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# NORMAL FORMS WITH EXPONENTIALLY SMALL REMAINDER AND GEVREY NORMALIZATION FOR VECTOR FIELDS WITH A NILPOTENT LINEAR PART 

by Patrick BONCKAERT \& Freek VERSTRINGE


#### Abstract

We explore the convergence/divergence of the normal form for a singularity of a vector field on $\mathbb{C}^{n}$ with nilpotent linear part. We show that a Gevrey- $\alpha$ vector field $X$ with a nilpotent linear part can be reduced to a normal form of Gevrey- $1+\alpha$ type with the use of a Gevrey- $1+\alpha$ transformation. We also give a proof of the existence of an optimal order to stop the normal form procedure. If one stops the normal form procedure at this order, the remainder becomes exponentially small.

Résumé. - Nous investiguons la convergence/divergence de la forme normale d'une singularité d'un champ de vecteurs de $\mathbb{C}^{n}$ avec une partie linéaire nilpotente. Nous prouvons que chaque champ de vecteurs Gevrey- $\alpha$ avec une partie linéaire nilpotente peut être réduit à une forme normale Gevrey- $1+\alpha$ en utilisant une transformation Gevrey- $1+\alpha$. Nous prouvons également que si on arrête la procédure de normalisation à un certain ordre optimal, le reste de la forme normale devient exponentiellement petit.


## 1. Introduction and statement of the results

We consider vector fields where the linear part at a singularity is nilpotent, with no restriction on the dimension. This, for example, includes the case of a coupled Takens-Bogdanov system, see e.g., [8]. See [9] for an introduction to the subject.

We briefly give some history of the subject. In [12] the planar case $y \frac{\partial}{\partial x}+$ $\ldots$ was considered and a formal normal form $(y+a(x)) \frac{\partial}{\partial x}+b(x) \frac{\partial}{\partial y}$ was derived. It was shown in [10], also in the planar case, that an analytic vector

[^0]field with nilpotent linear part $y \frac{\partial}{\partial x}$ can be analytically transformed to a normal form. Other results related to the planar case are in [1].

More recently it was shown in [6] that the analytic vector fields with linear part $y \frac{\partial}{\partial x}+z \frac{\partial}{\partial y}$ can be Gevrey-1 reduced to a normal form using a specific normal form procedure that is also described later on in this article. This framework was extended in [6] to the case of quasihomogeneous vector fields. In [6] it is explained what the generalization of the so called small denominators are for non-diagonal linear vector fields (and more general quasihomogeneous vector fields); and some results of convergence and Gevrey-1 normalization are explained (See also theorems 2.3 and 2.4).

In [5] the optimal order to cut off the normal form procedure was determined. At this order the remainder becomes exponentially small. This was done for vector fields with a semi-simple linear part with eigenvalues that satisfy $\left|\langle\lambda, \delta\rangle-\lambda_{j}\right| \geqslant \gamma /|\delta|^{\tau}$ for each $\delta \in \mathbb{N}^{n}, j \in\{1, \ldots, n\}$ for which $\langle\lambda, \delta\rangle-\lambda_{j} \neq 0$ and a certain $\gamma, \tau>0$; resonant eigenvalues are allowed. A similar result was also obtained for vector fields with a completely degenerate linear part (i.e., all eigenvalues are zero) and at most one non-trivial Jordan block of order 2 or 3 . The difficulty starts with Jordan blocks of higher dimension or multiple Jordan blocks because it seems rather difficult to compute the generalized small denominators directly using the method at hand. In this article we will succeed in computing these generalized small denominators using the $\operatorname{sl}(2, \mathbb{C})$ structure that is hidden in the problem.

Somewhat before that, we have the results of [2] on the formal structure of the normal forms with a nilpotent linear part, using representation theory of $\mathrm{sl}(2, \mathbb{C})$. More recently [11] and [8] have also made contributions to this subject in the multidimensional case, on the formal level.

The purpose of this article is to combine both ideas : we will show how to use representation theory of $\operatorname{sl}(2, \mathbb{C})$ in nilpotent cases in order to calculate the small denominators in the framework of [6] and hence obtain qualitative information on the growth of coefficients appearing inside the normal form procedure.

Let us state the main results. We say that the linear part $N$ of a vector field $X$ is nilpotent at 0 if it acts as a nilpotent linear operator on the space of polynomials of degree $\delta$, for each $\delta \in \mathbb{N} \backslash\{0\}$. Note that this means, up to a linear change of the coordinates, that the linear part of the vector field can be written as $N=\sum_{i=1}^{n-1} a_{i} x_{i+1} \frac{\partial}{\partial x_{i}}$, for certain $a_{1}, \ldots, a_{n-1} \in \mathbb{C}$. We say that a vector field $X=\sum_{i=1}^{n} X_{i} \frac{\partial}{\partial x_{i}}$ is Gevrey- $\alpha$ if each $X_{i}$ is a formal power series $\sum_{l \geqslant 0} \sum_{|k|=l} a_{k} x^{k}$ and there exists a $C, r>0$ such that $\sum_{|k|=l}\left|a_{k}\right| \leqslant C r^{l}(l!)^{\alpha}$; the sum is taken over $k \in \mathbb{N}^{n}$, and $|k|=k_{1}+\ldots+k_{n}$.

Theorem 1.1. - Every formal Gevrey- $\alpha$ vector field $X=N+R$, where $N$ is a nilpotent linear part and $R$ is a part of higher order, admits a formal Gevrey- $1+\alpha$ transformation to Gevrey- $1+\alpha$ normal form. If the Gevrey- $\alpha$ vector field is formally linearizable, then the transformation and corresponding normal form are Gevrey- $\alpha$.

