# Linear Ordinary Differential Equations and $\mathcal{D}$-modules, solving and reduction methods. 

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#### Abstract

Systems of linear ODEs form a monoidal category, and objects within this category include homomorphisms, symmetric and exterior powers, as well as dual equations.

Considering symmetries of equation we find equations that can be solved by algebraic manipulations and quadratures. This includes a number of equations with semisimple symmetry algebras, as well as equations with solvable symmetry algebras. Unlike most previous results, there is no requirement on the symmetry algebras of having dimension equal to the order of the equations, in some cases even a single symmetry is sufficient to solve the equation.


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- Til Bodøjentan: gleda mæ til neste jentetur, spautgaven...


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## Chapter 1

## Introduction

This thesis is devoted to the study of linear ordinary differential equations. Main results are found in Chapters 5 and 6 , where we obtain methods to decompose and solve equations with both solvable and semisimple Lie algebras of symmetries. We prove that for a number of such equations one can obtain solutions through combining algebraic methods and quadrature. Also, there is no restriction on the order of an equation with respect to the dimension of its symmetry algebra, for this to work. Given the right conditions it may even be sufficient with a single symmetry to solve an equation through eigenvalue decomposition and quadrature.

Our starting point is to connect to systems of linear ODEs algebraic objects which we will denote differential modules, or $\mathcal{D}$-modules.

The notion of a differential module appears in differential algebra, see e.g. [15, 16], but differential Galois theory and Picard-Vessiot theory deals with modules over differential fields, and mainly the study of differential field extensions by solutions of ODEs. This approach may be used to state whether solutions are algebraic with respect to the base field, study the extensions and address inverse problems in differential Galois theory etc.

Our approach has geometrical roots, dating back to Sophus Lie, and points in a different direction with respect to applications. Throughout this thesis we will consider systems of linear ODEs with coefficients in the differential $\mathbb{R}$-algebra $\mathcal{A}=C^{\infty}(\mathbb{R})$, with derivation being the usual derivative in the variable $x, \delta_{\mathcal{A}}=\frac{d}{d x}$.

There is a one-to-one correspondence

$$
\text { System of linear ODEs } \quad \Leftrightarrow \quad \mathcal{D} \text {-module }(E, \delta) \text { over }\left(\mathcal{A}, \delta_{\mathcal{A}}\right)
$$

by the isomorphism of vector spaces

$$
\text { Solution space of the ODE } \Leftrightarrow \operatorname{ker} \delta \subset E \text {. }
$$

A straightforward explanation of this correspondence can be found in Section 2.1 , based only on the pure algebraic definition of $\mathcal{D}$-modules.

Sophus Lie initiated a geometric approach to differential equations, where one uses symmetries of equations to study their properties and to reduce and solve them. Viewing ODEs as submanifolds of an appropriate jet space provides a geometrical framework which is widely used to study geometric properties and symmetries of equations. From this framework $\mathcal{D}$-modules emerge in the following way.

A linear ODE is a linear sub-bundle in jet space, with a linear connection determined by the Cartan distribution. The $\mathcal{D}$-module corresponding to an equation can be identified with the $\mathcal{A}$-module of sections in the linear bundle, with derivation determined by the lifting of $\frac{d}{d x}$ by the linear connection. This relation is accounted for in Section 2.4.

Chapter 2 is devoted to establish the relation between equations and $\mathcal{D}$ modules, and appropriate algebraic tools. In Section 2.1 we learn precisely how to connect the two objects, and in combination with Section 2.3.
$\mathcal{D}$-modules form a monoidal category, and Section 2.2 describes this category and algebraic constructions within the category. That forms the basic algebraic framework used to produce the results of this thesis, with the key result being Theorem 2.2.3. It is our main tool, and allows us to lift properties and results that apply to the vector space $\operatorname{ker} \delta$ to the whole module $E$. Finally, in Section 2.5, we introduce a third view on $\mathcal{D}$-modules, through differential operators. From a practical viewpoint this is an important addition to the theory, providing us with a generic way to calculate within our $\mathcal{D}$-modules.

In Chapter 3 we investigate equations with Euclidean, symplectic, complex and Hermitian structures. For second order equations we determine classes of equations with such structures. Also, we encounter $\mathcal{D}$-modules with $S_{2}$-representations due to solutions of the Yang - Baxter equation.

Chapter 4 contains results on symmetries of equations in general. Again due to Theorem 2.2.3 we establish how to apply results from the theory of representations of Lie algebras into vector spaces to $\mathcal{D}$-modules, and whence to equations with these symmetry algebras. Chapters 5 and 6 are based on this observation.

In Section 4.2 we further see how to use symmetry operators when working with $\mathcal{D}$-modules. Proposition 4.2.2 determines how a symmetry operator of an equation induces a $\delta$-invariant endomorphism of the corresponding $\mathcal{D}$ module, and Theorem 4.2.2 explains how it acts inside ker $\delta$.

In Chapter 5 we investigate equations with solvable symmetry algebras and eigenvalue decompositions of $\mathcal{D}$-modules. We find conditions for when an equation can be solved by using a single symmetry, Theorem 5.1.1. We determine equations with solvable symmetry algebras that are solvable in terms of quadratures, see Theorem 5.2.5.

In Chapter 6 we encounter semisimple symmetry algebras. For a semisimple Lie algebra $\mathfrak{g}$ there is an associated symmetry ring $\mathcal{D}(\mathfrak{g})$, an analouge of the the Grothendieck ring of isomorphism classes of finite dimensional vector space representations of $\mathfrak{g}$, and its symmetry ring is generated by a finite number of elements just as its Grothendieck ring is. The generators are isomorphism classes of $\mathcal{D}$-modules with symmetry algebra $\mathfrak{g}$, Theorem 6.1.2.

As a consequence, any $\mathcal{D}$-module with symmetry algebra $\mathfrak{g}$ is polynomial in $\mathcal{D}$-modules isomorphic to the generators, meaning that solutions of the generator differential equations generate all solutions of the original equation.

In particular, any equation with an $\mathfrak{s l}_{2}$-algebra of symmetries has solution space spanned by powers of solutions of second order model equations of $\mathfrak{s l}_{2}$, Schrödinger equations, Theorem 6.2.4. Solutions may in many cases be obtained by algebraic methods and quadrature, and an algorithmic approach is outlined.

## Chapter 2

## Equations as $\mathcal{D}$-modules

The goal of this chapter is to make a connection between differential modules, denoted $\mathcal{D}$-modules, and linear ODEs. These $\mathcal{D}$-modules are free modules of finite rank over the $\mathbb{R}$-algebra of smooth functions in one real variable, with a derivation $\delta$ over $\frac{d}{d x}$. They form a monoidal category. We will establish a range of results concerning $\mathcal{D}$-modules, that can be directly applied to the corresponding equations.

### 2.1 Connecting modules and equations

### 2.1.1 $\mathcal{D}$-modules over a general algebra.

Fix an algebra $A$ over a field $K$, and a derivation $\delta_{A}: A \rightarrow A$. A pair $\left(\mathcal{A}, \delta_{\mathcal{A}}\right)$ is called a differential algebra.

Definition 2.1.1. A $\mathcal{D}$-module over $\left(A, \delta_{A}\right)$ is a pair $(E, \delta)$ where $E$ is a module over $A$ and the map

$$
\delta: E \rightarrow E
$$

is a derivation over $\delta_{A}$, i. e. it is (i) $K$-linear, and satisfies a Leibniz rule (ii) with respect to $\delta_{A}$ :

$$
\begin{array}{llr}
(i) & \delta\left(e_{1}+e_{2}\right)=\delta\left(e_{1}\right)+\delta\left(e_{2}\right) & e_{1}, e_{2} \in E \\
\text { (ii) } & \delta(a e)=\delta_{A}(a) e+a \delta(e) & a \in A, e \in E \tag{2.1}
\end{array}
$$

Throughout the text we will restrict ourselves to consider free $\mathcal{D}$-modules, i.e. free modules over an algebra $A$, that are also $\mathcal{D}$-modules. $\mathcal{D}$-modules
$(E, \delta)$ over a fixed pair $\left(A, \delta_{A}\right)$ constitute the objects of a category which we will denote $\mathcal{C}$, and morphisms are $A$-homomorphisms of modules that commute with the respective derivations.

Proposition 2.1.1. $\mathcal{C}$ is monoidal with respect to the tensor product of modules with the induced derivation $\delta$ over $\delta_{A}$ as defined in Definition 2.2.1.

Note that $\left(A, \delta_{A}\right)$ is a unit object in $\mathcal{C}$.
Section 2.2 is devoted to the case when $A=C^{\infty}(\mathbb{R})$ and $\delta_{A}=\frac{d}{d x}$, and contains a number of constructions and results for that category. More or less all definitions and results in that section apply to $\mathcal{D}$-modules over a gerneral $\left(A, \delta_{A}\right)$, with the exception of Theorem 2.2.3, and some particular statements involving properties of $C^{\infty}(\mathbb{R})$.

### 2.1.2 $\mathcal{D}$-modules and linear ODEs.

Now let $\mathcal{A}=C^{\infty}(\mathbb{R})$, the $\mathbb{R}$-algebra of smooth functions in one real variable. The pair $\left(\mathcal{A}, \delta_{A}=\frac{d}{d x}\right)$ is a differential ring, $\frac{d}{d x}$ being a derivation of $\mathcal{A}$.

Definition 2.1.2. By a $\mathcal{D}$-module over $\left(C^{\infty}(\mathbb{R}), \frac{d}{d x}\right)$ we mean a pair $(E, \delta)$ where $E$ is a free module of rank $=n<\infty$ over $C^{\infty}(\mathbb{R})$ and the map

$$
\delta: E \rightarrow E
$$

is a derivation over $\delta_{A}=\frac{d}{d x}$, i. e. it is (i) $\mathbb{R}$-linear, and satisfies a Leibniz rule (ii) with respect to $\frac{d}{d x}$ :
(i) $\delta\left(e_{1}+e_{2}\right)=\delta\left(e_{1}\right)+\delta\left(e_{2}\right), \quad e_{1}, e_{2} \in E$,
(ii) $\quad \delta(a e)=\frac{d a}{d x} e+a \delta(e), \quad a \in A, e \in E$.

From the definition we can immediately deduce a correspondence between a $\mathcal{D}$-module $(E, \delta)$ of rank $n$ and a system oflinear ordinary differential equations.

Theorem 2.1.1. Given a rank $=n \mathcal{D}$-module $(E, \delta)$ as in Definition 2.1.2. Then the $\mathbb{R}$-vector space $E^{\#}=\operatorname{ker} \delta \subset E$ is isomorphic to the solution space of an $n \times n$ system of linear first order differential equations.

Proof. $E$ is a free module of rank $n$ over $\mathcal{A}$, so there is a basis $\left\{e_{1}, \ldots, e_{n}\right\}$ of $E$ over $\mathcal{A}$. The action of $\delta$ on $E$ can be written in matrix form

$$
\delta(\underline{e})=A \underline{e}
$$

where $A=\left(a_{i j}(x)\right)$, entries $a_{i j} \in \mathcal{A}$, and $\underline{e}=\left[e_{1}, \ldots, e_{n}\right]^{T}$. Considering a general element

$$
\begin{equation*}
h=h_{1}(x) e_{1}+\ldots+h_{n}(x) e_{n} \in E, \tag{2.3}
\end{equation*}
$$

coefficients $h_{i}(x) \in \mathcal{A}$. Then $\delta$ applied to $h$ is

$$
\begin{align*}
\delta(h) & =\sum_{i=1}^{n}\left(h_{i}^{\prime} e_{i}+h_{i} \delta e_{i}\right)=\sum_{i=1}^{n}\left(h_{i}^{\prime} e_{i}+h_{i} \sum_{j=1}^{n} a_{i j} e_{j}\right)  \tag{2.4}\\
& =\sum_{s=1}^{n}\left(h_{s}^{\prime}+\sum_{i=1}^{n} a_{i s} h_{i}\right) e_{s} .
\end{align*}
$$

Thus,

$$
\delta(h)=0
$$

if and only if the coefficient functions $h_{1}(x), \ldots, h_{n}(x)$ satisfy the system

$$
\begin{equation*}
\underline{h}^{\prime}+A^{T} \underline{h}=0, \tag{2.5}
\end{equation*}
$$

where $\underline{h}=\left[h_{1}(x), \ldots, h_{n}(x)\right]^{T}$.
The map

$$
\phi: \underline{h} \text { solution of }(2.5) \mapsto h=\sum_{i=1}^{n} h_{i} e_{i} \in \operatorname{ker} \delta
$$

is an isomorphism of vector spaces.

### 2.2 The monoidal category of linear ODEs

This section is dedicated to investigate the algebraic framework in which $\mathcal{D}$ modules belong, and we will establish some useful properties and constructions. Throughout we may keep in mind the relation between $\mathcal{D}$-modules and systems of ODEs as described in Theorem 2.1.1, and recall that every time we consider some $\mathcal{D}$-module it corresponds to a particular equation.

Proposition 2.2.1. $\mathcal{D}$-modules in the sense of Definition 2.1.2, over the fixed differential algebra

$$
\left(\mathcal{A}=C^{\infty}(\mathbb{R}), \delta_{A}=\frac{d}{d x}\right)
$$

constitute the objects of a category which we will denote $\mathcal{L O D E}$, linear ODEs.
For objects $\left(E_{1}, \delta_{1}\right),\left(E_{2}, \delta_{2}\right)$ in $\mathcal{L O D E}$, morphisms, $\operatorname{Mor}\left(\left(E_{1}, \delta_{1}\right),\left(E_{2}, \delta_{2}\right)\right)$, are $\mathcal{A}$-homomorphisms $F$ such that the diagram (2.6) commutes

$$
\begin{array}{ccc}
E_{1} \xrightarrow{F} & E_{2} \\
\delta_{1} \downarrow & &  \tag{2.6}\\
E_{1} \xrightarrow[F]{ } & & { }^{2} \\
\delta_{2}
\end{array}
$$

i.e. $\delta_{2} \circ F=F \circ \delta_{1}$.

