3\}$. In fact we have $Im(ad_3 A) = span\{x_1^2 x_2 \partial/\partial x_3, x_1 x_2^2 \partial/\partial x_3, x_2^3 \partial/\partial x_3\}$ and we may verify that, if we set

$$g^3(x) = \left(\frac{1}{2} x_1^2 x_2 - \frac{1}{2} x_1 x_2^2 - \frac{7}{6} x_2^3 \right) \frac{\partial}{\partial x_3}, \quad h^3(x) = 5x_1^3 \frac{\partial}{\partial x_3},$$

we can write $X = A - [A, g^3] + h^3, h^3 \in Ker(ad_3 A)$. Thus, after performing the coordinate transformation

$$x = \phi(y) = \begin{pmatrix} y_1 \\ y_2 \\ y_3 + (1/2)y_1^2 y_2 - (1/2)y_1 y_2^2 - (7/6)y_2^3 \end{pmatrix}$$

the field X takes the desired form (1.3), (1.4):

$$T_\phi X(y) = \left(y_1 \frac{\partial}{\partial y_1} - y_2 \frac{\partial}{\partial y_2} + 3y_3 \frac{\partial}{\partial y_3} \right) + 5y_1^3 \frac{\partial}{\partial y_3},$$

$$\left[y_1 \frac{\partial}{\partial y_1} - y_2 \frac{\partial}{\partial y_2} + 3y_3 \frac{\partial}{\partial y_3}, 5y_1^3 \frac{\partial}{\partial y_3} \right] = 0.$$

(ii) On \mathbb{R}^4 consider the field

$$X(x) = -x_2 \frac{\partial}{\partial x_1} + x_1 \frac{\partial}{\partial x_2} + (-x_4 + 5x_1^3 - 6x_1^2 x_2 + 5x_1 x_2^2 + x_2^3) \frac{\partial}{\partial x_3}$$

$$+ (x_3 - 2x_1^3 + 3x_1^2 x_2 + 3x_1 x_2^2 + 11x_2^3) \frac{\partial}{\partial x_4}.$$

Note that $X \in H^{1,r}(\mathbb{R}^4)$, with $r = (1, 1, 3, 3)$, and that its linear part A is a semisimple field with double, purely imaginary eigenvalues $\lambda_{1,2} = \pm i$. An easy calculation shows that the kernel of $ad_3 A^T$ is equal to $\text{span}\{(x_1^3 + x_1 x_2^2) \partial / \partial x_3 + (x_1^2 x_2 + x_2^3) \partial / \partial x_4, (x_1^2 x_2 + x_2^3) \partial / \partial x_3 - (x_1^3 + x_1 x_2^2) \partial / \partial x_4\}$, and that the range of $ad_3 A$ is equal to $\text{span}\{3x_1^2 x_2 \partial / \partial x_3 + x_1^3 \partial / \partial x_4, (-x_1^3 + x_1 x_2^2) \partial / \partial x_3 + x_1^2 x_2 \partial / \partial x_4, (-2x_1^2 x_2 + x_2^3) \partial / \partial x_3 + x_1 x_2^2 \partial / \partial x_4, -3x_1 x_2^2 \partial / \partial x_3 + x_2^3 \partial / \partial x_4, -x_1^3 \partial / \partial x_3 + 3x_1^2 x_2 \partial / \partial x_4, -x_2^3 \partial / \partial x_3 - 3x_1 x_2^2 \partial / \partial x_4\}$. Let $X^3 = (5x_1^3 - 6x_1^2 x_2 + 5x_1 x_2^2 + x_2^3) \partial / \partial x_3 + (-2x_1^3 + 3x_1^2 x_2 + 3x_1 x_2^2 + 11x_2^3) \partial / \partial x_4$ be the homogeneous part of degree 3 (with respect to the standard dilation) of X . Then, by the previous remark, we can write $X^3 = -[A, g^3] + h^3$, with $g^3 \in H^{3,1}(\mathbb{R}^4) \cap H^{1,r}(\mathbb{R}^4)$, and $h^3 \in \text{Ker}(ad_3 A^T)$. Indeed, this equality is satisfied with

$$g^3(x) = (2x_1^3 - 5x_1^2 x_2 - 4x_2^3) \frac{\partial}{\partial x_3} + (3x_1^3 + x_2^3) \frac{\partial}{\partial x_4},$$

$$h^3(x) = (7x_1^3 + 7x_1 x_2^2) \frac{\partial}{\partial x_3} + (7x_1^2 x_2 + 7x_2^3) \frac{\partial}{\partial x_4}.$$