Theorem 1.2. - Suppose that $X$ is an analytic vector field in a neighbourhood of 0 of $\mathbb{C}^{m}$ that vanishes at the origin. This is

$$
X=L+\sum_{\delta \geqslant 2} f_{\delta},
$$

where $L$ is linear and nilpotent, $f_{\delta}$ is a bounded $\delta$-linear symmetric form and

$$
\left\|f_{\delta}\left(x_{1}, \ldots, x_{\delta}\right)\right\| \leqslant \frac{c\left\|x_{1}\right\| \ldots\left\|x_{\delta}\right\|}{\rho^{\delta}}
$$

with $c, \rho>0$ independent of $\delta$. Suppose that $Q$ is an invertible linear transformation with norm $\|Q\|=\sup _{\|x\|=1}\|Q(x)\|$ for which $N=Q L Q^{-1}$ is as in Proposition 4.3. Let $\nu=\sup \frac{e^{2} p!}{p^{p+1 / 2} e^{-p}}$,

$$
\begin{aligned}
C & =\frac{\sqrt{m}}{\rho^{2}}\left\{\left(\frac{5}{2} m+2\right) c\|Q\|^{2} \cdot\left\|Q^{-1}\right\|^{2}+3 \rho\|Q\| \cdot\left\|Q^{-1}\right\|\right\} \\
M_{0} & =\frac{10}{9}\left\{\left(\nu \sqrt{\frac{27}{8 e}}\right)+(2 e)^{2}\right\},
\end{aligned}
$$

$w=\frac{1}{e C}$ and $p_{\mathrm{opt}}=\left[\frac{1}{e C}\right]$. There exists a polynomial coordinate transformation id $+u_{p_{\text {opt }}}$ of degree at most $p_{\text {opt }}$ such that the vector field expressed in new coordinates becomes $Y=L+R_{p_{\mathrm{opt}}}+T_{p_{\mathrm{opt}}}$, where $R_{p_{\mathrm{opt}}}$ is a polynomial of degree $p_{\mathrm{opt}}$ and $T_{p_{\mathrm{opt}}}=O\left(|x|^{p_{\mathrm{opt}}+1}\right)$ is analytic and exponentially small i.e., for $0<\epsilon$ small enough we have the estimate

$$
\sup _{\|x\| \leqslant \varepsilon}\left\|T_{p_{\mathrm{opt}}}(x)\right\| \leqslant M_{0} \varepsilon^{2} e^{-w / \varepsilon}
$$

Moreover $R_{p_{\text {opt }}}\left(e^{t \bar{L}} x\right)=e^{t \bar{L}} R_{p_{\text {opt }}}(x)$, where $\bar{L}=Q^{-1} N^{*} Q$.
Remark that the cases where $N=y \frac{\partial}{\partial x}$ and $N=\left(y \frac{\partial}{\partial x}+z \frac{\partial}{\partial y}\right)$ have already been treated in [5]. This theorem provides a generalization and a geometric explanation using representation theory of $\operatorname{sl}(2, \mathbb{C})$ of these examples. Considering the results in [10] and [7] one could wonder whether or not the given normal form actually converges, when $X$ is analytic (i.e., $\alpha=0$ ). We think however that, in general, this is not the case.

In Section 2 of this article we repeat some results of the framework created in [6], in order to be self-contained. In Section 3 we state some results
on the representation theory of $\mathrm{sl}(2, \mathbb{C})$. In Section 4 we prove some propositions that lead to the main results, stated as Theorem 1.1 and Theorem 1.2 above and proven in Section 5.

## 2. Background and notation

We recall some standard preliminaries about the used normal form procedure. We follow the outline of [6], although similar ideas appear in [3] and [5].

### 2.1. Setting

Let $X=N+R$ be a local formal vector field in the neighbourhood of the origin, $N$ its linear part and $R$ the part of order $\geqslant 2$. We will look for a coordinate transform $\Phi^{-1}=I+U, U$ of order $\geqslant 2$, such that the pullback $\Phi_{*}(X)=X^{\prime}=N+R^{\prime}$. A minor calculation shows that

$$
\begin{align*}
& \Phi_{*}(X)=N+R^{\prime} \\
\Leftrightarrow & X \circ \Phi^{-1}=D \Phi^{-1} \cdot X^{\prime} \\
\Leftrightarrow & (S+R) \circ(I+U)=D(I+U) \cdot\left(N+R^{\prime}\right) \\
\Leftrightarrow & R^{\prime}+[U, N]=R(I+U)-D U \cdot R^{\prime} \tag{2.1}
\end{align*}
$$

Now we are going to determine the terms of order $\delta$ for the formal series $U=U_{2}+U_{3}+U_{4}+\ldots$ and $R^{\prime}=R_{2}^{\prime}+R_{3}^{\prime}+R_{4}^{\prime}+\ldots$ by induction. Therefore suppose that we already know $U_{2}, \ldots, U_{\delta-1}$ and $R_{2}^{\prime}, \ldots, R_{\delta-1}^{\prime}$. We take the projection of the terms of order $\delta$ in (2.1) and obtain:

$$
R_{\delta}^{\prime}+\left[U_{\delta}, N\right]=R H S_{\delta}\left(U_{<\delta}, R_{<\delta}^{\prime}\right)
$$

where $R H S_{\delta}$ stands for the projection of order $\delta$ of the right hand side of (2.1) and depends only on $U_{l}, R_{l}^{\prime}$ with index strictly smaller than $\delta$. Therefore it is natural to introduce the Lie-operator

$$
d_{0, \delta}: \mathcal{V}_{\delta} \longrightarrow \mathcal{V}_{\delta}: U \mapsto[U, N] ;
$$