Proof. The only category property of composition of $\mathcal{D}$-module morphisms we need to check is that the composition $\psi \circ \phi$ of $\mathcal{A}$-homomorphisms
$\phi: E_{1} \rightarrow E_{2}$ and $\psi: E_{2} \rightarrow E_{3}$ really satisfies the necessary commutator relations. But

$$
\begin{align*}
(\psi \circ \phi) \circ \delta_{1} & =\psi \circ\left(\phi \circ \delta_{1}\right)=\psi \circ\left(\delta_{2} \circ \phi\right) \\
& =\left(\psi \circ \delta_{2}\right) \circ \phi=\delta_{3} \circ(\psi \circ \phi), \tag{2.7}
\end{align*}
$$

so $\phi \in \operatorname{Mor}\left(\left(E_{1}, \delta_{1}\right),\left(E_{2}, \delta_{2}\right)\right), \psi \in \operatorname{Mor}\left(\left(E_{2}, \delta_{2}\right),\left(E_{3}, \delta_{3}\right)\right)$ implies that $\psi \circ \phi \in \operatorname{Mor}\left(\left(E_{1}, \delta_{1}\right),\left(E_{3}, \delta_{3}\right)\right)$.

All tensorial constructions of $\mathcal{D}$-modules belong in this category, and each correponds to an equation. We may take the tensor product of two modules in the category, and define a corresponding $\mathcal{D}$-module with an induced $\delta$ as follows.

Definition 2.2.1. The product of two $\mathcal{D}$-modules $\left(E_{1}, \delta_{1}\right)$ and $\left(E_{2}, \delta_{2}\right)$ is the object $\left(E_{1} \otimes_{A} E_{2}, \delta\right)$ where

$$
\delta: E_{1} \otimes_{A} E_{2} \longrightarrow E_{1} \otimes_{A} E_{2}
$$

is defined by the requirement that it is a derivation over $\left(A, \delta_{A}\right)$ and that

$$
\delta\left(e_{1} \otimes e_{2}\right)=\delta_{1}\left(e_{1}\right) \otimes e_{2}+e_{1} \otimes \delta_{2}\left(e_{2}\right)
$$

on decomposable elements, $e_{1} \in E_{1}, e_{2} \in E_{2}$.

Theorem 2.2.1. $\mathcal{L O D E}$ is monoidal with respect to the tensor product of modules with the induced $\delta$ as defined in Definition 2.2.1.
$\left(\mathcal{A}, \delta_{\mathcal{A}}\right)$ is a unit object in the category.
Proof. We may take basis element $1 \in \mathcal{A}$. For any object $(E, \delta)$ in the category we have the following isomorphism $l: \mathcal{A} \otimes E \rightarrow E$ defined by

$$
1 \otimes e \mapsto e, e \in E
$$

and requiring $\mathcal{A}$-linearity. Note that $l$ is a morphism in the category:

$$
\begin{align*}
\left(\delta_{E} \circ l\right)(b \otimes e) & =\delta_{E}(e)=l\left(b \otimes \delta_{E}(e)\right) \\
& =l \circ\left(\delta_{\mathcal{A}} b \otimes e+b \otimes \delta_{E} e\right)  \tag{2.8}\\
& =l \circ \delta_{\mathcal{A} \otimes E}(b \otimes e)
\end{align*}
$$

Corollary 2.2.1. The unit object $\left(\mathcal{A}, \delta_{\mathcal{A}}\right)$ in the category $\mathcal{L O D E}$ corresponds to the first order equation

$$
\begin{equation*}
y^{\prime}=0 \tag{2.9}
\end{equation*}
$$

Proof. Obviously $f \in \operatorname{ker} \delta \subset \mathcal{A}$ if and only if $f^{\prime}=0$.
We may also investigate products of morphisms.
Given $\mathcal{D}$-modules $\left(E_{1}, \delta_{E_{1}}\right),\left(E_{2}, \delta_{E_{2}}\right),\left(F_{1}, \delta_{F_{1}}\right)$ and $\left(F_{2}, \delta_{F_{2}}\right)$ and two module homomorphisms $\phi: E_{1} \rightarrow F_{1}, \psi: E_{2} \rightarrow F_{2}$ we may consider their usual tensor product

$$
\begin{equation*}
\phi \otimes \psi: E_{1} \otimes E_{2} \rightarrow F_{1} \otimes F_{2}, \tag{2.10}
\end{equation*}
$$

which on decomposable elements is

$$
\phi \otimes \psi:\left(e_{1} \otimes e_{2}\right) \mapsto \phi\left(e_{1}\right) \otimes \psi\left(e_{2}\right) .
$$

Proposition 2.2.2. Given morphisms $\phi \in \operatorname{Mor}\left(\left(E_{1}, \delta_{E_{1}}\right),\left(F_{1}, \delta_{F_{1}}\right)\right)$ and $\psi \in \operatorname{Mor}\left(\left(E_{2}, \delta_{E_{2}}\right),\left(F_{2}, \delta_{F_{2}}\right)\right)$. Then their product is again a morphism, $\phi \otimes \psi \in \operatorname{Mor}\left(\left(E, \delta_{E}\right),\left(F, \delta_{F}\right)\right)$, where $E=E_{1} \otimes_{\mathcal{A}} E_{2}, F=F_{1} \otimes_{\mathcal{A}} F_{2}$ and $\delta_{E}, \delta_{F}$ the induced derivations on the products $E, F$.

Proof. We need only check that $\phi \otimes \psi$ satisfies the necessary composition property $(\phi \otimes \psi) \circ \delta_{E}=\delta_{F} \circ(\phi \otimes \psi)$. Writing $\delta_{E}=\delta_{E_{1}} \otimes I_{E_{2}}+I_{E_{1}} \otimes \delta_{E_{2}}$ and $\delta_{F}=\delta_{F_{1}} \otimes I_{F_{2}}+I_{F_{1}} \otimes \delta_{F_{2}}$ in product notation we see that

$$
\begin{align*}
(\phi \otimes \psi) \circ \delta_{E} & =(\phi \otimes \psi) \circ\left(\delta_{E_{1}} \otimes I_{E_{2}}+I_{E_{1}} \otimes \delta_{E_{2}}\right) \\
& =\left(\phi \circ \delta_{E_{1}}\right) \otimes \psi+\phi \otimes\left(\psi \circ \delta_{E_{2}}\right) \\
& =\left(\delta_{F_{1}} \circ \phi\right) \otimes \psi+\phi \otimes\left(\delta_{F_{2}} \circ \psi\right)  \tag{2.11}\\
& =\left(\delta_{F_{1}} \otimes I_{F_{2}}+I_{F_{1}} \otimes \delta_{F_{2}}\right) \circ(\phi \otimes \psi) \\
& =\delta_{F} \circ(\phi \otimes \psi) .
\end{align*}
$$

Thus, $(\phi \otimes \psi) \in \operatorname{Mor}\left(\left(E, \delta_{E}\right),\left(F, \delta_{F}\right)\right)$.
Some $\mathcal{D}$-modules come with a bit of extra structure, we will encounter both algebras and Lie algebras, the $\mathcal{D}$-module versions are as follows.

Definition 2.2.2. A $\mathcal{D}$-algebra $(E, \delta)$ is a $\mathcal{D}$-module with a product

$$
m: E \otimes_{\mathcal{A}} E \longrightarrow E
$$

such that

$$
\begin{array}{ccc}
E \otimes E \xrightarrow{m} & E \\
\delta \downarrow & & \downarrow \delta  \tag{2.12}\\
E \otimes E \xrightarrow[m]{ } & E
\end{array}
$$

commutes, which satisfies the associativity condition

on $E$.

Note: If $(E, \delta)$ is a $\mathcal{D}$-algebra, then $\operatorname{ker} \delta \subset E$ is an $\mathbb{R}$-algebra.
Definition 2.2.3. A $\mathcal{D}$-Lie-algebra $(E, \delta)$ is a $\mathcal{D}$-module with a bracket

$$
[\cdot, \cdot]: E \times E \longrightarrow E
$$

which is
(1) $\mathcal{A}$-linear in both arguments, skew-symmetric and satisfies the Jacobi identity, i.e.

$$
[X,[Y, Z]]+[Z,[X, Y]]+[Y,[Z, X]]=0
$$

and,
(2) the bracket operation is $\delta$-invariant, i.e.

$$
\delta[X, Y]=[\delta X, Y]+[X, \delta Y] .
$$

Note: If $(E, \delta)$ is a $\mathcal{D}$-Lie algebra, then the solution space $\operatorname{ker} \delta$ is really a Lie algebra over $\mathbb{R}$ in the usual sense, with respect to the restriction of the bracket to ker $\delta$.

A natural construction to consider in the category $\mathcal{L O D E}$ is homomorphisms of modules.

Proposition 2.2.3. Given $\mathcal{D}$-modules $\left(E_{1}, \delta_{1}\right)$ and $\left(E_{2}, \delta_{2}\right)$, then $\left(\operatorname{Hom}_{A}\left(E_{1}, E_{2}\right), \delta\right)$ with

$$
\delta: \operatorname{Hom}_{A}\left(E_{1}, E_{2}\right) \longrightarrow \operatorname{Hom}_{A}\left(E_{1}, E_{2}\right)
$$

defined by

$$
\delta: F \mapsto \delta_{2} \circ F-F \circ \delta_{1}
$$

for $F \in \operatorname{Hom}_{A}\left(E_{1}, E_{2}\right)$, is an object in the category $\mathcal{L O D E}$.
Proof. Given $a \in A, F \in \operatorname{Hom}_{A}\left(E_{1}, E_{2}\right)$ and $e \in E_{1}$ we see that

$$
\begin{align*}
\delta F(a e) & =\left(\delta_{2} \circ F-F \circ \delta_{1}\right)(a \cdot e) \\
& =\delta_{2}(a \cdot F(e))-F\left(\delta_{A}(a) \cdot e+a \cdot \delta_{1}(e)\right)  \tag{2.14}\\
& =\delta_{A}(a) \cdot F(e)+a \delta_{2}(F(e))-\delta_{A}(a) \cdot F(e)-a \cdot F\left(\delta_{1}(e)\right) \\
& =a \cdot\left(\delta_{2} \circ F-F \circ \delta_{1}\right)(e)=a \cdot \delta F(e) .
\end{align*}
$$

So, $\delta F$ is an $A$-homomorphism. Second, we have that

$$
\begin{align*}
\delta(a F) & =\delta_{2} \circ(a \cdot F)-a \cdot F \circ \delta_{1} \\
& =\delta_{A} a \cdot F+a \cdot \delta_{2} \circ F-a \cdot F \circ \delta_{1}  \tag{2.15}\\
& =\delta_{A} a \cdot F+a \cdot \delta F .
\end{align*}
$$

Thus $\left(\operatorname{Hom}_{A}\left(E_{1}, E_{2}\right), \delta\right)$ is an object in $\mathcal{L O D E}$.

We pause to introduce the notation

$$
E^{\#}
$$

for taking the kernel of $\delta$ in $E$, i. e. $E_{1}^{\#}=\operatorname{ker} \delta_{1} \subset E_{1}$, and $E_{2}^{\#}=\operatorname{ker} \delta_{2} \subset$ $E_{2}$. Note that a $\delta$-invariant homomorphism $F \in\left(\operatorname{Hom}_{A}\left(E_{1}, E_{2}\right)\right)^{\#}$ is a morphism in the category, with

$$
\begin{equation*}
\delta_{2} \circ F=F \circ \delta_{1} . \tag{2.16}
\end{equation*}
$$

How can we interpret morphisms at this stage? If $\left(E_{1}, \delta_{1}\right)$ and $\left(E_{2}, \delta_{2}\right)$ correspond to ODE systems (2.17) and (2.18) respectively,

$$
\begin{align*}
& \underline{h}^{\prime}+A_{1} \underline{h}=0  \tag{2.17}\\
& \underline{u}^{\prime}+A_{2} \underline{u}=0 \tag{2.18}
\end{align*}
$$

then the vector spaces $E_{1}^{\#} \subset E_{1}$ and $E_{2}^{\#} \subset E_{2}$ are isomorphic to the solution spaces of the systems (2.17) and (2.18) respectively. By requirement

$$
\begin{equation*}
F\left(\operatorname{ker} \delta_{1}\right) \subset \operatorname{ker} \delta_{2} \tag{2.19}
\end{equation*}
$$

so a morphism $F \in\left(\operatorname{Hom}_{A}\left(E_{1}, E_{2}\right)\right)^{\#}$ maps solutions of system (2.17) into solutions of system (2.18). Thus solutions of the induced homomorphism equation

$$
\begin{equation*}
\left(\operatorname{Hom}_{\mathcal{A}}\left(E_{1}, E_{2}\right), \delta\right) \tag{2.20}
\end{equation*}
$$

give precisely linear maps that transfer solutions of equation $E_{1}$ to solutions of $E_{2}$. If system (2.17) is an $n \times n$ system and (2.18) is an $m \times m$ system, then the homomorphism equation is an $(n m) \times(n m)$ system.

Definition 2.2.4. For a $\mathcal{D}$-module $(E, \delta)$ in $\mathcal{L O D E}$, the $O D E$ corresponding to the induced $\mathcal{D}$-module

$$
\begin{equation*}
\left(E n d_{\mathcal{A}}(E), \delta\right) \tag{2.21}
\end{equation*}
$$

we will call the Lie equation, or symmetry equation of $(E, \delta)$.
The whole Chapter 4 is devoted to the study of $\delta$-invariant endomorphisms, symmetries, but we will briefly add the following result about endomorphisms here.

Proposition 2.2.4. $\left(E n d_{\mathcal{A}}(E), \delta\right)$ is
(i) an associative $\mathcal{D}$-algebra with respect to composition of endomorphisms, and,
(ii) a $\mathcal{D}$-Lie algebra with respect to commutators of endomorphisms.