Therefore, the coordinate transformation

$$x = \phi(y) = \begin{pmatrix} y_1 \\ y_2 \\ y_3 + 2y_1^3 - 5y_1^2 y_2 - 4y_2^3 \\ y_4 + 3y_1^3 + y_2^3 \end{pmatrix}$$

gives to the field X the form (1.3), (1.4):

$$\begin{aligned}
 T_\phi X(y) &= -y_2 \frac{\partial}{\partial y_1} + y_1 \frac{\partial}{\partial y_2} + (-y_4 + 7y_1^3 + 7y_1y_2^2) \frac{\partial}{\partial y_3} \\
 &\quad + (y_3 + 7y_1^2y_2 + 7y_2^3) \frac{\partial}{\partial y_4}, \\
 \left[-y_2 \frac{\partial}{\partial y_1} + y_1 \frac{\partial}{\partial y_2} - y_4 \frac{\partial}{\partial y_3} + y_3 \frac{\partial}{\partial y_4}, (7y_1^3 + 7y_1y_2^2) \frac{\partial}{\partial y_3} \right. \\
 &\quad \left. + (7y_1^2y_2 + 7y_2^3) \frac{\partial}{\partial y_4} \right] = 0.
 \end{aligned}$$

Remark 3.6. – The constructive proof of Theorem 3.1 permits one to compute a normal form for X satisfying (1.3), (1.4) in a finite number of steps $j \leq j_m$. Such a procedure cannot be applied to vector fields that are not homogeneous of degree one with respect to a dilation δ_ε^r . In fact in this latter case, since we cannot use Theorem 2.1, any transformation ϕ_k of the form (3.5) introduces additional higher order terms that are no more necessarily homogeneous of degree one with respect to the dilation δ_ε^r . Thus we would produce an infinite sequence $\{\phi_k\}_k, k \in \mathbb{N}$, and it would be necessary to study the convergence to a function of the corresponding products $\{\phi_1 \circ \dots \circ \phi_k\}_k$.

Remark 3.7. – Theorem 3.1 can also be derived from a result of K.T. Chen [6, Thm. 8.1] who proved the existence of a formal transformation (i.e. a transformation given by a formal series) that puts a formal vector field of a graded Lie algebra into the form (1.3), (1.4). Our approach is a constructive treatment of the formal result of Chen and yields an explicit procedure to compute a polynomial coordinate change that gives to a field $X \in H^{1,r}(\mathbb{R}^n)$ the normal form (1.3), (1.4).

Remark 3.8. – The decomposition (1.3), (1.4) of a field $X \in H^{1,r}(\mathbb{R}^n)$ produces a normal form for all the elements of the group orbit $T_\phi X, \phi \in GH^{1,r}(\mathbb{R}^n)$. This normal form is not unique. Indeed its linear part can be uniquely determined if we require it to be in Jordan canonical form (or in real canonical form in the case it has some non real eigenvalue) which can be done by performing, after obtaining (1.3), a further coordinate linear transformation $y = \psi(z), \psi \in GH^{1,r}(\mathbb{R}^n)$, that clearly preserves (1.4). However, the choice of the complementary spaces V_k to $Im(ad_k A)$ that satisfy condition (3.9) of the lemma, is in general far from being unique and therefore some arbitrariness in the form of the non linear part is unavoidable.

3.2. An application

We present here a corollary of Theorem 3.1 that provides a Poincaré type result for homogeneous vector fields of degree one with respect to a dilation δ_ε^r , $r = (r_1, \dots, r_n)$ of the type $(i_1, \dots, i_m; j_1, \dots, j_m)$ as defined in (2.1). In fact a well known theorem proved by Poincaré shows that a sufficient condition for the existence of a (formal) coordinate change that takes a given vector field X into a linear field is that the eigenvalues $\lambda_1, \dots, \lambda_n$ of the linear part of X do not satisfy any resonance relation of order $k \geq 2$, i.e. any relation of the form

$$\lambda_i = \sum_{s=1}^n \nu_s \lambda_s, \quad \nu_s \in \mathbb{N}, \quad \sum_{s=1}^n \nu_s = k.$$

For vector fields $X \in H^{1,r}(\mathbb{R}^n)$ it turns out to be sufficient to check that the eigenvalues do not satisfy only a finite number of resonance relations.