where

$$
\begin{aligned}
\mathcal{V}_{\delta} & =\left\{\left.\sum_{i=1}^{n} P_{i} \frac{\partial}{\partial x_{i}} \right\rvert\, P_{i} \in \mathcal{P}_{\delta+1}\right\} \\
\mathcal{P}_{\delta} & =\{P \mid P \text { is a homogeneous polynomial of degree } \delta\}
\end{aligned}
$$

and decompose the space $\mathcal{V}_{\delta}$ of vector fields of degree $\delta$ as $\mathcal{V}_{\delta}=\operatorname{Im}\left(d_{0, \delta}\right) \oplus$ $\mathcal{W}_{\delta}$, where $\mathcal{W}_{\delta}$ is a particular choice of a complementary space that is induced by an inner product. This is explained in detail in the next section. Remark that we sometimes drop the $\delta$ in the notation whenever no confusion is possible.

### 2.2. The choice of the complementary subspaces $\mathcal{W}_{\delta}$

In order to define the complementary spaces $\mathcal{W}_{\delta}$ we need the adjoint action of $d_{0}$ with respect to an inner product. Therefore we have the following

Definition 2.1. - We define an inner product on $\mathcal{P}_{\delta}$, the space of polynomials of degree $\delta$ as

$$
\left\langle\sum_{|\alpha|=\delta} a_{\alpha} x^{\alpha}, \sum_{|\beta|=\delta} b_{\beta} x^{\beta}\right\rangle=\sum_{|\alpha|=\delta} a_{\alpha} \bar{b}_{\alpha} \frac{\alpha!}{|\alpha|!} .
$$

This induces an inner product on the space $\mathcal{V}_{\delta-1}$ of vector fields of degree $\delta-1$ as follows:

$$
\begin{equation*}
\left\langle\sum_{k=1}^{n} V_{k} \frac{\partial}{\partial x_{k}}, \sum_{k=1}^{n} W_{k} \frac{\partial}{\partial x_{k}}\right\rangle=\sum_{k=1}^{n}\left\langle V_{k}, W_{k}\right\rangle_{\delta} \tag{2.2}
\end{equation*}
$$

where the $V_{k}, W_{k}$ are elements of $\mathcal{P}_{\delta}$.
Now we define $d_{0}^{*}$ as the adjoint action of $d_{0}$ with respect to the above inner product. We repeat that $d_{0}^{*}$ is defined as the unique linear map satisfying $\left\langle d_{0}^{*}(V), W\right\rangle=\left\langle V, d_{0}(W)\right\rangle$, for all $V, W \in \mathcal{V}_{\delta}$. From linear algebra we know that:
(1) The operators $\square_{\delta}=d_{0} d_{0}^{*}$ are self-adjoint.
(2) The operators $\square_{\delta}$ are diagonizable.
(3) The operators $\square_{\delta}$ have real positive eigenvalues.
(4) $\mathcal{V}_{\delta}=\operatorname{Ker}\left(\square_{\delta}\right) \oplus \operatorname{Im}\left(\square_{\delta}\right)=\operatorname{Ker}\left(d_{0}^{*}\right) \oplus \operatorname{Im}\left(d_{0}\right)$

We will from now on choose the complementary subspace $W_{\delta}$ as $\operatorname{Ker}\left(d_{0}^{*}\right)$ $=\operatorname{Ker}\left(\square_{\delta}\right)$.

We recall from [6] a nice way to calculate the adjoint operator $d_{0}^{*}$. Let us first define the isomorphism:

$$
\phi: \mathcal{V}_{\delta} \longrightarrow \mathcal{P}_{\delta+1}^{n}: \sum_{k=1}^{n} V_{k} \frac{\partial}{\partial x_{k}} \mapsto\left(V_{1}, V_{2}, \ldots, V_{n}\right)
$$

Lemma 2.2 ([6], p. 691). - Suppose that $V=\sum_{k=1}^{n} V_{k} \frac{\partial}{\partial x_{k}} \in \mathcal{V}_{\delta}$. Then we have

$$
\phi\left(d_{0}^{*}(V)\right)=\left(\begin{array}{cccc}
N^{*}-\left(\frac{\partial N_{1}}{\partial x_{1}}\right)^{*} & -\left(\frac{\partial N_{2}}{\partial x_{1}}\right)^{*} & \ldots & -\left(\frac{\partial N_{n}}{\partial x_{1}}\right)^{*} \\
-\left(\frac{\partial N_{1}}{\partial x_{2}}\right)^{*} & N^{*}-\left(\frac{\partial N_{2}}{\partial x_{2}}\right)^{*} & \cdots & -\left(\frac{\partial N_{n}}{\partial x_{2}}\right)^{*} \\
\vdots & & & \vdots \\
-\left(\frac{\partial N_{1}}{\partial x_{n}}\right)^{*} & \cdots & -\left(\frac{\partial N_{n-1}}{\partial x_{n}}\right)^{*} & N^{*}-\left(\frac{\partial N_{n}}{\partial x_{n}}\right)^{*}
\end{array}\right)\left(\begin{array}{c}
V_{1} \\
\vdots \\
\vdots \\
V_{n}
\end{array}\right) .
$$

We give an example in the case that $n=2$ and $N=x_{2} \frac{\partial}{\partial x_{1}}$. In this case we have

$$
\phi\left(d_{0}^{*}\left(V_{1} \frac{\partial}{\partial x_{1}}+V_{2} \frac{\partial}{\partial x_{2}}\right)\right)=\left(\begin{array}{cc}
N^{*} & 0 \\
-1 & N^{*}
\end{array}\right)\binom{V_{1}}{V_{2}}
$$

Hence $d_{0}^{*}\left(V_{1} \frac{\partial}{\partial x_{1}}+V_{2} \frac{\partial}{\partial x_{2}}\right)=N^{*} V_{1} \frac{\partial}{\partial x_{1}}+\left(-V_{1}+N^{*} V_{2}\right) \frac{\partial}{\partial x_{2}}$.