Proof. In ( $i$ ) multiplication of endomorphisms $\phi, \psi$ is defined by

$$
\begin{equation*}
m(\phi \otimes \psi)=\phi \circ \psi \tag{2.22}
\end{equation*}
$$

Thus

$$
\begin{align*}
m(\delta(\phi \otimes \psi)) & =m(\delta \phi \otimes \psi+\phi \otimes \delta \psi) \\
& =\delta \phi \circ \psi+\phi \circ \delta \psi \\
& =(\delta \circ \phi-\phi \circ \delta) \circ \psi+\phi \circ(\delta \circ \psi-\psi \circ \delta)  \tag{2.23}\\
& =\delta \circ \phi \circ \psi-\phi \circ \psi \circ \delta=\delta(m(\phi \otimes \psi)),
\end{align*}
$$

and whence $\delta \circ m=m \circ \delta$. In (ii) the bracket operation is defined by

$$
\begin{equation*}
[\phi, \psi]=m(\phi \otimes \psi-\psi \otimes \phi) \tag{2.24}
\end{equation*}
$$

so it follows from the previous computation that the bracket operation and $\delta$ satisfy

$$
\delta[\phi, \psi]=[\delta \phi, \psi]+[\phi, \delta \psi]
$$

Proposition 2.2.5. For any $\mathcal{D}$-module $(E, \delta)$ in the category $\mathcal{L O D E}$, the dual module

$$
\left(E^{*}=\operatorname{Hom}_{A}(E, A), \delta\right)
$$

is also an object in $\mathcal{L O D E}$. The equation $\left(E^{*}, \delta\right)$ is the adjoint of equation $(E, \delta)$.

Proof. Assume $(E, \delta)$ is of rank $n$ with a basis $\left\{e_{1}, \ldots, e_{n}\right\}$ over $\mathcal{A}$, and that $\delta \underline{e}=A \underline{e}$ on matrix form. Taking the dual basis $\left\{e_{1}^{*}, \ldots, e_{n}^{*}\right\}$ as basis in $E^{*}$ and applying the definition of the induced $\delta$ on basis elements $e_{i}^{*} \in E^{*}$ yields

$$
\begin{align*}
\left(\delta e_{i}^{*}\right)\left(e_{j}\right) & =\left(\delta \circ e_{i}^{*}\right)\left(e_{j}\right)-\left(e_{i}^{*} \circ \delta\right)\left(e_{j}\right) \\
& =\delta\left(\delta_{i}^{j}\right)-e_{i}^{*}\left(\sum_{s=1}^{n} a_{j s} e_{s}\right)=0-\sum_{s=1}^{n} a_{j s} \delta_{i}^{s}=-a_{j i} \tag{2.25}
\end{align*}
$$

where $\delta_{i}^{j}$ denotes the Kronecker delta. Whence

$$
\begin{equation*}
\delta \underline{e}^{*}=-A^{T} \underline{e}^{*} \tag{2.26}
\end{equation*}
$$

on matrix form and the corresponding system of ODEs is

$$
\begin{equation*}
\underline{h}^{\prime}-A \underline{h}=0 \tag{2.27}
\end{equation*}
$$

which is the adjoint system.
Proposition 2.2.6. Given a $\delta$-invariant representation of a group $G$ into a $\mathcal{D}$-module $(E, \delta)$, i.e. a group representation

$$
r: G \rightarrow E n d_{\mathcal{A}}(E)
$$

such that $\delta \circ r(g)=r(g) \circ \delta$ for any $g \in G$. Then the set of $G$-invariant elements in $E$,

$$
\Sigma_{G}(E)=\{e \in E \mid r(g)(e)=e, \forall g \in G\}
$$

is a sub-D-module of $E$ with respect to restriction of $\delta$, i.e. $\left(\Sigma_{G}(E), \delta\right) \in$ $\operatorname{Ob}(\mathcal{L O D E})$.

Proof. $\Sigma_{G}(E)$ is obviously a sub-module of $E$, and we need only prove that $\delta\left(\Sigma_{G}(E)\right) \subset \Sigma_{G}(E)$. Let $e \in \Sigma_{G}(E)$. Then

$$
\delta(e)=\delta(r(g)(e))=r(g)(\delta(e))
$$

for any $g \in G$, thus $\delta e \in \Sigma_{G}(E)$.
There are some well-known constructions that are of this type. Both symmetric and anti-symmetric tensors, $S^{n}(E)$ and $\bigwedge^{n}(E)$ arise as invariant submodules of $E^{\otimes n}$ with respect to representations of the symmetric group $G=S_{n}$. Regarding representations of $S_{n}$ we have the following general result. Let $E$ be a finite rank module over $\mathcal{A}$. A homomorphism

$$
\begin{equation*}
\tau: E^{\otimes 2} \rightarrow E^{\otimes 2} \tag{2.28}
\end{equation*}
$$

with the condition $\tau^{2}=1$ determines a homomorphism

$$
\begin{equation*}
\tau_{i}=I_{i-1} \otimes \tau \otimes I_{n-i-1}: E^{\otimes n} \rightarrow E^{\otimes n} \tag{2.29}
\end{equation*}
$$

for $n \leq 2$, with $\tau$ acting only on the $i$ th and $(i+1)$ th copy of $E$ in $E^{\otimes n}$.

Theorem 2.2.2. $\tau_{1}, \ldots, \tau_{n-1}$ generate a representation of $S_{n}$ into $E^{\otimes n}$ iff $\tau$ satisfies the so called Yang-Baxter equation :

$$
\begin{equation*}
(\tau \otimes 1)(1 \otimes \tau)(\tau \otimes 1)=(1 \otimes \tau)(\tau \otimes 1)(1 \otimes \tau) \tag{2.30}
\end{equation*}
$$

on $E^{\otimes 3}$.
The symmetric power $S^{n}(E)=\Sigma_{\tau}^{n}\left(E^{\otimes n}\right)$ consists of elements in $E^{\otimes n}$ invariant with respect to the action of $S_{n}$ given by the twist-solution of the Yang-Baxter equation,

$$
\tau(f \otimes g)=g \otimes f
$$

Note that for the twist $\tau$

$$
\begin{aligned}
\delta(\tau(f \otimes g)) & =\delta(g \otimes f)=\delta g \otimes f+g \otimes \delta f \\
& =\tau(\delta f \otimes g+f \otimes \delta g)=\tau(\delta(f \otimes g))
\end{aligned}
$$

so this $\tau$ generates a $\delta$-invariant representation of $S_{2}$ into $E^{\otimes n}$.
Similarly, $\bigwedge^{n}(E)=\Sigma_{n}^{\tau}\left(E^{\otimes n}\right) \subset E^{\otimes n}$ where $\tau$ is minus twist,

$$
\tau(f \otimes g)=-g \otimes f
$$

Proposition 2.2.7. The symmetrization of the $k$ th tensor product of $E$, $S^{k}(E) \subset E^{\otimes k}$, is again an object in $\mathcal{L O D E}$ together with the restriction of $\delta$ on $E^{\otimes k}$ to

$$
\delta: S^{k}(E) \longrightarrow S^{k}(E)
$$

The restriction of $\delta$ acts on decomposable elements of $S^{k}(E)$ by

$$
\begin{equation*}
\delta\left(\theta_{1} \cdot \ldots \cdot \theta_{k}\right)=\delta\left(\theta_{1}\right) \cdot \ldots \cdot \theta_{k}+\ldots+\theta_{1} \cdot \ldots \cdot \delta\left(\theta_{k}\right), \tag{2.31}
\end{equation*}
$$

where $\theta_{i} \in E$, and $\cdot$ is the symmetric product.
Proposition 2.2.8. Any exterior power $\wedge^{k}(E) \subset E^{\otimes k}$ of a $\mathcal{D}$-module $(E, \delta)$ is an object in $\mathcal{L O D E}$ with $\delta$ being the restriction of $\delta$ on $E^{\otimes k}$ to

$$
\delta: \wedge^{k}(E) \longrightarrow \wedge^{k}(E)
$$

The restriction of $\delta$ acts on decomposable $k$-forms by

$$
\begin{equation*}
\delta\left(\omega_{1} \wedge \ldots \wedge \omega_{k}\right)=\delta\left(\omega_{1}\right) \wedge \ldots \wedge \omega_{k}+\ldots+\omega_{1} \wedge \ldots \wedge \delta\left(\omega_{k}\right) \tag{2.32}
\end{equation*}
$$

where $\omega_{i} \in E$.

In Chapter 3 we will see examples of non-trivial representations of $S_{2}$ for second order equations.

The theorem below will end this section, and will be our most important tool in the study of equations. It allows us to move between studying the $\mathcal{D}$-module itself, and the solution space of the corresponding differential equation. It will enable us to extend properties and structures on the $\mathbb{R}$-vector space $E^{\#}$ to the $\mathcal{D}$-module $E$, and apply to $\mathcal{D}$-modules. As an example, results from the theory of representations of Lie algebras into vector spaces.

Theorem 2.2.3. For any $(E, \delta) \in O b(\mathcal{L O D E})$

$$
\begin{equation*}
E \cong E^{\#} \otimes_{\mathbf{R}} \mathcal{A} \tag{2.33}
\end{equation*}
$$

by an $\mathcal{D}$-module isomorphism

$$
\begin{equation*}
\phi: E^{\#} \otimes_{\mathbf{R}} \mathcal{A} \rightarrow E \tag{2.34}
\end{equation*}
$$

defined by $\phi: v_{i} \otimes 1 \mapsto v_{i}$ for any basis $\left\{v_{1}, \ldots, v_{n}\right\}$ of $E^{\#}$.
Proof. First note that we may rephrase the statement, it is equivalent to the following:
For $(E, \delta) \in O b(\mathcal{L O D E})$ any basis of $E^{\#}$ over $\mathbb{R}$ is a basis of $E$ over $\mathcal{A}$.
Let $E$ be as above, rank $n$, with $\delta$-matrix $A$. Every element of $E^{\#}$ is on the form $h=\sum_{i=1}^{n} h_{i}(x) e_{i}$ where $\underline{h}$ solves (2.5). From the theory of ODEs we know that there exist a fundamental set of solutions of the system (2.5). Let $\underline{h}_{1}=\left[h_{11}(x), \ldots, h_{1 n}(x)\right]^{T}, \ldots, \underline{h}_{n}=\left[h_{n 1}(x), \ldots, h_{n n}(x)\right]^{T}$ be such a set. Then

$$
\left\{\gamma_{1}=\sum_{i} h_{1 i}(x) e_{i}, \ldots, \gamma_{n}=\sum_{i} h_{n i}(x) e_{i}\right\}
$$

is a basis of $E^{\#}$ over $\mathbb{R}$. The matrix $H=\left(h_{i j}(x)\right)$ is the Wronskian of the system (2.5), hence its determinant is non-zero everywhere, and

$$
\underline{\gamma}=H \underline{e}
$$

constitutes a basis of $E$ over $\mathcal{A}$. Any basis of $E^{\#}$ over $\mathbb{R}$ is on the form as the set $\left\{\gamma_{i}\right\}$ above, hence a basis of $E$ over $\mathcal{A}$.

Corollary 2.2.2. Given $\left(E_{1}, \delta_{1}\right),\left(E_{2}, \delta_{2}\right)$ in $\operatorname{Ob}(\mathcal{L O D E})$, then,
(i) $\left(E_{1} \otimes_{\mathcal{A}} E_{2}\right)^{\#}=E_{1}^{\#} \otimes_{\mathbb{R}} E_{2}^{\#}$, and,
(ii) $\left(\operatorname{Hom}_{\mathcal{A}}\left(E_{1}, E_{2}\right)\right)^{\#}=\operatorname{Hom}_{\mathbb{R}}\left(E_{1}^{\#}, E_{2}^{\#}\right)$.

Proof. This follows directly by combining the theorem above with the definitions of the induced $\delta$-s on $E_{1} \otimes E_{2}$ and $\operatorname{Hom}_{\mathcal{A}}\left(E_{1}, E_{2}\right)$ respectively.

### 2.3 Primitive element bases

Given an equation we wish to study, we need to know how to identify and work with the corresponding $\mathcal{D}$-module. Considering a system resolved into single equation, a convenient way to describe the corresponding module is to introduce the notion of a primitive element in $E$. We will use this approach throughout the rest of the thesis.

Definition 2.3.1. Let $(E, \delta) \in O b(\mathcal{L O D E})$, with $\operatorname{rank}_{\mathcal{A}}(E)=n$. An element $e \in E$ with the property that

$$
\begin{equation*}
\mathcal{B}=\left\{e_{1}=e, e_{2}=\delta e, e_{3}=\delta^{2} e, \ldots, e_{n}=\delta^{n-1} e\right\} \tag{2.35}
\end{equation*}
$$

is a basis of $E$ over $\mathcal{A}$ is called a primitive element of $E$, and $\mathcal{B}$ a primitive element basis of $E$.