COROLLARY 3.9. – *Let X be a real analytic vector field in $H^{1,r}(\mathbb{R}^n)$. Denote by*

$$\lambda_1 = \lambda_{i_0}, \dots, \lambda_{i_1}, \lambda_{(i_1+1)}, \dots, \lambda_{i_2}, \dots, \lambda_{(i_{m-1}+1)}, \dots, \lambda_{i_m} = \lambda_n,$$

the eigenvalues of the linear part $A = \text{diag}\{A_1, \dots, A_m\}$ of X ordered so that $\lambda_{(i_{j-1}+1)}, \dots, \lambda_{i_j}$, $j = 1, \dots, m$, are the eigenvalues of the j -th block A_j of A . Suppose that $\lambda_1, \dots, \lambda_n$ do not satisfy any relation of the form

$$\lambda_i = \sum_{s=1}^n \nu_s \lambda_s, \tag{3.14}$$

with $\nu_s \in \mathbb{N}$ such that

$$r_i = \sum_{s=1}^n \nu_s r_s, \quad \sum_{s=1}^n \nu_s = k, \quad \text{for some } k \ (2 \leq k \leq j_m). \tag{3.15}$$

Then there exists a polynomial change of coordinates $\phi \in GH^{1,r}(\mathbb{R}^n)$ such that

$$T_\phi X = A. \tag{3.16}$$

A relation among the eigenvalues of the form (3.14), with $\nu_s \in \mathbb{N}$ satisfying (3.15), will be called a resonance with respect to the dilation δ_ε^r .

Proof. – The proof of lemma 3.2 shows that the semisimple part of each operator $ad_k A$ is given by the operator $ad_k S$ (with S being the semisimple part of A) and that the eigenvalues of $ad_k S$ have the form (3.13) $\mu = \lambda_i - \sum_{s=1}^n \nu_s \lambda_s$, with $\sum_{s=1}^n \nu_s r_s = r_i$, $\sum_{s=1}^n \nu_s = k$. Since the non resonance relations with respect to δ_ε^r of the eigenvalues $\lambda_1, \dots, \lambda_n$ imply that the eigenvalues of $ad_k S$ (and therefore of $ad_k A$) are all different from zero, the operators $ad_k A$ are invertible for $2 \leq k \leq j_m$. Hence the transformation $\phi = \phi_2 \circ \dots \circ \phi_{j_m}$ constructed in Theorem 3.1 gives X the form (3.16). \square

Example 3.10. – (i) On \mathbb{R}^3 consider the vector field

$$X(x) = (2x_1 + x_2) \frac{\partial}{\partial x_1} + 2x_2 \frac{\partial}{\partial x_2} + (x_1^2 + 2x_1x_2 - x_2^2) \frac{\partial}{\partial x_3}.$$

Observe that $X \in H^{1,r}(\mathbb{R}^n)$, with $r = (1, 1, 2)$. The eigenvalues of the linear part of X are $\lambda_1 = \lambda_2 = 2, \lambda_3 = 0$. Note that there is no resonance of order 2 with respect to δ_ε^r :

$$\lambda_3 \neq 2\lambda_1, \quad \lambda_3 \neq \lambda_1 + \lambda_2, \quad \lambda_3 \neq 2\lambda_2.$$

Therefore, by Corollary 3.9, we can find a coordinate change $\phi \in GH^{1,r}(\mathbb{R}^n)$ which transforms the field X into its linear part. In fact, if we set

$$x = \phi(y) = \begin{pmatrix} y_1 \\ y_2 \\ y_3 + (1/4)y_1^2 + (3/80)y_1y_2 - (11/32)y_2^2 \end{pmatrix},$$

we obtain

$$T_\phi X(y) = (2y_1 + y_2) \frac{\partial}{\partial y_1} + 2y_2 \frac{\partial}{\partial y_2}.$$

It may be noted that it was possible to transform the field X into its linear part even if there were resonances of any order of the form (3.14) not satisfying (3.15) (*i.e.* resonances in the classical sense) due to the eigenvalue $\lambda_3 = 0$ and in particular there were the following resonances of order 2:

$$\lambda_1 = \lambda_2 + \lambda_3, \quad \lambda_2 = \lambda_1 + \lambda_3.$$

(ii) On \mathbb{R}^2 consider the vector field

$$X(x) = (x_1 + x_2^2) \frac{\partial}{\partial x_1} + (-x_2 + x_1^2) \frac{\partial}{\partial x_2}.$$

The eigenvalues of the linear part are $\lambda_1 = 1, \lambda_2 = -1$. Observe that there is no classical resonance of order 2 since the eigenvalues do not satisfy any relation of the form (3.14) with $\nu_1 + \nu_2 = 2$:

$$\lambda_1, \lambda_2 \neq 0, \quad \lambda_1 \neq 2\lambda_2, \quad \lambda_2 \neq 2\lambda_1.$$

Thus we can find a coordinate transformation ϕ which removes completely the quadratic terms. If we set

$$x = \phi(y) = \begin{pmatrix} y_1 - (1/3)y_2^2 \\ y_2 + (1/3)y_1^2 \end{pmatrix},$$

we obtain

$$T_\phi X(y) = \begin{pmatrix} y_1 + \frac{2}{3}y_1^2y_2 + o(|y_1|^3, |y_2|^3) \\ -y_2 - \frac{2}{3}y_1y_2^2 + o(|y_1|^3, |y_2|^3) \end{pmatrix} \begin{pmatrix} \frac{\partial}{\partial y_1} \\ \frac{\partial}{\partial y_2} \end{pmatrix}.$$