### 2.3. Resonant terms and small denominators

When solving equation (2.1), we decompose $R H S_{\delta}=Q_{\delta} \oplus T_{\delta}$, where $Q_{\delta} \in \operatorname{Ker}\left(\square_{\delta}\right)$ and $T_{\delta} \in \operatorname{Im}\left(\square_{\delta}\right)=\operatorname{Im}\left(d_{0}\right)$. Now let $\Lambda$ be the set of eigenvalues with multiplicity of the operator $\square_{\delta}$ and $\Lambda^{*}$ be the set of nonzero eigenvalues. Since $\square_{\delta}$ is diagonizable, it is possible to decompose $T_{\delta}$ in a base of eigenvectors of $\square_{\delta}$. More precisely:

$$
T_{\delta}=\square_{\delta}\left(V_{\delta}\right)=\sum_{\lambda \in \Lambda^{*}} \square_{\delta}\left(V_{\delta, \lambda}\right)=\sum_{\lambda \in \Lambda^{*}} \lambda V_{\delta, \lambda} .
$$

If we define $W_{\delta, \lambda}=d_{0}^{*}\left(V_{\delta, \lambda}\right)$ and $W_{\delta}=\sum_{\lambda \in \Lambda^{*}} W_{\delta, \lambda}$, then

$$
d_{0}\left(W_{\delta}\right)=d_{0}\left(\sum_{\lambda \in \Lambda^{*}} W_{\delta, \lambda}\right)=\sum_{\lambda \in \Lambda^{*}} \square_{\delta}\left(V_{\delta, \lambda}\right)=\sum_{\lambda \in \Lambda^{*}} \lambda V_{\delta, \lambda}=T_{\delta} .
$$

Moreover since

$$
\left\langle W_{\delta, \lambda}, W_{\delta, \lambda}\right\rangle=\left\langle d_{0}^{*}\left(V_{\delta, \lambda}\right), d_{0}^{*}\left(V_{\delta, \lambda}\right)\right\rangle=\left\langle V_{\delta, \lambda}, \square_{\delta}\left(V_{\delta, \lambda}\right)\right\rangle=\lambda\left\langle V_{\delta, \lambda}, V_{\delta, \lambda}\right\rangle
$$

it follows that we have the estimates:

$$
\left\|T_{\delta}\right\|^{2}=\left\|d_{0}\left(W_{\delta}\right)\right\|^{2}=\sum_{\lambda \in \Lambda^{*}} \lambda^{2}\left\|V_{\delta, \lambda}\right\|^{2}=\sum_{\lambda \in \Lambda^{*}} \lambda\left\|W_{\delta, \lambda}\right\|^{2} \geqslant\left(\min _{\lambda \in \Lambda^{*}}(\sqrt{\lambda})\left\|W_{\delta}\right\|\right)^{2}
$$

This estimate makes clear that the $\lambda$ 's will play the role of the small denominators. We explain now what we mean by ' $S$ satisfies a diophantine condition'. Therefore, following [6], p. 675, we introduce the numbers $a_{\delta}=\min _{\lambda \in \Lambda^{*}}(\sqrt{\lambda})$ and define the numbers $\eta_{\delta}$, for $\delta \geqslant 0$, recursively by (let $\eta_{0}=1$ )

$$
a_{\delta} \eta_{\delta}=\max _{\delta_{1}+\ldots+\delta_{r}=\delta} \eta_{\delta_{1}} \ldots \eta_{\delta_{r}}
$$

where the maximum is taken over the set where at least two of the $\delta_{i}$ 's are strictly positive. We say that $S$ satisfies a diophantine condition if the $\eta_{\delta} \leqslant c M^{\delta}$ for certain positive constants $c, M$.

We will say that " $S$ satisfies a Siegel condition of order $(\tau, \gamma)$ " or more shortly "of order $\tau$ " whenever we have the estimates:

$$
a_{\delta} \geqslant \frac{\gamma}{\delta^{\tau}}
$$

for a certain $\gamma>0$ and $\tau \geqslant 0$. These conditions are important because of the following three theorems:

Theorem 2.3 ([6], Theorem 5.6 p. 676 and Remark 6.7 on p. 686). Suppose that $X=N+R$ is a formal Gevrey- $\alpha$ vector field that is formally linearizable to its linear part $N$. Suppose that $N$ satisfies a diophantine condition, then $X$ is Gevrey- $\alpha$ linearizable.

Theorem 2.4 ([6], Theorem 6.2 p. 683 and Remark 6.7 on p. 686). Suppose that $X=N+R$ is a Gevrey- $\alpha$ vector field that has a formal normal form $X^{\prime}=N+R^{\prime}$ by means of the procedure explained in this section, and suppose that the linear part $N$ of $X$ satisfies a Siegel condition of order $\tau$, then $X^{\prime}$ and $\Phi$ are formal power series of type Gevrey- $(1+\tau+\alpha)$.