In a primitive element basis as $\mathcal{B}$ above the action of $\delta$ is completely described by $n$ functions $a_{i}(x) \in \mathcal{A}$ where

$$
\begin{equation*}
\delta^{n} e=\sum_{i=1}^{n} a_{i}(x) \delta^{i-1} e \tag{2.36}
\end{equation*}
$$

In this basis the matrix form of the action of $\delta$ becomes

$$
\begin{equation*}
\delta \underline{e}=A \underline{e}, \tag{2.37}
\end{equation*}
$$

where

$$
\delta\left[\begin{array}{c}
e \\
\delta e \\
\cdot \\
\cdot \\
\cdot \\
\delta^{n-1} e
\end{array}\right]=\left[\begin{array}{cccccc}
0 & 1 & 0 & 0 & . & 0 \\
0 & 0 & 1 & 0 & . & 0 \\
\cdot & \cdot & \cdot & \cdot & . & \cdot \\
. & \cdot & \cdot & \cdot & 1 & 0 \\
0 & \cdot & \cdot & \cdot & 0 & 1 \\
a_{1} & a_{2} & \cdot & . & a_{n-1} & a_{n}
\end{array}\right]\left[\begin{array}{c}
e \\
\delta e \\
\cdot \\
\cdot \\
\cdot \\
\delta^{n-1} e
\end{array}\right]
$$

The advantage of this approach is that the corresponding equation system

$$
\begin{equation*}
\underline{h}^{\prime}+A^{t} \underline{h}=0 \tag{2.38}
\end{equation*}
$$

with $\underline{h}=\left[h_{1}(x), \ldots, h_{n}(x)\right]^{T}$, may then be resolved into a single equation

$$
\begin{equation*}
y^{(n)}+\left(a_{n} y\right)^{(n-1)}-\left(a_{n-1} y\right)^{(n-2)}+\ldots+(-1)^{n+1} a_{1} y=0 . \tag{2.39}
\end{equation*}
$$

That is, $h_{n}=y(x)$ and $h_{k}=H_{k}(y)$ for $k=1, \ldots, n-1$ where

$$
\begin{equation*}
H_{k}(y)=(-1)^{n-k}\left(y^{(n-k)}+\sum_{l=1}^{n-k}(-1)^{l+1}\left(a_{n+1-l} y\right)^{(n-k-l)}\right) . \tag{2.40}
\end{equation*}
$$

Written as an operator, with $\partial=\frac{d}{d x}, H_{k}$ is

$$
\begin{equation*}
H_{k}=(-1)^{n-k}\left(\partial^{(n-k)}+\sum_{l=1}^{n-k}(-1)^{l+1} \partial^{(n-k-l)} a_{n+1-l}\right) . \tag{2.41}
\end{equation*}
$$

We may sum up as follows.
Proposition 2.3.1. Given an equation on the form (2.39), then the corresponding $\mathcal{D}$-module $(E, \delta) \in O b(\mathcal{L O D E})$ has a primitive element basis

$$
\left\{e, \delta e, \delta^{2} e, \ldots, \delta^{n-1} e\right\}
$$

with $\delta: E \rightarrow E$ determined by the coefficients $a_{i}$ in the equation through the relation

$$
\delta^{n} e=\sum_{i=1}^{n} a_{i}(x) \delta^{i-1} e .
$$

The kernel of $\delta$ is

$$
\begin{equation*}
E^{\#}=\left\{h_{y}=\sum_{i=1}^{n} H_{k}(y) \delta^{i-1} e \mid y \text { solves (2.39) }\right\} \tag{2.42}
\end{equation*}
$$

with the operators $H_{k}$ on the form (2.41).
Remark: A natural concern is whether, starting with an equation on the form

$$
\begin{equation*}
y^{(n)}+f_{n}(x) y^{(n-1)}+f_{n-1}(x) y^{(n-2)}+\ldots+f_{1}(x) y=0, \tag{2.43}
\end{equation*}
$$

it is a problem to write it on the form (2.39), in order to be able to write down the structure of the corresponding $\mathcal{D}$-module $(E, \delta)$.

This is not a problem. To express the coefficients $a_{i}(x)$ in terms of the $f_{i}$-s and their derivatives we need only start with the highest coefficient $a_{n}=f_{n}$, and nest our way down to $a_{1}$. At each stage $a_{i}$ is given in terms of derivatives of the functions $f_{i}, f_{i+1}, \ldots, f_{n}$, and the formulas are recovered by performing the derivations in the expression (2.39) and collect terms of the same degree of derivatives of $y$ and compare with the form (2.43). The equations are on the form

$$
\begin{aligned}
f_{n} & =a_{n} \\
f_{n-1} & =a_{n}^{\prime}-a_{n-1} \\
f_{n-2} & =\binom{n-1}{1} a_{n}^{\prime \prime}-a_{n-1}^{\prime}+a_{n-2} \\
\vdots & \\
f_{1}= & =a_{n}^{(n-2)}-a_{n-1}^{(n-3)}+\ldots+(-1)^{n-1} a_{1}
\end{aligned}
$$

Knowing $E^{\#}$, i.e. knowing solutions of the corresponding equation, means that we can produce solutions of equations corresponding to such $\mathcal{D}$-modules as

$$
E^{\otimes n}, S^{k}(E), \wedge^{l}(E)
$$

since $\left(E^{\otimes n}\right)^{\#}=\left(E^{\#}\right)^{\otimes n}, S^{k}(E)^{\#}=S^{k}\left(E^{\#}\right)$ and $\bigwedge^{l}(E)^{\#}=\bigwedge^{l}\left(E^{\#}\right)$ are completely described when $E^{\#}$ is described. We will use this repeatedly.

We may use primitive element basises to precisely describe solutions of symmetric powers of second order equations. We will use this later, particularly in Section 6.2.
Let $(E, \delta)$ correspond to an equation

$$
\begin{equation*}
y^{\prime \prime}+\left(a_{2} y\right)^{\prime}-a_{1} y=0 \tag{2.44}
\end{equation*}
$$

with primitive element basis $\left\{e_{1}, e_{2}\right\}$. Let $\left\{e_{1}^{*}, e_{2}^{*}\right\}$ denote the dual basis of $\left(E^{*}, \delta\right)$. This basis generates a basis of the module $S^{k}\left(E^{*}\right)$

$$
\begin{equation*}
\left\{\alpha_{l}=\left(e_{1}^{*}\right)^{k-l+1} \cdot\left(e_{2}^{*}\right)^{l-1}, l=1, \ldots, k+1\right\} \tag{2.45}
\end{equation*}
$$

To find the equation $S^{k}\left(E^{*}\right)$ for some $k$ simply apply $\delta$ to a general element

$$
\begin{equation*}
\theta=\sum_{l=1}^{k+1} g_{l}(x) \alpha_{l} \tag{2.46}
\end{equation*}
$$

in $S^{k}\left(E^{*}\right)$. Recall that $\delta e_{1}^{*}=-a_{1} e_{1}^{*}$, and $\delta e_{2}^{*}=-e_{1}^{*}-a_{2} e_{1}^{*}$. Thus

$$
\begin{equation*}
\delta: \alpha_{l} \mapsto-a_{1}(k-l+1) \alpha_{l+1}-a_{2}(l-1) \alpha_{l}-(l-1) \alpha_{l-1} . \tag{2.47}
\end{equation*}
$$

Setting

$$
\delta \theta=0
$$

and collecting basis terms $\alpha_{s}$ yields a system of $k+1$ equations,

$$
\begin{equation*}
g_{s}^{\prime}-a_{1}(k-s+2) g_{s-1}-a_{2}(s-1) g_{s}-s g_{s+1}= \tag{2.48}
\end{equation*}
$$

for $s=1, \ldots, k+1$. This system resolves into a single equation in $g_{1}=y(x)$. We may conclude the following about $S^{k}\left(E^{*}\right)^{\#}$.

Proposition 2.3.2. Let $(E, \delta)$ be the $\mathcal{D}$-module corresponding to an arbitrary second order equation (2.44). For each $k \geq 1$ the kernel $S^{k}\left(E^{*}\right)^{\#}$ consists of elements

$$
\begin{equation*}
\theta_{y}=y \alpha_{1}+y^{\prime} \alpha_{2}+\sum_{l=3}^{k+1} g_{l}(y) \alpha_{l} \tag{2.49}
\end{equation*}
$$

where

$$
\begin{equation*}
g_{l}=\frac{1}{l-1}\left[g_{l-1}^{\prime}-a_{2}(l-2) g_{l-1}-a_{1}(k-l+3) g_{l-2}\right] \tag{2.50}
\end{equation*}
$$

for $l=2, . ., k+1$ and $g_{1}=y$ solves the $S^{k}\left(E^{*}\right)$ equation, i.e. the equation in $y$ we obtain from setting

$$
\delta \theta_{y}=0
$$

for $\theta_{y}$ on the form (2.49), with $g_{l}-s$ expressed in derivatives of $y$.
A list of symmetric powers of second order equations is easily produced, and particular hierarchies of this sort will be investigated in Section 6.2. We may immediately deduce the following result concerning solutions of such an hierarchy of equations.

Theorem 2.3.1. Given a set of fundamental solutions $\{u, v\}$ of a second order equation corresponding to a $\mathcal{D}$-module $(M, \delta)$. Then

$$
\left\{u^{k}, u^{k-1} v, \ldots, u v^{k-1}, v^{k}\right\}=\left\{u^{i} v^{k-i}\right\}_{i=0}^{k}
$$

is a fundamental set of solutions of the equation corresponding to the $\mathcal{D}$ module ( $\left.S^{k}(M), \delta\right)$ for any $k \geq 2$.

Proof. We may freely choose $(E, \delta)$ such that $(M, \delta)=\left(E^{*}, \delta\right)$ as in Proposition 2.3.2. Given the solutions $\{u, v\}$ we know that $\left\{\theta_{u}, \theta_{v}\right\}$ span $\left(E^{*}\right)^{\#}$, thus

$$
\left\{\theta_{u}^{k}, \theta_{u}^{k-1} \cdot \theta_{v}, \ldots, \theta_{u} \cdot \theta_{v}^{k-1}, \theta_{v}^{k}\right\}
$$

span $S^{k}\left(E^{*}\right)^{\#}$ over $\mathbb{R}$. Also, $S^{k}\left(E^{*}\right)^{\#}$, is closed with respect to the symmetric product, and we need only collect the $\alpha_{1}=\left(e_{1}^{*}\right)^{k}$ term in products $\theta_{u}^{k-i} \cdot \theta_{v}^{i}$ to state that

$$
\begin{equation*}
\theta_{u}^{k-i} \cdot \theta_{v}^{i}=\theta_{u^{k-i} v^{i}} \tag{2.51}
\end{equation*}
$$

hence, $\left\{u^{i} v^{k-i}\right\}_{i=0}^{k}$ span the solution space of $S^{k}\left(E^{*}\right)$.

### 2.4 Geometric image of ODEs in Jet space

Consider a vector bundle $B \xrightarrow{\beta} \mathbb{R}$ of rank $m$ with its $\mathcal{A}=C^{\infty}(\mathbb{R})$-module of sections $C^{\infty}(\beta)=\left\{s \in C^{\infty}(\mathbb{R}, \mathcal{E}) \mid \beta \circ s=I_{\mathbb{R}}\right\}$. The corresponding bundle $J^{k}(\beta) \xrightarrow{\pi_{k}} \mathbb{R}$ of $k$-jets of sections of $\beta$ is of rank $m(k+1)$ over $\mathbb{R}$, and is equipped with the Cartan distribution.

A system of linear $k$-th order ordinary differential equations is a linear subbundle

$$
\begin{equation*}
\mathcal{E} \xrightarrow{\alpha} \mathbb{R} \subset J^{k}(\beta) \xrightarrow{\pi_{k}} \mathbb{R} \tag{2.52}
\end{equation*}
$$

of codimension $m$ such that the Cartan distribution on $J^{k}(\beta)$ when restricted to $\mathcal{E}$, and denoted $\mathcal{C}_{\mathcal{E}}$
(i) is 1-dimensional, and
(ii) projects isomorphically to $\mathbb{R}$

We denote the $\mathcal{A}$-module of sections in the bundle $\alpha$ by $C^{\infty}(\alpha)$. We have a linear connection in the bundle $\alpha$,

$$
\begin{equation*}
\nabla: \mathcal{D}(\mathbb{R}) \longrightarrow \operatorname{Der}\left(C^{\infty}(\alpha)\right) \tag{2.53}
\end{equation*}
$$

where $\operatorname{Der}\left(C^{\infty}(\alpha)\right)$ denotes derivations of $C^{\infty}(\alpha)$ over $\frac{d}{d x}$, i. e. $\mathbb{R}$-linear maps

$$
D: C^{\infty}(\alpha) \rightarrow C^{\infty}(\alpha)
$$

such that

$$
D(f s)=f^{\prime} \cdot s+f \cdot D(s)
$$

for any $f \in \mathcal{A}, s \in C^{\infty}(\alpha)$.
$\nabla$ is defined by the requirement that it lifts $\frac{d}{d x}$ on the base $\mathbb{R}$ to a generator $X \in \mathcal{D}(\mathcal{E})$ of $\mathcal{C}_{\mathcal{E}}$ on $\mathcal{E}$. Consider $\left.s=\underline{s}(x) \in C^{\infty}(\alpha)\right)$ as a curve in $\mathcal{E}$. Then, geometrically, $\nabla_{\frac{d}{d x}}$ on acts on $s$ by

$$
\begin{equation*}
\bar{A}_{-t} \circ s \circ A_{t}=s+\nabla_{\frac{d}{d x}}(s) \cdot t+o(t) \tag{2.54}
\end{equation*}
$$

where $A_{t}$ is the flow generated by $\frac{d}{d x}$ on $\mathbb{R}$, and $\bar{A}_{t}$ is the flow generated by $X$ on $\mathcal{E}$. Thus, constant sections of $\nabla$, i.e. sections $s$ such that

$$
\nabla_{Y}(s)=0, \quad \forall Y \in \mathcal{D}(\mathbb{R})
$$

are precisely the integral curves of $\mathcal{C}_{\mathcal{E}}$ on $\mathcal{E}$.

The pair $\left(C^{\infty}(\alpha), \delta=\nabla\left(\frac{d}{d x}\right)\right)$ is a $\mathcal{D}$-module over $\left(\mathcal{A}, \delta_{\mathcal{A}}\right)$, and we have the correspondence.

$$
\left(C^{\infty}(\alpha), \nabla_{\frac{d}{d x}}\right) \Leftrightarrow(E, \delta) \in O b(\mathcal{L O D E})
$$

If we wish to study a single $k$ th order linear equation

$$
\begin{equation*}
y^{(k)}+f_{k} y^{(k-1)}+\ldots+f_{1} y=0 \tag{2.55}
\end{equation*}
$$

the corresponding linear bundle is

$$
\mathcal{E} \xrightarrow{\alpha} \mathbb{R} \quad \subset \quad J^{k}(\mathbb{R}) \xrightarrow{\pi_{k}} \mathbb{R}
$$

where

$$
\mathcal{E}=\left\{p_{k}=-f_{k} p_{k-1}-\ldots-f_{1} p_{0}\right\} \subset J^{k}(\mathbb{R})
$$

with coordinates $\left(x, p_{0}, \ldots, p_{k-1}\right)$, taking standard coordinates $\left(x, p_{0}, \ldots p_{k}\right)$ on $J^{k}(\mathbb{R})$.
Denote $F\left(x, p_{0}, \ldots, p_{k-1}\right)=-f_{k} p_{k-1}-\ldots-f_{1} p_{0}$. The vector field

$$
\mathcal{D}=\partial_{x}+p_{1} \partial_{p_{0}}+\ldots p_{k-1} \partial_{p_{k-2}}+F \partial_{p_{k-1}}
$$

is a generator of the Cartan distribution on $\mathcal{E}$, and its integral curves are on the form

$$
\phi(x)=\left(x, y(x), y^{\prime}(x), \ldots, y^{(k-1)}(x)\right)
$$

where $y=y(x)$ is a solution of equation(2.55).
Here $\delta=\nabla_{\frac{d}{d x}}$, where $\nabla$ lifts $\frac{d}{d x}$ on the base to $\mathcal{D}$ in the bundle $\mathcal{E} \xrightarrow{\alpha} \mathbb{R}$.