Note that $(2/3)y_1^2y_2 \partial/\partial y_1, -(2/3)y_1y_2^2 \partial/\partial y_2$ are resonant terms of order 3 in the classical sense: the eigenvalues satisfy the following relations of the form (3.14)

$$\lambda_1 = 2\lambda_1 + \lambda_2, \quad \lambda_2 = \lambda_1 + 2\lambda_2.$$

Thus these terms cannot be removed by further coordinate transformations. In this case the field X is not homogeneous of degree one with respect to any dilation. Therefore the lack of classical resonances of order 2 is not sufficient to guarantee the existence of a coordinate change that transforms X into its linear part.

4. REPRESENTATION OF SOLUTIONS FOR A CLASS OF NONLINEAR SYSTEMS

In this section we first derive a general representation formula for the solutions of an autonomous system of differential equations

$$\dot{x} = X(x), \quad X \in H^{1,r}(\mathbb{R}^n). \quad (4.1)$$

Then we use this result to obtain a representation for the trajectories of an n -dimensional, single input, affine control system

$$\begin{aligned} \dot{x} &= X_0(x) + Bu, \quad X_0 \in H^{1,r}(\mathbb{R}^n), \\ B &= \begin{pmatrix} b_1 \\ \vdots \\ b_n \end{pmatrix} \in \mathbb{R}^n, \quad b_i = 0, \quad i = 1, \dots, i_{m-1}, \end{aligned} \quad (4.2)$$

($i_{m-1} + 1$ denoting as in (2.1) the smallest index i such that $r_i = r_n$) in terms of integrals of the control u , which is analogous to the standard representation of the trajectories of a linear control system $\dot{x} = Ax + bu$, (A, b , denoting $n \times n, n \times 1$ matrices). In fact, a system of the form (4.2) coincides with a linear system when δ_ε^r is the standard dilation δ_ε^1 with $r_1 = \dots = r_n = 1$. Thus, systems of the form (4.2) can be regarded as a natural extension of the linear control systems.

4.1. Representation of the flow generated by a homogeneous field of degree one with respect to an arbitrary dilation

THEOREM 4.1. – *Let N be a vector field in $H^{1,r}(\mathbb{R}^n)$ of the form*

$$N = \widehat{N} + Y, \quad \widehat{N}, Y \in H^{1,r}(\mathbb{R}^n), \quad (4.3)$$

where \widehat{N} is a linear nilpotent field and Y is a strictly non linear field. If $\phi_k(t, p)$, $k \in \mathbb{N}$, denotes the Picard approximations of the solution to $\dot{x} = N(x), x(0) = p$, recursively defined by

$$\begin{aligned} \phi_0(t, p) &= p, \\ \phi_{k+1}(t, p) &= p + \int_0^t N(\phi_k(s, p)) ds, \quad k = 0, 1, \dots, \end{aligned}$$

then there exists an integer \bar{k} such that

$$(\exp tN)(p) = \phi_{\bar{k}}(t, p), \quad t \in \mathbb{R}, \quad p \in \mathbb{R}^n. \quad (4.4)$$

Proof. – Denote by $\phi_{k,i}$ the i -th component of the k -th Picard approximation ϕ_k . We will show, by induction on the index i , $1 \leq i \leq n$, that there exists a non decreasing n -tuple of positive integers $k_1 \leq k_2 \leq \dots \leq k_n$, such that

$$\phi_{k,i}(t, p) = \phi_{k+1,i}(t, p), \quad t \in \mathbb{R}, \quad p \in \mathbb{R}^n, \quad \forall k \geq k_i. \quad (4.5)$$

This implies that, if we set $\bar{k} = k_n$, we have

$$(\exp tN)(p) = \lim_{k \rightarrow \infty} \phi_k(t, p) = \phi_{\bar{k}}(t, p), \quad t \in \mathbb{R}, \quad p \in \mathbb{R}^n,$$

which is what we need to prove.