Theorem 2.5 ([5], Corollary 1.9 p. 7 and [6], Theorem 6.11 p. 689). Suppose that $X$ is an analytic vector field in a neighbourhood of 0 of $\mathbb{C}^{m}$ that vanishes at the origin. This is

$$
X=L+\sum_{\delta \geqslant 2} f_{\delta}
$$

where $L$ is linear, $f_{\delta}$ is a bounded $\delta$-linear symmetric form and

$$
\left\|f_{\delta}\left(x_{1}, \ldots, x_{\delta}\right)\right\| \leqslant \frac{c\left\|x_{1}\right\| \ldots\left\|x_{\delta}\right\|}{\rho^{\delta}}
$$

with $c, \rho>0$ independent of $\delta$. Suppose that $Q$ is an invertible linear transformation with norm $\|Q\|=\sup _{\|x\|=1}\|Q(x)\|$ for which $N=Q L Q^{-1}$ satisfies the Siegel condition of order $(\tau, \gamma)$. Let $\nu=\sup \frac{e^{2} p!}{p^{p+1 / 2} e^{-p}}, a=1 / \gamma$,

$$
\begin{aligned}
C & =\frac{\sqrt{m}}{\rho^{2}}\left\{\left(\frac{5}{2} m+2\right) a c\|Q\|^{2} \cdot\left\|Q^{-1}\right\|^{2}+3 \rho\|Q\| \cdot\left\|Q^{-1}\right\|\right\} \\
M_{\tau} & =\frac{10}{9}\left\{\left(\nu \sqrt{\frac{27}{8 e}}\right)^{1+\tau}+(2 e)^{2+2 \tau}\right\}
\end{aligned}
$$

$b=\frac{1}{1+\tau}, w=\frac{1}{e C^{b}}$ and $p_{\mathrm{opt}}=\left[\frac{1}{e C^{b}}\right]$. There exists a polynomial coordinate transformation id $+u_{p_{\mathrm{opt}}}$ of degree at most $p_{\mathrm{opt}}$ such that the vector field
expressed in new coordinates becomes $Y=L+R_{p_{\text {opt }}}+T_{p_{\text {opt }}}$, where $R_{p_{\text {opt }}}$ is a polynomial of degree $p_{\mathrm{opt}}$ and $T_{p_{\mathrm{opt}}}=O\left(|x|^{p_{\mathrm{opt}}+1}\right)$ is analytic and exponentially small i.e., for $0<\epsilon$ small enough we have the estimate

$$
\sup _{\|x\| \leqslant \varepsilon}\left\|T_{p_{\mathrm{opt}}}(x)\right\| \leqslant M_{\tau} \varepsilon^{2} e^{-w / \varepsilon^{b}}
$$

Moreover $R_{p_{\text {opt }}}\left(e^{t \bar{L}} x\right)=e^{t \bar{L}} R_{p_{\text {opt }}}(x)$, where $\bar{L}=Q^{-1} N^{*} Q$.
Remark 2.6. - If $Q$ is a unitary transformation, then the above estimates hold with $\|Q\|,\left\|Q^{-1}\right\|$ replaced by 1 ; moreover then also $\bar{L}=N^{*}$ holds.

## 3. Representations of $\operatorname{sl}(2, \mathbb{C})$

We briefly recall the definition of a Lie algebra, a representation of a Lie algebra and some related algebraic concepts.

Definition 3.1. - A Lie algebra $(\mathfrak{g},[]$,$) is a vector space \mathfrak{g}$ provided with a multiplication [,] $: \mathfrak{g} \times \mathfrak{g} \mapsto \mathfrak{g}:(x, y) \mapsto[x, y]$ that satisfies the relations

- $\left[g_{1}, g_{2}\right]=-\left[g_{2}, g_{1}\right]$,
- $\left[g_{1},\left[g_{2}, g_{3}\right]\right]+\left[g_{2},\left[g_{3}, g_{1}\right]\right]+\left[g_{3},\left[g_{1}, g_{2}\right]\right]=0$.

We list the following concepts:
(1) A Lie algebra $\mathfrak{g}$ is called simple iff $[\mathfrak{g}, \mathfrak{g}]=\mathfrak{g}$.
(2) A Lie algebra homomorphism is a linear map $L: \mathfrak{g} \longrightarrow \mathfrak{h}$ between two Lie algebra's preserving the product structure : $L\left(\left[g_{1}, g_{2}\right]\right)=$ $\left[L\left(g_{1}\right), L\left(g_{2}\right)\right]$.
(3) A Lie algebra representation of $\mathfrak{g}$ is a Lie algebra homomorphism $L: \mathfrak{g} \longrightarrow g l(V)$, where $V$ is a vector space and $g l(V)$ is the group of linear transformations from $V$ to $V$.
(4) A Lie algebra representation $L: \mathfrak{g} \longrightarrow g l(V)$ is irreducible, iff there exist no subspace $W$ different from $V$ or $\{0\}$ such that $L(g)(w) \in$ $W$, for every $w \in W$ and every $g \in \mathfrak{g}$. A subspace $W$ with this property defines a subrepresentation.

We will need one of the key results of representations of simple Lie algebra's. A proof can be found e.g., in [4].

Theorem 3.2. - Every finite dimensional representation of a simple Lie algebra $\mathfrak{g}$ can be written as a direct sum of irreducible representations of $\mathfrak{g}$.

We now recall some basic facts of the representations of the simple Lie algebra $\operatorname{sl}(2, \mathbb{C})$. Let us first recall the definition.

Definition 3.3. - We define the Lie algebra $\mathrm{sl}(2, \mathbb{C})$ as the subalgebra of $g l_{2}(\mathbb{C})$ of matrices with trace 0 . It is generated by the matrices

$$
N=\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right), M=\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right), H=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right) .
$$

Remark 3.4. - Any Lie algebra generated by three elements $N, M, H$ and subject to the relations

$$
[H, N]=2 N,[H, M]=-2 M,[N, M]=H
$$

is isomorphic with $\mathrm{sl}(2, \mathbb{C})$. Moreover it is now clear that $\mathrm{sl}(2, \mathbb{C})$ is a simple Lie algebra.