### 2.5 Differential operator view on $\mathcal{D}$-modules.

There is a third way to approach $\mathcal{D}$-modules that correspond to linear ODEs, and this view is an important adding on the practical side in the study of equations. It will be a tool for working with the modules in practice, and important for the use of symmetries to solve equations.

Let $\mathcal{K}$ be the ring of linear differential operators over $\mathbb{R}$. An operator

$$
\begin{equation*}
P=\partial^{k}+c_{k}(x) \partial^{k-1}+\ldots+c_{1}(x) \tag{2.56}
\end{equation*}
$$

where $\partial=\frac{d}{d x}$ defines a $\mathcal{D}$-module $\left(E_{P}, \delta\right) \in \operatorname{Ob}(\mathcal{L O D E})$ with

$$
\begin{equation*}
E_{P} \stackrel{\text { def }}{=} \mathcal{K} /(\mathcal{K} \circ P), \tag{2.57}
\end{equation*}
$$

and

$$
\begin{equation*}
\delta: E_{P} \longrightarrow E_{P} \quad \text { defined by } \delta:[X] \mapsto[\partial \circ X] . \tag{2.58}
\end{equation*}
$$

Obviously this operation is well defined with respect to choice of representative $X \in \mathcal{K}$ modulo $(\mathcal{K} \circ P)$, and it is a derivation over $\partial$. For $P$ as above

$$
\begin{equation*}
\left\{e_{1}=e=[1], e_{2}=\delta e=[\partial], \ldots, e_{k}=\delta^{k-1} e=\left[\partial^{k-1}\right]\right\} \tag{2.59}
\end{equation*}
$$

is a primitive element basis of $E_{P}$ over $\mathcal{A}$. In $E_{P}$

$$
\begin{equation*}
\left[\partial^{k}\right] \equiv-c_{1} e-\ldots-c_{k} \delta^{k-1} e \tag{2.60}
\end{equation*}
$$

thus

$$
\operatorname{ker} \delta \cong \operatorname{ker} P^{t}
$$

To return to the situation in Section 2.3, considering an equation

$$
\begin{equation*}
L(y)=\left(\partial^{k}+c_{k}(x) \partial^{k-1}+\ldots+c_{1}(x)\right)(y)=0 \tag{2.61}
\end{equation*}
$$

we get the corresponding $\mathcal{D}$-module $(E, \delta)$ from taking

$$
\begin{equation*}
E=E_{L^{t}} \tag{2.62}
\end{equation*}
$$

so that

$$
\begin{equation*}
\operatorname{ker} \delta \cong \operatorname{ker} L \tag{2.63}
\end{equation*}
$$

## Chapter 3

## Classic Geometries and ODEs

Theorem 2.2.3 gives us a way to transfer the classical structures from $\mathbb{R}$-vector spaces to $\mathcal{D}$-modules, and hence to ODEs. We study equations equipped with the classic Euclidean, Symplectic, Complex and Hermitian structures.

### 3.1 Euclidean structures

Definition 3.1.1. By a harmonic oscillator we mean a $\mathcal{D}$-module $(E, \delta) \in O b(\mathcal{L O D E})$ equipped with an $\delta$-invariant positive symmetric 2 -form $g \in\left(S^{2}\left(E^{*}\right)\right)^{\#}$.

Theorem 3.1.1. (1) For any linear $O D E$ there exists a quadratic 1st integral.
(2) Any two Harmonic Oscillators $(E, \delta, g)$ and $\left(E^{\prime}, \delta^{\prime}, g^{\prime}\right)$ of the same dimension are equivalent in the sense that there exists an isomorphism

$$
A: E \longrightarrow E^{\prime}
$$

such that
(i) $A \circ \delta=\delta^{\prime} \circ A$ and
(ii) $g^{\prime}(A x, A y)=g(x, y)$

Proof. (1) An ODE of degree $n$ specifies a module $E$ of dimension $n$ as in Theorem 2.2.3, hence there exists a $\delta$-invariant basis $\left\{\gamma_{i}\right\}$ of $E$ as described in the proof of the theorem. Then $g=\left(\gamma_{1}^{*}\right)^{2}+\ldots+\left(\gamma_{n}^{*}\right)^{2} \in S^{2}\left(E^{\#}\right)$ and it is obviously positive definite. $g(h, h)=c$ is our quadratic first integral, for
$h \in E^{\#}$.
(2) Let $\left\{\gamma_{i}\right\}$ and $\left\{\gamma_{i}^{\prime}\right\}$ be basises of $E^{\#}$ and $\left(E^{\prime}\right)^{\#}$ respectively as in Theorem 2.2.3. In these basises $g$ and $g^{\prime}$ are given by orthogonally diagonalizable $\mathbb{R}$ matrices $G$ and $G^{\prime}$. Let $\left\{\eta_{i}^{*}\right\}$ and $\left\{\eta_{i}^{\prime *}\right\}$ be basises of $\left(E^{*}\right)^{\#}$ and $\left(E^{* *}\right)^{\#}$ such that $G$ and $G^{\prime}$ are diagonal. Then the map $A$ is given by

$$
A: \eta_{i} \mapsto\left(\frac{g\left(\eta_{i}, \eta_{i}\right)}{g^{\prime}\left(\eta_{i}^{\prime}, \eta_{i}^{\prime}\right)} \cdot\right)^{1 / 2} \eta_{i}^{\prime}
$$

Since $\left\{\eta_{i}\right\}$ and $\left\{\eta_{i}^{\prime}\right\}$ are basises of $E$ and $E^{\prime}$ over $\mathcal{A}$, expand $A$ as an $\mathcal{A}$ homomorphism $E \rightarrow E^{\prime}$.

Before moving to more specific results on Euclidean structures we include the following property of $\delta$-invariant symmetric bilinear forms.

Proposition 3.1.1. Given a $\mathcal{D}$-module $(E, \delta)$ in the category $\mathcal{L O D E}$. For any $g \in S^{2}\left(E^{*}\right)$, and arbitrary $X, Y \in E$ the following are equivalent

$$
\begin{equation*}
\delta g=0 \Leftrightarrow g(\delta X, Y)+g(X, \delta Y)=g(X, Y)^{\prime} . \tag{3.1}
\end{equation*}
$$

Proof. Given $g=\sum_{i, j=1}^{n} g_{i j} e_{i}^{*} \cdot e_{j}^{*}$ in $S^{2}\left(E^{*}\right)$ and arbitrary $X, Y \in E$. Then

$$
\begin{align*}
\delta g(X, Y) & =\left[\sum g_{i j}^{\prime} e_{i}^{*} \cdot e_{j}^{*}+\sum g_{i j}\left(\delta e_{i}^{*} \cdot e_{j}^{*}+e_{i}^{*} \cdot \delta e_{j}^{*}\right)\right](X, Y) \\
& =\sum\left(g_{i j}^{\prime} e_{i}^{*}(X) \cdot e_{j}^{*}(Y)+g_{i j}\left[\delta\left(e_{i}^{*}(X)\right) e_{j}^{*}(Y)+e_{i}^{*}(X) \cdot \delta\left(e_{j}^{*}(Y)\right)\right]\right) \\
& -\sum g_{i j}\left[e_{i}^{*}(\delta X) e_{j}^{*}(Y)+e_{i}^{*}(X) e_{j}^{*}(\delta Y)\right] \\
& =g(X, Y)^{\prime}-g(\delta X, Y)-g(X, \delta Y) . \tag{3.2}
\end{align*}
$$

### 3.2 Euclidean equations of second order

We will take a closer look at 2nd order equations and Euclidean structures. Consider a general equation of second order

$$
\begin{equation*}
y^{\prime \prime}+\left(a_{2}(x) y\right)^{\prime}-a_{1}(x) y=0 \tag{3.3}
\end{equation*}
$$

corresponding to a $\mathcal{D}$-module $(E, \delta)$ with primitive element basis $\left\{e_{1}=e, e_{2}=\delta e\right\}$ where $\delta$ is described by

$$
\delta^{2} e=a_{1} e+a_{2} \delta e .
$$

We want to study the induced module $\left(S^{2}\left(E^{*}\right), \delta\right)$ and look for positive $\delta$ invariant symmetric forms. Taking the dual basis $\left\{e_{1}^{*}, e_{2}^{*}\right\}$ of $E^{*}$ we recall that the induced $\delta$ in the dual module is given by

$$
\delta \underline{e}^{*}=\left[\begin{array}{cc}
0 & -a_{1}  \tag{3.4}\\
-1 & -a_{2}
\end{array}\right] \underline{e}^{*} .
$$

Constructing a basis $\left\{\left(e_{1}^{*}\right)^{2}, e_{1}^{*} \cdot e_{2}^{*},\left(e_{2}^{*}\right)^{2}\right\}$ of $S^{2}\left(E^{*}\right)$ by taking symmetric products in the basis elements of $E^{*}$ and calculating the induced $\delta$ gives us a full description of $\left(S^{2}\left(E^{*}\right), \delta\right)$.

$$
\delta:\left[\begin{array}{c}
\left(e_{1}^{*}\right)^{2}  \tag{3.5}\\
e_{1}^{*} \cdot e_{2}^{*} \\
\left(e_{2}^{*}\right)^{2}
\end{array}\right] \mapsto\left[\begin{array}{ccc}
0 & -2 a_{1} & 0 \\
-1 & -a_{2} & -a_{1} \\
0 & -2 & -2 a_{2}
\end{array}\right]\left[\begin{array}{c}
\left(e_{1}^{*}\right)^{2} \\
e_{1}^{*} \cdot e_{2}^{*} \\
\left(e_{2}^{*}\right)^{2}
\end{array}\right]
$$

Thus the system of equations $\left(S^{2}\left(E^{*}\right), \delta\right)$ is

$$
\begin{array}{ll}
s_{1}^{\prime}-s_{2} & =0 \\
s_{2}^{\prime}-2 a_{1} s_{1}-a_{2} s_{2}-2 s_{3} & =0 \\
s_{3}^{\prime}-a_{1} s_{2}-2 a_{2} s_{3} & =0 \tag{3.8}
\end{array}
$$

for

$$
g=s_{1}(x)\left(e_{1}^{*}\right)^{2}+s_{2}(x) e_{1}^{*} \cdot e_{2}^{*}+s_{3}(x)\left(e_{2}^{*}\right)^{2} \in S^{2}\left(E^{*}\right)
$$

So, obviously, we may attack the problem directly, and see that the system (3.6) - (3.8) can be resolved into a single governing equation

$$
\begin{equation*}
s^{\prime \prime \prime}+\left(-3 a_{2}\right) s^{\prime \prime}+\left(2 a_{2}^{2}-a_{2}^{\prime}-4 a_{1}\right) s^{\prime}+\left(4 a_{1} a_{2}-2 a_{1}^{\prime}\right) s=0 \tag{3.9}
\end{equation*}
$$

by setting $s_{1}=s(x)$. Equation (3.6) implies that

$$
s_{2}=s_{1}^{\prime}=s^{\prime}
$$

and

$$
s_{3}=\frac{1}{2}\left(s^{\prime \prime}-a_{2} s^{\prime}-2 a_{1} s\right)
$$

by (3.7). Then (3.8) becomes (3.9), which we will denote the 2 nd symmetric power of the equation $\left(E^{*}, \delta\right)$. We may conclude that any element $g$ in the kernel $S^{2}\left(E^{*}\right)^{\#} \subset S^{2}\left(E^{*}\right)$ is on the form

$$
\begin{equation*}
g=s\left(e_{1}^{*}\right)^{2}+s^{\prime}\left(e_{1}^{*} \cdot e_{2}^{*}\right)+\frac{1}{2}\left(s^{\prime \prime}-a_{2} s^{\prime}-2 a_{1} s\right)\left(e_{2}^{*}\right)^{2}, \tag{3.10}
\end{equation*}
$$

where $s=s(x)$ is a solution of (3.9).
There is a second approach to the quest of finding $\delta$-invariant symmetric bilinear forms of an equation; we may use Proposition 3.1.1 to deduce properties of positive, symmetric bilinear forms on a general second order equation (3.3). Let $\left\{e_{1}^{*}, e_{2}^{*}\right\}$ be the dual basis of the primitive element basis $\{e, \delta e\}$ of $(E, \delta)$ as before. Consider $g=\sum_{i, j=1}^{2} g_{i j} e_{i}^{*} \cdot e_{j}^{*}$ in $S^{2}\left(E^{*}\right)$. We will require throughout that $g$ is positive.

Step 1 We may start with the assumption that $g$ is normalized on the primitive element, i. e.

$$
g(e, e)=g_{11}(x)=1
$$

We have the requirement that $g$ is positive, so if $g(e, e)=\alpha^{2}(x)>0, \alpha^{2} \neq 1$, we may perform a change of primitive element basis

$$
\begin{equation*}
\widetilde{e}=\frac{1}{\alpha} e . \tag{3.11}
\end{equation*}
$$

Then $\delta \widetilde{e}=-\frac{\alpha^{\prime}}{\alpha^{2}} e+\frac{1}{\alpha} \delta e$. Writing the transformation in matrix form yields

$$
\left[\begin{array}{c}
\widetilde{e}  \tag{3.12}\\
\delta \widetilde{e}
\end{array}\right]=\left[\begin{array}{cc}
\frac{1}{\alpha} & 0 \\
-\frac{\alpha^{\prime}}{\alpha^{2}} & \frac{1}{\alpha}
\end{array}\right]\left[\begin{array}{c}
e \\
\delta e
\end{array}\right] .
$$

The determinant of the transformation matrix is $\frac{1}{\alpha^{2}(x)} \neq 0$, thus this is a change of basis. On the level of equations a transformation that changes the primitive element by a non-zero factor as above corresponds to a change of variable transformation of (3.3):

$$
y=\frac{1}{\alpha} u .
$$

Thus we may assume that $g$ is normalized in $e$, up to a change of variable in the original ODE.