Since \widehat{N} is a linear nilpotent element of $H^{1,r}(\mathbb{R}^n)$, it follows that (with the usual notation of Section 2) $\widehat{N} = \text{diag}\{\widehat{N}_1, \dots, \widehat{N}_m\}$, where each \widehat{N}_j is a $(i_j - i_{j-1}) \times (i_j - i_{j-1})$ nilpotent block. Denote by $n_j, j = 1, \dots, m$, the nilpotency order of \widehat{N}_j , i.e. the smallest positive integer n such that $\widehat{N}_j^n = 0$. Note that, by Theorem 2.5, ϕ_k are homogeneous functions of degree one with respect to δ_x^r , in the p -variable. From the definition of ϕ_k and the general form of an element of $H^{1,r}(\mathbb{R}^n)$, it follows that, if we set $k_1 = k_2 = \dots = k_{i_1} = n_1$, we have

$$\begin{aligned} \begin{pmatrix} \phi_{k,1}(t, p) \\ \vdots \\ \phi_{k,i_1}(t, p) \end{pmatrix} &= \sum_{s=0}^k \frac{t^s}{s!} \widehat{N}_1^s \begin{pmatrix} p_1 \\ \vdots \\ p_{i_1} \end{pmatrix} = \sum_{s=0}^{k+1} \frac{t^s}{s!} \widehat{N}_1^s \begin{pmatrix} p_1 \\ \vdots \\ p_{i_1} \end{pmatrix} \\ &= \begin{pmatrix} \phi_{k+1,1}(t, p) \\ \vdots \\ \phi_{k+1,i_1}(t, p) \end{pmatrix}, \quad t \in \mathbb{R}, \quad p \in \mathbb{R}^n, \quad \forall k \geq k_i, \quad i = 1, \dots, i_1. \end{aligned}$$

Suppose now that we have obtained $k_1 \leq k_2 \leq \dots \leq k_{i_j}, 1 \leq j < m$, such that relation (4.5) is satisfied for $i = 1, \dots, i_j$. Then

$$\phi_{k,i}(t, p) = \phi_{k+1,i}(t, p), \quad t \in \mathbb{R}, \quad p \in \mathbb{R}^n, \quad i = 1, \dots, i_j, \quad k \geq k_{i_j}. \tag{4.6}$$

Moreover, it can be easily verified that, for any integer $k > 1$, the following equality holds

$$\begin{aligned} \begin{pmatrix} \phi_{k,i_j+1}(t, p) \\ \vdots \\ \phi_{k,i_{(j+1)}}(t, p) \end{pmatrix} &= \sum_{s=0}^k \frac{t^s}{s!} \widehat{N}_{j+1}^s \begin{pmatrix} p_{i_j+1} \\ \vdots \\ p_{i_{(j+1)}} \end{pmatrix} + \frac{t^k}{k!} \widehat{N}_{j+1}^{k-1} \begin{pmatrix} Y_{i_j+1}(p) \\ \vdots \\ Y_{i_{(j+1)}}(p) \end{pmatrix} \\ &+ \int_0^t \int_0^{t_{k-2}} \int_0^{t_{k-3}} \dots \int_0^{t_1} \widehat{N}_{j+1}^{k-2} \begin{pmatrix} Y_{i_j+1}(\phi_1(v, p)) \\ \vdots \\ Y_{i_{(j+1)}}(\phi_1(v, p)) \end{pmatrix} dv dt_1 \dots dt_{k-2} \\ &+ \dots + \int_0^t \begin{pmatrix} Y_{i_j+1}(\phi_{k-1}(v, p)) \\ \vdots \\ Y_{i_{(j+1)}}(\phi_{k-1}(v, p)) \end{pmatrix} dv. \end{aligned}$$

Hence, using relation (4.5) and nilpotency of \widehat{N}_{j+1} , it follows that, if we set $k_{i_j+1} = k_{i_j+2} = \dots = k_{i_{(j+1)}} = k_{i_j} + n_{j+1}$, we have

$$\begin{aligned} \begin{pmatrix} \phi_{k_{i_j+1}}(t, p) \\ \vdots \\ \phi_{k_{i_{(j+1)}}}(t, p) \end{pmatrix} &= \sum_{s=0}^{n_{(j+1)}-1} \frac{t^s}{s!} \widehat{N}_{j+1}^s \begin{pmatrix} p_{i_j+1} \\ \vdots \\ p_{i_{(j+1)}} \end{pmatrix} \\ + \int_0^t \int_0^{t_{n_{(j+1)}-1}} \dots \int_0^{t_1} \widehat{N}_{j+1}^{n_{(j+1)}-1} &\begin{pmatrix} Y_{i_j+1}(\phi_{k_{i_j}}(v, p)) \\ \vdots \\ Y_{i_{(j+1)}}(\phi_{k_{i_j}}(v, p)) \end{pmatrix} dv dt_1 \dots dt_{n_{(j+1)}-1} \\ &+ \dots + \int_0^t \begin{pmatrix} Y_{i_j+1}(\phi_{k_{i_j}}(v, p)) \\ \vdots \\ Y_{i_{(j+1)}}(\phi_{k_{i_j}}(v, p)) \end{pmatrix} dv, \quad \forall k \geq k_{i_{(j+1)}}, \end{aligned}$$