The following theorem is well-known: a proof can be found e.g., in [4].
Theorem 3.5. - For every $n$ the representation of $\operatorname{sl}(2, \mathbb{C})$ defined by

$$
\begin{aligned}
& \left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right) \mapsto \widetilde{N}_{n}=\left(\begin{array}{cccccc}
0 & n & 0 & 0 & \ldots & 0 \\
0 & 0 & n-1 & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \ldots & 0 & 2 & 0 \\
0 & 0 & \ldots & 0 & 0 & 1 \\
0 & 0 & \ldots & 0 & 0 & 0
\end{array}\right) \\
& \left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right) \mapsto \widetilde{M}_{n}=\left(\begin{array}{cccccc}
0 & 0 & 0 & 0 & \ldots & 0 \\
1 & 0 & 0 & 0 & \ldots & 0 \\
0 & 2 & 0 & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \ldots & n-1 & 0 & 0 \\
0 & 0 & \ldots & 0 & n & 0
\end{array}\right) \\
& \left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right) \mapsto \widetilde{H}_{n}=\left(\begin{array}{cccccc}
n & 0 & 0 & 0 & \ldots & 0 \\
0 & n-2 & 0 & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \ldots & 0 & -(n-2) & 0 \\
0 & 0 & \ldots & 0 & 0 & -n
\end{array}\right)
\end{aligned}
$$

and acting on $\mathbb{C}^{n+1}$ is irreducible. Moreover, any other irreducible representation of $\operatorname{sl}(2, \mathbb{C})$ is isomorphic with one of these representations.

## 4. Construction of some particular $\operatorname{sl}(2, \mathbb{C})$ representations

In this section we focus on the construction of some particular $\operatorname{sl}(2, \mathbb{C})$ representations. In order to make the computations a bit more transparent, we use the correspondence between matrices and vector fields by a bijection $\phi: \sum_{i} \sum_{j} a_{i, j} x_{j} \frac{\partial}{\partial x_{i}} \mapsto A=\left(a_{i j}\right)$. Now suppose that we have two vector fields $A_{v}=\sum_{j=1}^{n} \sum_{i=1}^{n} a_{i j} x_{i} \frac{\partial}{\partial x_{j}}$ and $B_{v}=\sum_{j=1}^{n} \sum_{i=1}^{n} b_{i j} x_{i} \frac{\partial}{\partial x_{j}}$ with corresponding matrices $A$ and $B$, then the Lie bracket transforms as $\phi\left(\left[A_{v}, B_{v}\right]\right)=A B-B A$.

We start now with the construction. Therefore we define the following vector fields:

$$
\begin{aligned}
N_{n} & :=\alpha_{1} x_{2} \frac{\partial}{\partial x_{1}}+\alpha_{2} x_{3} \frac{\partial}{\partial x_{2}}+\ldots+\alpha_{n} x_{n+1} \frac{\partial}{\partial x_{n}} \\
M_{n} & =\alpha_{1} x_{1} \frac{\partial}{\partial x_{2}}+\alpha_{2} x_{2} \frac{\partial}{\partial x_{3}}+\ldots+\alpha_{n} x_{n} \frac{\partial}{\partial x_{n+1}} \\
H_{n} & :=\left[N_{n}, M_{n}\right] .
\end{aligned}
$$

It is important to note that $M_{n}$ is the adjoint of $N_{n}$ with respect to the inner product (2.2). We will use the same notation $N_{n}, M_{n}$ and $H_{n}$ for the associated matrices and drop the index $n$ where no confusion is possible. We want to choose the coefficients $\alpha_{1}, \ldots, \alpha_{n}$ in such a way that they are nonzero and that the triple $N, M, H$ is isomorphic to the Lie algebra $\operatorname{sl}(2, \mathbb{C})$. Therefore it is sufficient to ensure that the relations described in Remark 3.4 are satisfied. The third relation is automatic from the construction. We focus on the first relation. In matrix notation this relation becomes $H N-$ $N H-2 N=0$. Now remark that

$$
\begin{gathered}
N=\left(\begin{array}{ccccc}
0 & \alpha_{1} & 0 & \ldots & 0 \\
0 & 0 & \alpha_{2} & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \ldots & 0 & \alpha_{n} \\
0 & 0 & \ldots & 0 & 0
\end{array}\right), M=\left(\begin{array}{ccccc}
0 & 0 & 0 & \ldots & 0 \\
\alpha_{1} & 0 & 0 & \ldots & 0 \\
0 & \alpha_{2} & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \ldots & \alpha_{n} & 0
\end{array}\right), \\
H=\left(\begin{array}{ccccc}
\alpha_{1}^{2} & 0 & 0 & \ldots & 0 \\
0 & \alpha_{2}^{2}-\alpha_{1}^{2} & 0 & \ldots & 0 \\
0 & 0 & \alpha_{3}^{2}-\alpha_{2}^{2} & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \ldots & \alpha_{n}^{2}-\alpha_{n-1}^{2} & 0 \\
0 & 0 & \ldots & 0 & -\alpha_{n}^{2}
\end{array}\right) .
\end{gathered}
$$