Step 2 For $g(e, e)=1$ applying Proposition 3.1.1 immediately determines $g(e, \delta e)$ by

$$
g(e, \delta e)+g(\delta e, e)=2 g(e, \delta e)=g(e, e)^{\prime}=0
$$

hence, $g(e, \delta e)=0$.
Step 3 By positivity of $g$ we have that $g(\delta e, \delta e)=\omega^{2}(x)$ for some nonzero $\omega(x)$. Using Proposition 3.1.1 again we get a requirement on $g(\delta e, \delta e)$ by

$$
g(\delta e, \delta e)+g\left(e, \delta^{2} e\right)=g(e, \delta e)^{\prime}=0
$$

But

$$
g\left(e, \delta^{2} e\right)=g\left(e, a_{1} e+a_{2} \delta e\right)=a_{1} g(e, e)+a_{2} g(e, \delta e)=a_{1}
$$

thus we get the requirement $\omega^{2}=-a_{1}$.
Step 4 The last relation we are able to get from applying the proposition determines a relation between $\omega$ and $a_{2}$. First we have

$$
g\left(\delta e, \delta^{2} e\right)+g\left(\delta^{2} e, \delta e\right)=2 g\left(\delta e, \delta^{2} e\right)=2 g\left(\delta e, a_{1} e+a_{2} \delta e\right)=2 a_{2} g(\delta e, \delta e)
$$

and due to Proposition 3.1.1,

$$
2 a_{2} g(\delta e, \delta e)=2 a_{2} \omega^{2}=g(\delta e, \delta e)^{\prime}=2 \omega^{\prime} \omega
$$

that is, $a_{2}=\frac{\omega^{\prime}}{\omega}$.
We may sum this up as follows. In matrix form, i. e. $g=\underline{e}^{* T} G \underline{e}^{*}, g$ is given by

$$
G=\left[\begin{array}{cc}
1 & 0 \\
0 & \omega^{2}
\end{array}\right]
$$

Proposition 3.2.1. An equation

$$
\begin{equation*}
y^{\prime \prime}+\left(\frac{\omega^{\prime}}{\omega} y\right)^{\prime}+\omega^{2} y=0 \tag{3.13}
\end{equation*}
$$

with $\omega \neq 0$ has a quadratic first integral

$$
\begin{equation*}
q=\left[\omega^{2}+\left(\frac{\omega^{\prime}}{\omega}\right)^{2}\right] y^{2}+2\left[\frac{\omega^{\prime}}{\omega}\right] y y^{\prime}+\left[y^{\prime}\right]^{2} . \tag{3.14}
\end{equation*}
$$

Proof. The calculations above determine that equation (3.13) has an associated positive, symmetric bilinear form $g=\left(e_{1}^{*}\right)^{2}+\omega^{2}\left(e_{2}^{*}\right)^{2} \in S^{2}\left(E^{*}\right)^{\#}$ on its solution space $E^{\#}$. That is, for any $h \in E^{\#}, g(h, h)=c$, constant. But any element $h \in E^{\#}$ is on the form

$$
\begin{equation*}
h=\left(-y^{\prime}-a_{2} y\right) e+y \delta e, \tag{3.15}
\end{equation*}
$$

where $y$ is a solution of (3.13). Thus, setting $q=g(h, h)$ gives the desired quadratic first integral.

There is another question to be considered here, namely, how to transform one equation with "potential" $\omega$ into another with "potential" $\widetilde{\omega}$ ? Theorem 3.1.1 in the beginning of this chapter shows the existence of a transformation between any two harmonic oscillators preserving the $\delta$-invariant Euclidean structure. But the construction in the proof depends on knowing solutions of our two equations, and deals only with existence. The following result for second order equations determines a transformation independent of knowing any solutions of the equations.

Theorem 3.2.1. Let $\left(E_{\omega}, \delta\right)$ be the $\mathcal{D}$-module corresponding to the equation

$$
\begin{equation*}
y^{\prime \prime}+\left(\frac{\omega^{\prime}}{\omega} y\right)^{\prime}+\omega^{2} y=0 \tag{3.16}
\end{equation*}
$$

The transformation of primitive element bases

$$
T_{\theta}:\left[\begin{array}{c}
e  \tag{3.17}\\
\delta e
\end{array}\right] \mapsto\left[\begin{array}{c}
\widetilde{e} \\
\delta \widetilde{e}
\end{array}\right]=\left[\begin{array}{cc}
\cos \theta & \frac{1}{\omega} \sin \theta \\
-\left(\omega+\theta^{\prime}\right) \sin \theta & \frac{\left(\omega+\theta^{\prime}\right)}{\omega} \cos \theta
\end{array}\right]\left[\begin{array}{c}
e \\
\delta e
\end{array}\right]
$$

yields an equation $\left(E_{\widetilde{\omega}}, \delta\right)$

$$
y^{\prime \prime}+\left(\frac{\widetilde{\omega}^{\prime}}{\widetilde{\omega}^{\prime}} y\right)^{\prime}+\widetilde{\omega}^{2} y=0
$$

where $\widetilde{\omega}=\omega+\theta^{\prime}$. The associated Euclidean structure is, in matrix form,

$$
\widetilde{G}=\left[\begin{array}{cc}
1 & 0 \\
0 & \widetilde{\omega}^{2}
\end{array}\right]
$$

and is non-degenerated for $\widetilde{\omega}=\omega+\theta^{\prime} \neq 0$.

Proof. The transformation can be divided into three steps. First one transforms $\{e, \delta e\}$ into a basis orthonormal with respect to $g$.

$$
\left[\begin{array}{l}
v_{1} \\
v_{2}
\end{array}\right]=\left[\begin{array}{cc}
1 & 0 \\
0 & \frac{1}{\omega}
\end{array}\right]\left[\begin{array}{c}
e \\
\delta e
\end{array}\right]
$$

The orthonormal basis is then "rotated" by an "angle" $\theta=\theta(x)$

$$
\left[\begin{array}{l}
w_{1} \\
w_{2}
\end{array}\right]=\left[\begin{array}{cc}
\cos \theta(x) & \sin \theta(x) \\
-\sin \theta(x) & \cos \theta(x)
\end{array}\right]\left[\begin{array}{l}
v_{1} \\
v_{2}
\end{array}\right]
$$

Now, taking $w_{1}$ as new primitive element we find, by applying $\delta$ that

$$
\delta w_{1}=\left(\theta^{\prime}+\omega\right) w_{2}
$$

The total transformation is

$$
T_{\theta}=\left[\begin{array}{cc}
1 & 0 \\
0 & \omega+\theta^{\prime}
\end{array}\right]\left[\begin{array}{cc}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right]\left[\begin{array}{ll}
1 & 0 \\
0 & \frac{1}{\omega}
\end{array}\right]=\left[\begin{array}{cc}
\cos \theta & \frac{1}{\omega} \sin \theta \\
-\left(\omega+\theta^{\prime}\right) \sin \theta & \frac{\left(\omega+\theta^{\prime}\right)}{\omega} \cos \theta
\end{array}\right]
$$

and has determinant $\frac{\left(\omega+\theta^{\prime}\right)}{\omega}$, which is non-zero for $\left(\omega+\theta^{\prime}\right) \neq 0$. Thus, if you wish to transform an equation $E_{\omega}$ to another $E_{\widetilde{\omega}}$, you need a $\theta(x)$ such that $\theta^{\prime}=\widetilde{\omega}-\omega$.

Example 3.2.1. The equation

$$
\begin{equation*}
u^{\prime \prime}=0 \tag{3.18}
\end{equation*}
$$

has the first quadratic integral

$$
\begin{equation*}
u^{2}-2 x u u^{\prime}+\left[x^{2}+1\right]\left(u^{\prime}\right)^{2} \tag{3.19}
\end{equation*}
$$

It is obtained by a change of variable $y=u \omega^{-1 / 2}$ from the equation

$$
\begin{equation*}
y^{\prime \prime}+\left(\frac{\omega^{\prime}}{\omega} y\right)^{\prime}+\omega^{2} y=0 \tag{3.20}
\end{equation*}
$$

with

$$
\begin{equation*}
\omega=\frac{1}{x^{2}+1} . \tag{3.21}
\end{equation*}
$$

### 3.3 3rd order Euclidean equations

Consider a third order equation

$$
\begin{equation*}
y^{\prime \prime \prime}+\left(a_{3} y\right)^{\prime \prime}-\left(a_{2} y\right)^{\prime}+a_{1} y=0 \tag{3.22}
\end{equation*}
$$

with corresponding $\mathcal{D}$-module $(E, \delta)$ and primitive element basis

$$
\left\{e, \delta e, \delta^{2} e\right\}
$$

with $\delta^{3} e=a_{1} e+a_{2} \delta e+a_{3} \delta_{2} e$. Using Proposition 3.1.1 repeatedly we can derive requirements for a symmetric bilinear form $g \in S^{2}\left(E^{*}\right)$ to be $\delta$-invariant. Additional requirements determine the coefficients $a_{i}$ in terms of $\alpha, \beta$, as in the step by step calculations leading to Proposition 3.2.1. Assuming that $g$ is normalized in $e$, written in matrix form, $g$ has to be on the form

$$
G=\left[\begin{array}{ccc}
1 & 0 & -\alpha^{2}  \tag{3.23}\\
0 & \alpha^{2} & \alpha^{\prime} \alpha \\
-\alpha^{2} & \alpha^{\prime} \alpha & \beta^{2}
\end{array}\right]
$$

where $\alpha=\alpha(x) \neq 0$ and $\beta=\beta(x) \neq 0$.
Theorem 3.3.1.

$$
\begin{equation*}
y^{\prime \prime \prime}+\left(a_{3} y\right)^{\prime \prime}-\left(a_{2} y\right)^{\prime}+a_{1} y=0 \tag{3.24}
\end{equation*}
$$

has a $\delta$-invariant symmetric bilinear form given by $G$ as in (3.23) for

$$
\begin{align*}
& a_{1}=-\left(\alpha^{2}\right)^{\prime}-\alpha^{2} \frac{\lambda^{\prime}}{\lambda}, \\
& a_{2}=\frac{\alpha^{\prime \prime}}{\alpha}+\frac{\alpha^{\prime}}{\alpha} \frac{\lambda^{\prime}}{\lambda}-\frac{\beta^{2}}{\alpha^{2}},  \tag{3.25}\\
& a_{3}=\frac{\alpha^{\prime}}{\alpha}-\frac{\lambda^{\prime}}{\lambda},
\end{align*}
$$

where

$$
\lambda=g(v, v)^{-1 / 2}=\left(\beta^{2}-\alpha^{4}-\left(\alpha^{\prime}\right)^{2}\right)^{-1 / 2}
$$

for

$$
v=\alpha^{2} e+\frac{\alpha^{\prime}}{\alpha} \delta e+\delta^{2} e,
$$

which is orthogonal to both e and $\delta e$.

### 3.4 Symplectic structures

We may equally study equations with symplectic structure on the solution space.

Definition 3.4.1. A symplectic equation is a $\mathcal{D}$-module $(E, \delta) \in O b(\mathcal{L O D E})$ of even rank $2 n$ equipped with a non-degenerated $\delta$-invariant $\omega \in \bigwedge^{2}\left(E^{*}\right)^{\#}$.

Thus, if we seek a symplectic structure on the solution space of an equation $(E, \delta)$ we should investigate the second exterior power of $E^{*}$.

### 3.4.1 Equations of second order.

Consider a second order equation

$$
\begin{equation*}
y^{\prime \prime}+\left(a_{2} y\right)^{\prime}-a_{1} y=0 \tag{3.26}
\end{equation*}
$$

corresponding to the $\mathcal{D}$-module $(E, \delta)$, with the usual primitive element basis $\left\{e_{1}=e, e_{2}=\delta e\right\}, \delta^{2} e=a_{1} e+a_{2} \delta e$. A general element in $\wedge^{2}\left(E^{*}\right)$ is on the form

$$
\begin{equation*}
\omega=\alpha(x) e_{1}^{*} \wedge e_{2}^{*} \tag{3.27}
\end{equation*}
$$

for some $\alpha \in \mathcal{A}$. Applying $\delta$ to $\omega$ yields

$$
\begin{align*}
\delta \omega & =\alpha^{\prime} e_{1}^{*} \wedge e_{2}^{*}+\alpha\left[\left(-a_{1} e_{2}^{*}\right) \wedge e_{2}^{*}+e_{1}^{*} \wedge\left(-e_{1}^{*}-a_{2} e_{2}^{*}\right)\right] \\
& =\left(\alpha^{\prime}-a_{2} \alpha\right) e_{1}^{*} \wedge e_{2}^{*} \tag{3.28}
\end{align*}
$$

so $\left(\bigwedge^{2}\left(E^{*}\right), \delta\right)$ corresponds to the equation

$$
\begin{equation*}
\alpha^{\prime}-a_{2} \alpha=0 \tag{3.29}
\end{equation*}
$$

The 2 -form is non-degenerated if and only if $\alpha \neq 0$. We may interpret this as a requirement on the coefficient $a_{2}$, and sum up as follows.

Theorem 3.4.1. For $\alpha(x) \neq 0$, any $a_{1}(x)$ in $\mathcal{A}$, the equation

$$
\begin{equation*}
y^{\prime \prime}+\left(\frac{\alpha^{\prime}}{\alpha} y\right)^{\prime}-a_{1} y=0 \tag{3.30}
\end{equation*}
$$

is a symplectic equation with the $\delta$-invariant 2-form

$$
\begin{equation*}
\omega=\alpha(x) e_{1}^{*} \wedge e_{2}^{*} \tag{3.31}
\end{equation*}
$$

determining the symplectic structure on $E^{\#}$.