which proves the inductive step. □

Remark 4.2. – It is clear, from the proof of Theorem 4.1, that the smallest integer \bar{k} such that (4.4) is satisfied, is given by $\bar{k} = \sum_{j=1}^m n_j$, where n_j denote the nilpotency orders of the blocks \widehat{N}_j of the linear part of the field N . Thus the solution to the Cauchy problem $\dot{x} = N(x)$, $x(0) = p$, with $N = \widehat{N} + Y \in H^{1,r}(\mathbb{R}^n)$ as in Theorem 4.1, can be computed using the following recursive formula

$$\begin{aligned} \begin{pmatrix} ((\exp tN)(p))_1 \\ \vdots \\ ((\exp tN)(p))_{i_1} \end{pmatrix} &= \sum_{s=0}^{n_1-1} \frac{t^s}{s!} \widehat{N}_1^s \begin{pmatrix} p_1 \\ \vdots \\ p_{i_1} \end{pmatrix}, \\ \begin{pmatrix} ((\exp tN)(p))_{i_j+1} \\ \vdots \\ ((\exp tN)(p))_{i_{(j+1)}} \end{pmatrix} &= \sum_{s=0}^{n_{(j+1)}-1} \frac{t^s}{s!} \widehat{N}_{j+1}^s \begin{pmatrix} p_{i_j+1} \\ \vdots \\ p_{i_{(j+1)}} \end{pmatrix} \\ + \int_0^t \int_0^{t_{n_{(j+1)}-1}} \dots &\dots \int_0^{t_1} \widehat{N}_{j+1}^{n_{(j+1)}-1} \begin{pmatrix} Y_{i_j+1}((\exp vN)(p)) \\ \vdots \\ Y_{i_{(j+1)}}((\exp vN)(p)) \end{pmatrix} dv dt_1 \dots dt_{n_{(j+1)}-1} \\ + \dots + \int_0^t &\begin{pmatrix} Y_{i_j+1}((\exp vN)(p)) \\ \vdots \\ Y_{i_{(j+1)}}((\exp vN)(p)) \end{pmatrix} dv, \quad j = 1, \dots, m-1. \end{aligned} \tag{4.7}$$

THEOREM 4.3. – Let X be a real analytic vector field in $H^{1,r}(\mathbb{R}^n)$, and $\phi \in GH^{1,r}(\mathbb{R}^n)$ a change of coordinates that transforms X into the canonical form $T_\phi X = S + N$ of Theorem 3.1. Then

$$(\exp tX)(p) = \phi(e^{tS}(\exp tN)(\phi^{-1}(p))), \quad t \in \mathbb{R}, \quad p \in \mathbb{R}^n, \quad (4.8)$$

where the flow $(\exp tN)(\phi^{-1}(p))$ can be computed by the recursive formula (4.7).

Proof. – Note that

$$(\exp tX)(p) = \phi((\exp tT_\phi X)(\phi^{-1}(p))). \quad (4.9)$$

Since $[S, N] = 0$ we can write

$$(\exp tT_\phi X)(q) = (\exp t(S + N))(q) = (\exp tS) \circ (\exp tN)(q),$$

which evaluated at $q = \phi^{-1}(p)$, gives (4.8) after applying ϕ and using (4.9). \square

Remark 4.4. – It may be noted that the formula (4.8) in Theorem 4.4, is analogous to the “generalized variation of constants formula”, obtained by A.A Agrachev, R.V. Gamkrelidze and A.V. Sarychev in [1], that express the flow generated by a vector field $X_\tau + Y_\tau$ as a perturbation of the flow generated by X_τ .

Example 4.5. – Consider the Cauchy problem on \mathbb{R}^4

$$\dot{x} = X(x), \quad x(0) = p = (p_1, p_2, p_3, p_4),$$

where $X \in H^{1,r}(\mathbb{R}^4)$, $r = (1, 1, 2, 2)$, is the vector field

$$X(x) = (2x_1 + x_2) \frac{\partial}{\partial x_1} + 2x_2 \frac{\partial}{\partial x_2} + (4x_3 + 5x_1^2 - x_1x_2 + x_2^2) \frac{\partial}{\partial x_3} + (-x_4 + x_1x_2 - 3x_2^2) \frac{\partial}{\partial x_4}.$$