Hence this relation becomes

$$
\left(\begin{array}{cccccc}
0 & b_{1} & 0 & 0 & \ldots & 0 \\
0 & 0 & b_{2} & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \ldots & 0 & 0 & b_{n} \\
0 & 0 & \ldots & 0 & 0 & 0
\end{array}\right)=0
$$

where $b_{i}=\alpha_{i}\left(-\alpha_{i-1}^{2}+2 \alpha_{i}^{2}-\alpha_{i+1}^{2}\right)-2 \alpha_{i}$ for $2 \leqslant i \leqslant n-1, b_{1}=\alpha_{1}\left(2 \alpha_{1}^{2}-\right.$ $\left.\alpha_{2}^{2}\right)-2 \alpha_{1}$ and $b_{n}=\alpha_{n}\left(-\alpha_{n-1}^{2}+2 \alpha_{n}^{2}\right)-2 \alpha_{n}$; and we need to solve the equations

$$
\left\{\begin{array}{l}
\alpha_{1}\left(2 \alpha_{1}^{2}-\alpha_{2}^{2}\right)=2 \alpha_{1} \\
\alpha_{2}\left(-\alpha_{1}^{2}+2 \alpha_{2}^{2}-\alpha_{3}^{2}\right)=2 \alpha_{2} \\
\quad \vdots \\
\alpha_{n}\left(-\alpha_{n-1}^{2}+2 \alpha_{n}^{2}\right)=2 \alpha_{n}
\end{array}\right.
$$

Since we suppose that none of the $\alpha_{i}$ vanish, this simplifies to

$$
\left(\begin{array}{cccccc}
2 & -1 & 0 & 0 & \ldots & 0 \\
-1 & 2 & -1 & 0 & \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & \ldots & 0 & -1 & 2 & -1 \\
0 & \ldots & 0 & 0 & -1 & 2
\end{array}\right)\left(\begin{array}{c}
\alpha_{1}^{2} \\
\alpha_{2}^{2} \\
\vdots \\
\alpha_{n-1}^{2} \\
\alpha_{n}^{2}
\end{array}\right)=\left(\begin{array}{c}
2 \\
2 \\
\vdots \\
2
\end{array}\right)\left(\begin{array}{c}
2 \\
2 \\
\vdots \\
2
\end{array}\right)
$$

One can verify that a solution is given by $\alpha_{i}^{2}=i(n+1-i)$ for $1 \leqslant i \leqslant n$. We choose the positive solutions and put $\alpha_{i}=\sqrt{i(n+1-i)}$. Then it is readily verified (repeat the above calculations) that also the second relation $[H, M]=-2 M$ from Remark 3.4 is satisfied. We have now proven the

Lemma 4.1. - Let $n \in \mathbb{N}$ and define $N_{n}=\sum_{i=1}^{n} \sqrt{i(n+1-i)} x_{i+1} \frac{\partial}{\partial x_{i}}$, then the triple $N_{n}, M_{n}:=N_{n}^{*}, H=\left[N_{n}, M_{n}\right]$ defines a Lie-algebra isomorphic to $\mathrm{sl}(2, \mathbb{C})$.

We are now ready to show that
Lemma 4.2. - Let $\delta \in \mathbb{N} \backslash\{0\}$. For a given $N_{n}=\sum_{i=1}^{n} \sqrt{i(n+1-i)}$ $x_{i+1} \frac{\partial}{\partial x_{i}}$ the associated triple $d_{0}, d_{0}^{*}, D=\left[d_{0}, d_{0}^{*}\right]$ defines an $\operatorname{sl}(2, \mathbb{C})$ representation. Here is $d_{0}$ the Lie operator $U \mapsto\left[N_{n}, U\right]$ acting on $\mathcal{V}_{\delta}$ and $d_{0}^{*}$ its adjoint.

Proof. - Put $\alpha_{i}=\sqrt{i(n+1-i)}$ and $I$ the identity operator. The Lie operator acting on vector fields $d_{0}\left(\sum_{i=1}^{n+1} V_{i} \frac{\partial}{\partial x_{i}}\right)$ can be expressed using
matrix notation as

$$
\left(\begin{array}{cccccc}
N & -\alpha_{1} I & 0 & 0 \ldots & 0 & \\
0 & N & -\alpha_{2} I & 0 \ldots & 0 & \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & \cdots & 0 & 0 & N & -\alpha_{n} I \\
0 & \cdots & 0 & 0 & 0 & N
\end{array}\right)\left(\begin{array}{c}
V_{1} \\
V_{2} \\
\vdots \\
V_{n} \\
V_{n+1}
\end{array}\right)
$$

and its adjoint $d_{0}^{*}\left(\sum_{i=1}^{n+1} V_{i} \frac{\partial}{\partial x_{i}}\right)$ as (see also Lemma 2.2 and remember that $\left.M=N^{*}\right)$

$$
\left(\begin{array}{cccccc}
M & 0 & 0 & 0 \ldots & 0 & \\
-\alpha_{1} I & M & 0 & 0 \ldots & 0 & \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & \ldots & 0 & -\alpha_{n-1} I & M & 0 \\
0 & \ldots & 0 & 0 & -\alpha_{n} I & M
\end{array}\right)\left(\begin{array}{c}
V_{1} \\
V_{2} \\
\vdots \\
V_{n} \\
V_{n+1}
\end{array}\right) .
$$