Recall that an element $h_{y} \in E^{\#}$ corresponding to a solution $y$ of (3.30) is on the form $h_{y}=\left(-y^{\prime}-\frac{\alpha^{\prime}}{\alpha} y\right) e_{1}+y e_{2}$. Given two solutions $y_{1}, y_{2}$,

$$
\begin{align*}
\omega_{\alpha}\left(h_{y_{1}}, h_{y_{2}}\right) & =\alpha\left[\left(-y_{1}^{\prime}-\frac{\alpha^{\prime}}{\alpha} y_{1}\right) y_{2}-\left(-y_{2}^{\prime}-\frac{\alpha^{\prime}}{\alpha} y_{2}\right) y_{1}\right]  \tag{3.32}\\
& =\alpha\left(y_{1}^{\prime} y_{2}-y_{1} y_{2}^{\prime}\right)
\end{align*}
$$

### 3.4.2 Symplectic 4th order equations

Denote by $(E, \delta)$ the D-module corresponding to the 4 th order equation

$$
\begin{equation*}
y^{(4)}+\left(a_{4}(x) y\right)^{(3)}-\left(a_{3}(x) y\right)^{\prime \prime}+\left(a_{2}(x) y\right)^{\prime}-a_{1}(x) y=0 \tag{3.33}
\end{equation*}
$$

and let $\left\{e_{i}=\delta^{i-1} e\right\}_{i=1}^{4}$ be the associated primitive element basis. A $\delta$ invariant non-degenerated form $\omega \in \bigwedge^{2}\left(E^{*}\right)$ provides a symplectic structure on the solution space of the equation. Consider a general element in $\bigwedge^{2}\left(E^{*}\right)$,

$$
\begin{equation*}
\omega=s_{1}(x) e_{1}^{*} \wedge e_{2}^{*}+s_{2}(x) e_{1}^{*} \wedge e_{3}^{*}+\ldots+s_{6}(x) e_{3}^{*} \wedge e_{4}^{*} \tag{3.34}
\end{equation*}
$$

in matrix form $\omega$ is

$$
M_{\omega}=\left[\begin{array}{cccc}
0 & s_{1} & s_{2} & s_{3}  \tag{3.35}\\
-s_{1} & 0 & s_{4} & s_{5} \\
-s_{2} & -s_{4} & 0 & s_{6} \\
-s_{3} & -s_{5} & -s_{6} & 0
\end{array}\right]
$$

To require

$$
\delta \omega=0
$$

yields equations

$$
\begin{align*}
& s_{1}=s(x)  \tag{3.36}\\
& s_{2}=s^{\prime}  \tag{3.37}\\
& s_{3}=s^{\prime \prime}-u  \tag{3.38}\\
& s_{4}=u(x)  \tag{3.39}\\
& s_{5}=u^{\prime}  \tag{3.40}\\
& s_{6}=u^{\prime \prime}-a_{4} u^{\prime}-a_{3} u+a_{1} s, \tag{3.41}
\end{align*}
$$

where $s(x)$ and $u(x)$ solve the system (3.42) and (3.43):

$$
\begin{equation*}
s^{(3)}-a_{4} s^{\prime \prime}-a_{3} s^{\prime}-a_{2} s-2 u^{\prime}+a_{4} u=0 \tag{3.42}
\end{equation*}
$$

$$
\begin{align*}
& u^{(3)}-2 a_{4} u^{\prime \prime}+\left(a_{4}^{2}-a_{4}^{\prime}-a_{3}\right) u^{\prime}  \tag{3.43}\\
& \quad+\left(a_{2}+a_{3} a_{4}-a_{3}^{\prime}\right) u+2 a_{1} s^{\prime}+\left(a_{1}^{\prime}-a_{1} a_{4}\right) s=0
\end{align*}
$$

We find examples of equations with a symplectic form generated by constant solutions $u, s$ of equations (3.42) and (3.43) simply by investigating (3.42) and (3.43) for $u, s$ set constant. Setting $s=1$, and $u=c \in \mathbb{R}$ reveals the conditions on the coefficients $a_{1}, \ldots a_{4}$ for $\omega$ to be $\delta$-invariant. Equation (3.42) becomes

$$
\begin{equation*}
-a_{2}+c a_{4}=0 \tag{3.44}
\end{equation*}
$$

and (3.43) becomes

$$
\begin{equation*}
\left(a_{2}+a_{3} a_{4}-a_{3}^{\prime}\right) c+\left(a_{1}^{\prime}-a_{1} a_{4}\right)=0 \tag{3.45}
\end{equation*}
$$

Thus,

$$
\begin{equation*}
a_{2}=c a_{4} \tag{3.46}
\end{equation*}
$$

whereas

$$
\begin{equation*}
a_{4}=\frac{c a_{3}^{\prime}-a_{1}}{c^{2}+c a_{3}-a_{1}}=\left(\ln \left(c^{2}+c a_{3}-a_{1}\right)\right)^{\prime} . \tag{3.47}
\end{equation*}
$$

The condition for $\omega$ to be non-degenerated is precisely

$$
\operatorname{det} M_{\omega}=\operatorname{det}\left[\begin{array}{cccc}
0 & 1 & 0 & -c  \tag{3.48}\\
1 & 0 & c & 0 \\
0 & -c & 0 & a_{1}-c a_{3} \\
c & 0 & -a_{1}+c a_{3} & 0
\end{array}\right]=\left(c^{2}+c a_{3}-a_{1}\right)^{2} \neq 0
$$

Thus we know the following, an equation

$$
\begin{equation*}
y^{(4)}+\left(\frac{\alpha^{\prime}}{\alpha} y\right)^{(3)}-\left(a_{3}(x) y\right)^{\prime \prime}+\left(c \frac{\alpha^{\prime}}{\alpha} y\right)^{\prime}-a_{1}(x) y=0 \tag{3.49}
\end{equation*}
$$

with

$$
\alpha=c^{2}+c a_{3}-a_{1} \neq 0
$$

has an inherent non-degenerate $\delta$-invariant symplectic form

$$
\omega=e_{1}^{*} \wedge e_{2}^{*}-c e_{1}^{*} \wedge e_{4}^{*}+c e_{2}^{*} \wedge e_{3}^{*}+\left(a_{1}-c a_{3}\right) e_{3}^{*} \wedge e_{4}^{*} .
$$

### 3.5 Complex and Hermitian structure

There is a natural way to introduce a complex structure on a $\mathcal{D}$-module $(E, \delta)$.
Definition 3.5.1. Let $(E, \delta)$ be a $\mathcal{D}$-module in $\mathcal{L O D E}$. A complex structure on the corresponding equation is a $\delta$-invariant $\mathcal{A}$-endomorphism $J \in E n d_{\mathcal{A}}(E)^{\#}$ such that

$$
J^{2}=-I d_{E}
$$

We may immediately deduce the following.
Proposition 3.5.1. Denote $\widetilde{\mathcal{A}}=\mathbb{C} \otimes_{\mathbb{R}} \mathcal{A}$. We may identify the it with smooth complex valued functions in one real variable, $C^{\infty}(\mathbb{R}, \mathbb{C})$. Given $(E, \delta)$ with a complex structure $J \in E n d_{\mathcal{A}}(E)^{\#}$.
(1) $E$ is an $\widetilde{\mathcal{A}}$-module, which we may denote $\widetilde{E}$, by the following definition:

$$
\begin{equation*}
[u(x)+i \cdot v(x)] e \stackrel{\text { def }}{=} u(x) \cdot e+v(x) \cdot J(e) \tag{3.50}
\end{equation*}
$$

for $u(x), v(x) \in \mathcal{A}$ and $e \in E$.
(2) If $J$ is $\delta$-invariant, i. e. $\delta J=0$, or, equivalently,

$$
\delta \circ J=J \circ \delta .
$$

Then

$$
\begin{equation*}
J^{\#}=\left.J\right|_{E \#}: E^{\#} \rightarrow E^{\#} \tag{3.51}
\end{equation*}
$$

is a complex structure on the vector space $E^{\#}$.
As a digression we may stop to note that $J$ is actually a symmetry of our base equation $E$ that satisfies the extra condition $J^{2}=-1$. Symmetries and the corresponding symmetry equations will be discussed extensively in chapter 4.

Proposition 3.5.2. Given a second order equation with complex structure $J,(E, \delta, J)$, its solution space $E^{\#}$ is isomorphic to $\mathbb{C}$ as vector space, and as a field.

Proof. We know that any basis of $E^{\#}$ is generated by two linearly independent solutions $u_{1}, u_{2}$ of the equation corresponding to $(E, \delta)$. Written in the primitive element basis of $E$ the basis elements are on the form
$h_{u}=\left(-u^{\prime}-a_{2} u\right) e+u \delta e$. Choose $u=u_{1}$. We know that $J\left(E^{\#}\right) \subset E^{\#}$, thus the linear independent set $\left\{h_{u}, J\left(h_{u}\right)\right\}$ is a basis of $E^{\#}$. Now,

$$
\phi: E^{\#} \rightarrow \mathbb{C}
$$

defined by

$$
h_{u} \mapsto 1 \quad \text { and } \quad J\left(h_{u}\right) \mapsto i
$$

and requiring $\mathbb{R}$-linearity, is an isomorphism of vectorspaces. Defining multiplication in $E^{\#}$ accordingly by

$$
\begin{align*}
h_{u}^{2} & =h_{u}, \\
h_{u} \cdot J\left(h_{u}\right) & =J\left(h_{u}\right), \quad \text { and }  \tag{3.52}\\
J\left(h_{u}\right)^{2} & =-h_{u},
\end{align*}
$$

yields that $\phi$ is also a multiplicative homomorphism, i.e.

$$
\phi\left(h_{1} \cdot h_{2}\right)=\phi\left(h_{1}\right) \cdot \phi\left(h_{2}\right),
$$

and thus

$$
E^{\#} \cong \mathbb{C}
$$

as fields by $\phi$.
Definition 3.5.2. Given a $\mathcal{D}$-module $(E, \delta)$ with a complex structure given by $J \in E n d \mathcal{A}_{\mathcal{A}}(E)^{\#}$, a Hermitian structure on $E$ is a $\delta$-invariant positive form $H \in S^{2}\left(E^{*}\right)$ that satisfies the compatibility condition

$$
H(J x, y)=-H(x, J y)=J H(x, y)
$$

for all $x, y \in E$.
Recall that any Hermitian form correponds to a pair $g, \omega$ of Euclidean and symplectic forms satisfying the relation

$$
H=g+i \omega
$$

In the next subsection we will see examples of equations with complex structure and compatible Euclidean and symplectic structures, hence Hermitian structure.

### 3.5.1 Second order equations.

We shall investigate when a second order equation

$$
\begin{equation*}
y^{\prime \prime}+\left(a_{2} y\right)^{\prime}-a_{1} y=0 \tag{3.53}
\end{equation*}
$$

has complex structure. Denote the corresponding $\mathcal{D}$-module $(E, \delta)$, with the usual primitive element basis $\left\{e_{1}=e, e_{2}=\delta e\right\}$. We may identify $E n d_{\mathcal{A}}(E)$ with $E^{*} \otimes E$ and write an endomorphism

$$
\begin{equation*}
J=F_{1}(x) e_{1}^{*} \otimes e_{1}+F_{2}(x) e_{1}^{*} \otimes e_{2}+F_{3}(x) e_{2}^{*} \otimes e_{1}+F_{4}(x) e_{2}^{*} \otimes e_{2}, \tag{3.54}
\end{equation*}
$$

or, in matrix form, $J=\underline{e}^{* T} M_{J}^{T} \underline{e}$,

$$
M_{J}=\left[\begin{array}{ll}
F_{1} & F_{3}  \tag{3.55}\\
F_{2} & F_{4}
\end{array}\right]
$$

Then

$$
M_{J}^{2}=\left[\begin{array}{cc}
F_{1}^{2}+F_{2} F_{3} & F_{3}\left(F_{1}+F_{4}\right)  \tag{3.56}\\
F_{2}\left(F_{1}+F_{4}\right) & F_{4}^{2}+F_{2} F_{3}
\end{array}\right] .
$$

Requiring $M_{J}^{2}=-I$ gives us four equations on the coefficients $F_{i}$, and one immediate requirement is that $F_{1}^{2}=F_{4}^{2}$. Thus, we may split the problem into two cases and get the following classes of endomorphisms.
Class (A), characterized by $F_{1}=F_{4}$ :

$$
J=\left[\begin{array}{cc}
0 & \alpha  \tag{3.57}\\
-\frac{1}{\alpha} & 0
\end{array}\right]
$$

where $\alpha=\alpha(x) \neq 0$.
Adding the requirement that $\delta J=0$ gives four new equations, and for $F_{i}$-s as above they are reduced to

$$
\begin{equation*}
a_{1}(x)=-\alpha^{2} \tag{3.58}
\end{equation*}
$$

and

$$
\begin{equation*}
a_{2}(x)=\frac{\alpha^{\prime}}{\alpha} . \tag{3.59}
\end{equation*}
$$

Class (B), characterized by $F_{1}=-F_{4}$ :

$$
J=\left[\begin{array}{cc}
\alpha & \beta  \tag{3.60}\\
-\frac{1+\alpha^{2}}{\beta} & -\alpha
\end{array}\right],
$$

where $\alpha=\alpha(x)$ is any function, and $\beta=\beta(x) \neq 0$. The requirement $\delta J=0$ with $F_{i}$-s as above is reduced to

$$
\begin{equation*}
a_{1}(x)=-\frac{\beta\left(\beta-\alpha^{\prime}\right)}{1+\alpha^{2}}, \tag{3.61}
\end{equation*}
$$

and

$$
\begin{equation*}
a_{2}(x)=\frac{\beta^{\prime}}{\beta}+\frac{2 \alpha\left(\beta-\alpha^{\prime}\right)}{1+\alpha^{2}} . \tag{3.62}
\end{equation*}
$$