First we compute the coordinate transformation that puts X into the canonical form (1.3),(1.4). If A denotes the linear part of X , and S the semisimple part of A , it can be easily seen that $(5x_1^2 - x_1x_2 + x_2^2) \partial/\partial x_3 \in \text{Ker}(ad_2 S)$, and that $\text{span}\{x_1x_2 \partial/\partial x_4, x_2^2 \partial/\partial x_4\}$ is a subspace of $\text{Im}(ad_2 A)$. Indeed we may check that, if we set $g^2(x) = ((1/5)x_1x_2 - (16/25)x_2^2) \partial/\partial x_4$, $h^2(x) = (5x_1^2 - x_1x_2 + x_2^2) \partial/\partial x_3$, we have

$X = A - [A, g^2] + h^2$. Thus, the coordinate change $x = \phi(y), \phi = I + g^2$, transforms X into the canonical form

$$T_\phi X = S + N, \quad N(y) = y_2 \frac{\partial}{\partial y_1} + (5y_1^2 - y_1 y_2 + y_2^2) \frac{\partial}{\partial y_3}.$$

Next, using formula (4.7) with $q = (q_1, q_2, q_3, q_4)$, we compute the flow generated by N

$$\begin{aligned} \begin{pmatrix} ((\exp tN)(q))_1 \\ ((\exp tN)(q))_2 \end{pmatrix} &= \begin{pmatrix} q_1 \\ q_2 \end{pmatrix} + t \begin{pmatrix} q_2 \\ 0 \end{pmatrix} \\ \begin{pmatrix} ((\exp tN)(q))_3 \\ ((\exp tN)(q))_4 \end{pmatrix} &= \begin{pmatrix} q_3 \\ q_4 \end{pmatrix} + \int_0^t \begin{pmatrix} 5(q_1 + vq_2)^2 - (q_1 + vq_2)q_2 + q_2^2 \\ 0 \end{pmatrix} dv \\ &= \begin{pmatrix} q_3 + (5q_1^2 - q_1q_2 + q_2^2)t + (5q_1q_2 - (1/2)q_2^2)t^2 + (5/3)q_2^2t^3 \\ q_4 \end{pmatrix}. \end{aligned}$$

Then, multiplying the exponential matrix e^{tS} by the above expression, after substituting in it $q = \phi^{-1}(p) = p - g^2(p)$, we obtain

$$(\exp tT_\phi X)(\phi^{-1}(p)) = \begin{pmatrix} e^{2t}(p_1 + tp_2) \\ e^{2t}p_2 \\ e^{4t}(p_3 + (5p_1^2 - p_1p_2 + p_2^2)t + (5p_1p_2 - (1/2)p_2^2)t^2 + (5/3)p_2^2t^3) \\ e^{-t}(p_4 - (1/5)p_1p_2 + (16/25)p_2^2) \end{pmatrix},$$

which gives

$$(\exp tX)(p) = \begin{pmatrix} e^{2t}(p_1 + tp_2) \\ e^{2t}p_2 \\ e^{4t}(p_3 + (5p_1^2 - p_1p_2 + p_2^2)t + (5p_1p_2 - (1/2)p_2^2)t^2 + (5/3)p_2^2t^3) \\ e^{-t}(p_4 - (1/5)p_1p_2 + (16/25)p_2^2) \\ + e^{4t}(((1/5)p_1p_2 - (16/25)p_2^2) + (1/5)p_2^2t) \end{pmatrix}.$$

4.2. Representation of the trajectories for a class of nonlinear control systems.

THEOREM 4.6. – *Let $\dot{x} = X_0(x) + Bu$, be an n -dimensional, single input, affine control system as in (4.2). Then, if we denote by A the linear part of the field X_0 , the trajectories $x(\cdot, u)$ of such a system can be computed*

with the formula

$$x(t, u) = (\exp tX_0) \left(x(0) + \int_0^t e^{-sA} B u(s) ds \right), \quad t \in \mathbb{R}. \quad (4.10)$$

Proof. – A direct calculation verifies that in order to show that (4.10) represents a trajectory of (4.2) we only need to prove

$$\left(\frac{\partial(\exp tX_0)}{\partial p}(p) \Big|_{p=x(0)+\int_0^t e^{-sA} B u(s) ds} \right) e^{-tA} B = B, \quad t \in \mathbb{R}. \quad (4.11)$$

Since X_0 is homogeneous of degree one with respect to δ_ε^r , its linear part is represented by a diagonal matrix composed of blocks of order $(i_k - i_{k-1}) \times (i_k - i_{k-1})$, $k = 1, \dots, m$ (with notation of Section 2) and therefore also the matrix representing e^{-tA} has the same form. It follows that $e^{-tA} B$ is a column vector having all components zero but the last $n - i_{m-1}$ like B . Hence in order to prove (4.11) it is sufficient to show

$$\left(\frac{\partial(\exp tX_0)}{\partial p}(p) \right) C = e^{tA} C, \quad p \in \mathbb{R}^n, \quad t \in \mathbb{R}, \quad (4.12)$$