Hence it is readily verified that the commutator $D=\left[d_{0}, d_{0}^{*}\right]=d_{0} d_{0}^{*}-d_{0}^{*} d_{0}$ can be expressed as

$$
\left(\begin{array}{ccccc}
H+\alpha_{1}^{2} I & 0 & 0 & \cdots & 0 \\
0 & H+\left(\alpha_{2}^{2}-\alpha_{1}^{2}\right) I & 0 & \cdots & 0 \\
0 & 0 & H+\left(\alpha_{3}^{2}-\alpha_{2}^{2}\right) I & \cdots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \cdots & H+\left(\alpha_{n}^{2}-\alpha_{n-1}^{2}\right) I & 0 \\
0 & 0 & \cdots & 0 & H-\alpha_{n}^{2} I
\end{array}\right),
$$

where $H=[N, M]$. Now the commutator $\left[D, d_{0}\right]$ simplifies as

$$
\left(\begin{array}{cccccc}
a_{1} & b_{1} & 0 & 0 \ldots & 0 & \\
0 & a_{2} & b_{2} & 0 & 0 \ldots & 0 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & 0 & 0 \ldots & a_{n} & b_{n} \\
0 & 0 & 0 & 0 \ldots & 0 & a_{n+1}
\end{array}\right)
$$

with

$$
\begin{aligned}
a_{i} & =H N-N H+\left(\alpha_{i}^{2}-\alpha_{i-1}^{2}\right)(N I-I N)=[H, N]=2 N, \\
b_{i} & =-\alpha_{i}\left(H+\alpha_{i}^{2} I-\alpha_{i-1}^{2} I\right)+\alpha_{i}\left(H+\alpha_{i+1}^{2} I-\alpha_{i}^{2} I\right)= \\
& =-\alpha_{i}\left(-\alpha_{i+1}^{2}+2 \alpha_{i}^{2}-\alpha_{i-1}^{2}\right) I \\
& =-2 \alpha_{i} I
\end{aligned}
$$

where we have put $\alpha_{0}=0$ and $\alpha_{n+1}=0$ in the above calculation. We also used the fact that the triple $N, M, H$ defines a Lie algebra isomorphic to
$\operatorname{sl}(2, \mathbb{C})$. Hence $\left[D, d_{0}\right]=2 d_{0}$.
Making analogous calculations, one verifies that also $\left[D, d_{0}^{*}\right]=-2 d_{0}^{*}$.
As a corollary of this lemma we can consider the case of multiple nilpotent blocks as follows. Remark that we allow zero blocks (i.e., $k_{i}=0$ for some $i$ ).

Proposition 4.3. - Let $k_{1}, k_{2}, \ldots, k_{n}$ be natural numbers and let $x^{i}$ be a $\left(k_{i}+1\right)$-dimensional variable $\left(x_{1}^{i}, \ldots, x_{k_{i}+1}^{i}\right)$, for $1 \leqslant i \leqslant n$. Let $N=N_{k_{1}}\left(x^{1}\right)+\ldots+N_{k_{n}}\left(x^{n}\right)$, where

$$
N_{k_{j}}\left(x^{j}\right)=\sum_{i=1}^{k_{j}} \sqrt{i(n-i+1)} x_{i+1}^{j} \frac{\partial}{\partial x_{i}^{j}}, \quad N_{0}=0 .
$$

Then the triple $N, M:=N^{*}, H=[N, M]$ defines a Lie algebra isomorphic to $\mathrm{sl}(2, \mathbb{C})$. Moreover let $d_{0}$ be the associated Lie operator, then also the triple $d_{0}, d_{0}^{*}, D=\left[d_{0}, d_{0}^{*}\right]$ defines a Lie algebra isomorphic to $\operatorname{sl}(2, \mathbb{C})$.

Proof. - Use the concept of a direct sum, lemma 4.1 and Lemma 4.2.

## 5. Proof of the main theorems

This is rather a summary of all the foregoing. From linear algebra we know that, up to a linear change of variables, it is no restriction to start with a vector field $X=N+R$ where $N$ is as in Proposition 4.3. Let now $d_{0}$ be the associated Lie operator. Let $\delta \in \mathbb{N} \backslash\{0\}$. We are interested in the calculation of eigenvalues of the associated operator $\square_{\delta}=d_{0} d_{0}^{*}$ acting on $\mathcal{V}_{\delta}$. According to Proposition 4.3, we know that the triple $d_{0}, d_{0}^{*}$ and $D=\left[d_{0}, d_{0}^{*}\right]$ defines a Lie algebra isomorphic to $\mathrm{sl}(2, \mathbb{C})$. It follows, using Theorem 3.2, that the associated representation acting on $\mathcal{V}_{\delta}$ can be split in a direct sum of irreducible representations. Hence, up to a linear coordinate transform $\varphi$ (acting on the space $\mathcal{V}_{\delta}$ ), we can suppose that we are dealing with a representation of the form

$$
\begin{aligned}
& N=\widetilde{N}_{1} \oplus \widetilde{N}_{2} \oplus \ldots \oplus \widetilde{N}_{l} \\
& M=\widetilde{M}_{1} \oplus \widetilde{M}_{2} \oplus \ldots \oplus \widetilde{N}_{l}, \\
& H=\widetilde{H}_{1} \oplus \widetilde{H}_{2} \oplus \ldots \oplus \widetilde{H}_{l}
\end{aligned}
$$

where $\widetilde{N}_{i}, \widetilde{M}_{i}$ and $\widetilde{H}_{i}$ are as in Theorem 3.5. Hence $\varphi$ transforms the operator $\square_{\delta}=d_{0} d_{0}^{*}$ into $N M$. The nonzero eigenvalues of the operator $N M=\widetilde{N}_{1} \widetilde{M}_{1} \oplus \ldots \oplus \widetilde{N}_{l} \widetilde{M}_{l}$ are positive integers because each $\widetilde{N}_{i} \widetilde{M}_{i}$ is a diagonal matrix containing integers on the diagonal. Hence the same is true for the operator $\square_{\delta}$. Now using Theorem 2.4 with $\tau=0$ and $\gamma=1$
(or Theorem 2.3 in case the vector field is formally linearizable) finishes the proof Theorem 1.1. Similarly, using Theorem 2.5 with $\tau=0$ and $\gamma=1$ finishes the proof Theorem 1.2.

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