for any column vector C having all components zero but the last $n - i_{m-1}$. It is a well-known fact, in the theory of differential equations, that the Jacobian $(\partial(\exp tX_0)/\partial p)(p)$ is equal to the fundamental matrix solution of the variational equation

$$\dot{v} = \left(\frac{\partial X_0}{\partial x}(x) \Big|_{x=(\exp tX_0)(p)} \right) v = Av + \left(\frac{\partial Y}{\partial x}(x) \Big|_{x=(\exp tX_0)(p)} \right) v, \quad (4.13)$$

where Y denotes the nonlinear part of X_0 . Let $\tilde{v}_j, j = 1, \dots, n - i_{m-1}$ be the solutions of $\dot{\tilde{v}} = A_m \tilde{v}$, $\tilde{v}(0) = \tilde{e}_j$, where $\{\tilde{e}_1, \dots, \tilde{e}_{n-i_{m-1}}\}$ denotes the canonical basis of $\mathbb{R}^{n-i_{m-1}}$ and A_m is the last block of the matrix A . Then define the functions $v_j, j = 1, \dots, n - i_{m-1}$ having zero the first i_{m-1} components and equal to the ones of \tilde{v}_j the last $n - i_{m-1}$ components. It can be easily seen that v_j are solutions of $\dot{v} = Av, v(0) = e_{j+i_{m-1}}$, (e_j being the canonical basis of \mathbb{R}^n) and therefore they constitute the last $n - i_{m-1}$ columns of the matrix e^{tA} . Observe now that the last $n - i_{m-1}$ columns of the Jacobian $(\partial Y/\partial x)(x)$ are all zero since from the definition of homogeneity of degree one it follows that any component of the non linear part Y of X_0 is independent on the last $n - i_{m-1}$ variables. Therefore it can be easily verified that v_j are also solutions of 4.13 and hence they

constitute the last $n - i_{m-1}$ columns of $(\partial(\exp tX_0)/\partial p)(p)$ for any $p \in \mathbb{R}^n$. This shows that the last columns of the matrices e^{tA} , $(\partial(\exp tX_0)/\partial p)(p)$ are the same, for any $p \in \mathbb{R}^n$, which proves (4.12). \square

Remark 4.7. – Observe that the first i_{m-1} components of a map $\psi \in GH^{1,r}(\mathbb{R}^n)$ are independent on the last $n - i_{m-1}$ variables, while the last $n - i_{m-1}$ components depend linearly on the last $n - i_{m-1}$ variables. Therefore, if the first i_{m-1} components of a column vector B are zero and the last $n - i_{m-1}$ are constant, the same is true for the components of the vector $(\partial\psi/\partial y)B$. Thus it can be easily verified, using also Theorem 2.1 of Section 2, that any coordinate transformation $\phi \in GH^{1,r}(\mathbb{R}^n)$ transforms an affine system as (4.2) into a system of the same form. Hence we can always perform a coordinate change $x = \phi(y)$, $\phi \in GH^{1,r}(\mathbb{R}^n)$ that transforms the field X_0 in (4.2) into the canonical form $T_\phi X_0 = S + N$, $[S, N] = 0$ of Theorem 3.1 and use the formula given in Theorem 4.3 to derive the following explicit representation of the trajectories of (4.2)

$$x(t, u) = \phi \left(e^{tS}(\exp tN) \left(\phi^{-1}(x(0)) + \int_0^t e^{-s(S+\widehat{N})} B u(s) ds \right) \right),$$

$$t \in \mathbb{R}, \tag{4.14}$$

where \widehat{N} denotes the linear part of the field N .

Remark 4.8. – The flow generated by a homogeneous vector field of degree one can be in general calculated by solving a cascade system of the type

$$\dot{x}_i = A_i x_i + Y_i(x_1, \dots, x_{i-1}), \quad x_i \in \mathbb{R}^{d_i}, \quad d_i \in \mathbb{N} \setminus \{0\}.$$

This implies that systems of the form (4.2) can be always integrated without any restriction on the constant vector field B . However, formula (4.10) allows to derive a general description of the attainable set of a nonlinear system of type (4.2) as the image, through the flow generated by X_0 , of the attainable set of its linearized system.

Remark 4.9. – The representation of the trajectories of control systems of the form (4.2) given by formulas (4.10), (4.14) can be used to study problems of local controllability and to construct asymptotically stabilizing feedback controls for affine nonlinear control systems. Some interesting results in this sense have been recently obtained by H. Hermes in [14] where it is shown that non-resonance conditions (as stated in Corollary 3.9 here) and small time local controllability (STLC) imply many standard

necessary conditions for the existence of a continuous, asymptotically stabilizing feedback control (ASFC), for an n -dimensional, single input affine control system.

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