We may sum up as follows.
Theorem 3.5.1. There are two classes of second order equations that possess complex structure.
(A) For $\alpha(x) \neq 0$,

$$
\begin{equation*}
y^{\prime \prime}+\left(\frac{\alpha^{\prime}}{\alpha} y\right)^{\prime}+\alpha^{2} y=0 \tag{3.63}
\end{equation*}
$$

with complex structure $J$ determined by (3.57).
(B) For $\beta(x) \neq 0$ and any $\alpha(x)$,

$$
\begin{equation*}
y^{\prime \prime}+\left[\left(\frac{\beta^{\prime}}{\beta}+\frac{2 \alpha\left(\beta-\alpha^{\prime}\right)}{1+\alpha^{2}}\right) y\right]^{\prime}+\left[\frac{\beta\left(\beta-\alpha^{\prime}\right)}{1+\alpha^{2}}\right] y=0 \tag{3.64}
\end{equation*}
$$

with complex structure $J$ determined by (3.60).
Corollary 3.5.1. The complex structures of class ( $A$ ) and ( $B$ ) as above are symmetries of the respective equations.
(i) For equation

$$
\begin{equation*}
y^{\prime \prime}+\left(\frac{\alpha^{\prime}}{\alpha} y\right)^{\prime}+\alpha^{2} y=0 \tag{3.65}
\end{equation*}
$$

of class $(A), \alpha \neq 0$, the complex structure $J$ acts on solutions $y(x)$ as follows

$$
\begin{equation*}
O_{J}: y \mapsto \frac{(\alpha y)^{\prime}}{\alpha^{2}}=\left(\frac{1}{\alpha}\right) y^{\prime}+\left(\frac{\alpha^{\prime}}{\alpha^{2}}\right) y . \tag{3.66}
\end{equation*}
$$

(ii) For equation

$$
\begin{equation*}
y^{\prime \prime}+\left[\left(\frac{\beta^{\prime}}{\beta}+\frac{2 \alpha\left(\beta-\alpha^{\prime}\right)}{1+\alpha^{2}}\right) y\right]^{\prime}+\left[\frac{\beta\left(\beta-\alpha^{\prime}\right)}{1+\alpha^{2}}\right] y=0 \tag{3.67}
\end{equation*}
$$

of class $(B), \beta \neq 0$, any $\alpha$, the complex structure $J$ acts on solutions $y(x)$ as follows

$$
\begin{equation*}
O_{J}: y \mapsto \frac{\alpha \beta\left(\beta-\alpha^{\prime}\right)+\beta^{\prime}\left(1+\alpha^{2}\right)}{\beta} y+\frac{\left(1+\alpha^{2}\right)}{\beta} y^{\prime} . \tag{3.68}
\end{equation*}
$$

Note that we recognize the form of the class (A) equations, it is precisely the same as the equations with Euclidean structure in Proposition 3.2.1. We know that there is the notion of compatible symplectic structure, i.e. the defining 2-form given by

$$
\begin{equation*}
\omega(x, y)=g(J x, y) \tag{3.69}
\end{equation*}
$$

where $g$ is Euclidean and $J$ is a complex structure. The triple corresponds to a Hermitian structure

$$
H=g+i \omega .
$$

Theorem 3.5.2. The equation

$$
\begin{equation*}
y^{\prime \prime}+\left(\frac{\alpha^{\prime}}{\alpha} y\right)^{\prime}+\alpha^{2} y=0 \tag{3.70}
\end{equation*}
$$

where $\alpha=\alpha(x) \neq 0$ has Euclidean, complex, symplectic and Hermitian structures.

$$
\begin{align*}
g & =\left(e_{1}^{*}\right)^{2}+\alpha^{2}\left(e_{2}^{*}\right)^{2}  \tag{3.71}\\
J & =-\frac{1}{\alpha} e_{1}^{*} \otimes e_{2}+\alpha e_{2}^{*} \otimes e_{1}  \tag{3.72}\\
\omega & =-\alpha e_{1}^{*} \wedge e_{2}^{*}  \tag{3.73}\\
H & =g+i \omega \tag{3.74}
\end{align*}
$$

or, in matrix form,

$$
\begin{aligned}
g & =\left[\begin{array}{cc}
1 & 0 \\
0 & \alpha^{2}
\end{array}\right], \\
J & =\left[\begin{array}{cc}
0 & \alpha \\
-\frac{1}{\alpha} & 0
\end{array}\right], \\
\omega & =\left[\begin{array}{cc}
0 & \alpha \\
-\alpha & 0
\end{array}\right] .
\end{aligned}
$$

Example 3.5.1. The well known equation

$$
\begin{equation*}
y^{\prime \prime}+y=0 \tag{3.75}
\end{equation*}
$$

is on the form as in Theorem 3.5.2 for $\alpha^{2}=1$. We may take $\alpha=-1$ and get the standard complex structure with matrix

$$
J=\left[\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right]
$$

Choosing $\alpha=1$ only alters the sign of $J$. The solution space is spanned by $u=\cos (x), v=\sin (x)$, which in turn determines the basis

$$
\begin{align*}
h_{\cos (x)} & =\sin (x) e+\cos (x) \delta e  \tag{3.76}\\
J\left(h_{\cos (x)}\right) & =-\cos (x) e+\sin (x) \delta e=h_{\sin (x)}
\end{align*}
$$

of $E^{\#}$ over $\mathbb{R}$, and of $E$ over $\mathcal{A}$.

### 3.6 Yang-Baxter solutions

In Section 2.2 we saw that if we have a $\delta$-invariant group action into a $\mathcal{D}$ module, then the invariant elements of this action constitute a $\mathcal{D}$-sub-module, which in turn corresponds to a new ODE. In Theorem 2.2.2 we saw that, given a $\mathcal{D}$-module $(E, \delta)$, any solution $\tau: E \otimes E \rightarrow E \otimes E$ of the YangBaxter equation with the property $\tau^{2}=1$ gives a representation of $S_{n}$ into $E^{\otimes n}$ for all $n \geq 2$. Recall that $\tau$ being plus and minus twist gave us submodules $S^{n}(E)$ and $\bigwedge^{n}(E)$ of $E^{\otimes n}$. The twist operation may be used to construct other solutions of the Yang-Baxter equation.

Proposition 3.6.1. Given a splitting of a $\mathcal{D}$-module $E=E_{0} \oplus E_{1}$, we may introduce the following map on $E^{\otimes 2}$ given by a combination of + and - twisting:

$$
\tau:\left\{\begin{aligned}
& f_{0} \otimes g_{i} \mapsto g_{i} \otimes f_{0}, \\
& f_{i} \otimes g_{0} \mapsto g_{0} \otimes f_{i}, \\
& i=0,1 \\
& f_{1} \otimes g_{1} \mapsto-g_{1} \otimes f_{1}
\end{aligned}\right.
$$

This $\tau$ is a solution if the Yang - Baxter equation (2.30), and thus induces a representation of the symmetric group $S_{n}$ into $E^{\otimes n}$.

Given a $\mathcal{D}$-module $(E, \delta)$ corresponding to the second order equation

$$
\begin{equation*}
y^{\prime \prime}+\left(a_{2} y\right)^{\prime}-a_{1} y=0 \tag{3.77}
\end{equation*}
$$

with primitive element basis $\left\{e_{1}, e_{2}\right\}$, we may try to find a splitting of $E$ by means of an operator $A: E \rightarrow E$ with the property

$$
A^{2}=1
$$

such that $E$ splits into two one-dimensional modules

$$
E=E_{0} \oplus E_{1}=\operatorname{ker}(A-1) \oplus \operatorname{ker}(A+1)
$$

Ensuring that the splitting preserves the $\mathcal{D}$-module structure we require that

$$
\begin{equation*}
\delta A=0 . \tag{3.78}
\end{equation*}
$$

We get two classes of non-trivial splittings.
Class (1) consists of equations

$$
\begin{equation*}
y^{\prime \prime}+\left(\frac{\alpha^{\prime}}{\alpha} y\right)^{\prime}-\alpha^{2} y=0 \tag{3.79}
\end{equation*}
$$

with $\alpha \neq 0$. This equation splits non-trivially into

$$
E=E_{0} \oplus E_{1}=\operatorname{ker}(A-1) \oplus \operatorname{ker}(A+1)
$$

for

$$
A=\alpha e_{2}^{*} \otimes e_{1}+\frac{1}{\alpha} e_{1}^{*} \otimes e_{2}
$$

with $E_{0}=<\eta_{0}=\alpha e_{1}+e_{2}>_{\mathcal{A}}$ and $E_{1}=<\eta_{1}=-\alpha e_{1}+e_{2}>_{\mathcal{A}}$. Restricting $\delta$ to $E_{0}$ and $E_{1}$ yields

$$
\begin{align*}
& \delta \eta_{0}=\left(\frac{\alpha^{\prime}}{\alpha}+\alpha\right) \eta_{0}, \quad \text { and }  \tag{3.80}\\
& \delta \eta_{1}=\left(\frac{\alpha^{\prime}}{\alpha}-\alpha\right) \eta_{1} . \tag{3.81}
\end{align*}
$$

This means that $E^{\#}=\left\{u \eta_{0}+v \eta_{1}=\alpha(u-v) e_{1}+(u+v) e_{2}\right\}$ where $u$ and $v$ solve equations

$$
\begin{array}{ll}
\left(E_{0}, \delta\right): & u^{\prime}+\left(\frac{\alpha^{\prime}}{\alpha}+\alpha\right) u=0 \\
\left(E_{1}, \delta\right): & v^{\prime}+\left(\frac{\alpha^{\prime}}{\alpha}-\alpha\right) v=0 \tag{3.83}
\end{array}
$$

respectively, i.e. sums of solutions $y=u+v$ give all solutions of equation $E$.
But we also have an action into $E \otimes E$ of the symmetric group $S_{2}$ given by $\tau$ as in Proposition 3.6.1. Invariants of this action $\Sigma_{\tau}^{2}(E)=\operatorname{ker}(\tau-1) \subset E \otimes E$ is generated over $\mathcal{A}$ by

$$
\gamma_{1}=\eta_{0} \otimes \eta_{0}, \quad \gamma_{2}=\eta_{0} \otimes \eta_{1}+\eta_{1} \otimes \eta_{0}
$$

Investigating derivatives we get that

$$
\begin{align*}
\delta \gamma_{1} & =\left(2 \frac{\alpha^{\prime}}{\alpha}+2 \alpha\right) \gamma_{1}  \tag{3.84}\\
\delta \gamma_{2} & =\left(2 \frac{\alpha^{\prime}}{\alpha}\right) \gamma_{2} \tag{3.85}
\end{align*}
$$

thus,

$$
\Sigma_{\tau}^{2}(E)=M_{1} \oplus M_{2} \subset E \otimes E
$$

where $\left(M_{1}, \delta\right)$ corresponds to the equation

$$
u^{\prime}+\left(2 \alpha+2 \frac{\alpha^{\prime}}{\alpha}\right) u=0
$$

and $\left(M_{2}, \delta\right)$ corresponds to the equation

$$
u^{\prime}+\left(2 \frac{\alpha^{\prime}}{\alpha}\right) u=0 .
$$

On the other hand, we may take $-\tau$ to generate an action of $S_{2}$, and $\Sigma_{-\tau}^{2}(E)=\operatorname{ker}(\tau+1) \subset E \otimes E$ is generated over $\mathcal{A}$ by

$$
\gamma_{3}=\eta_{1} \otimes \eta_{1}, \quad \gamma_{4}=\eta_{0} \otimes \eta_{1}-\eta_{1} \otimes \eta_{0}
$$

Taking derivatives we get that

$$
\begin{align*}
\delta \gamma_{3} & =\left(2 \frac{\alpha^{\prime}}{\alpha}-2 \alpha\right) \gamma_{3}  \tag{3.86}\\
\delta \gamma_{4} & =\left(2 \frac{\alpha^{\prime}}{\alpha}\right) \gamma_{4} \tag{3.87}
\end{align*}
$$

thus,

$$
\Sigma_{-\tau}^{2}(E)=M_{3} \oplus M_{4}
$$

where $\left(M_{3}, \delta\right)$ corresponds to the equation

$$
u^{\prime}+\left(2 \alpha-2 \frac{\alpha^{\prime}}{\alpha}\right) u=0
$$

and $\left(M_{4}, \delta\right)$ corresponds to the equation

$$
u^{\prime}+(4 \alpha) u=0
$$

We may sum up as follows:

$$
E \otimes E=\Sigma_{\tau}^{2}(E) \oplus \Sigma_{-\tau}^{2}(E)=M_{1} \oplus M_{2} \oplus M_{3} \oplus M_{4}
$$

and

$$
(E \otimes E)^{\#}=\left\{m_{1} \gamma_{1}+m_{2} \gamma_{2}+m_{3} \gamma_{3}+m_{4} \gamma_{4}\right\}
$$

where $m_{i}$ solves equation $M_{i}$.

Class (2) consists of equations

$$
\begin{equation*}
y^{\prime \prime}+\left[\left(\frac{\beta^{\prime}}{\beta}+2 \frac{\alpha\left(\beta-\alpha^{\prime}\right)}{\left(\alpha^{2}-1\right)}\right) y\right]^{\prime}-\frac{\beta\left(\beta-\alpha^{\prime}\right)}{\left(\alpha^{2}-1\right)} y=0 \tag{3.88}
\end{equation*}
$$

with $\beta \neq 0, \alpha \neq-1,1$. It splits with respect to the operator

$$
A=\alpha e_{1}^{*} \otimes e_{1}+\beta e_{2}^{*} \otimes e_{1}+\frac{\left(1-\alpha^{2}\right)}{\beta} e_{1}^{*} \otimes e_{2}-\alpha e_{2}^{*} \otimes e_{2}
$$

into a direct sum $E=E_{0} \oplus E_{1}=\operatorname{ker}(A-1) \oplus E_{1}$ as for class (1)-equations, with $E_{0}=<\eta_{0}=\beta e_{1}+(1-\alpha) e_{2}>_{\mathcal{A}}$, and $E_{1}=<\eta_{0}=\beta e_{1}-(1+\alpha) e_{2}>_{\mathcal{A}}$. The corresponding equations are

$$
\begin{array}{ll}
E_{0}: & u^{\prime}+\left(\frac{\beta^{\prime}}{\beta}+\frac{\left(\beta-\alpha^{\prime}\right)}{(\alpha+1)}\right) u=0, \\
E_{1}: & v^{\prime}+\left(\frac{\beta^{\prime}}{\beta}+\frac{\left(\beta-\alpha^{\prime}\right)}{(\alpha-1)}\right) v=0 . \tag{3.90}
\end{array}
$$

Thus

$$
E^{\#}=\left\{u \eta_{0}+v \eta_{1}=\beta(u+v) e_{1}+[(u-v)-\alpha(u+v)] e_{2}\right\}
$$

so solutions of equation $E$ are $y=(u-v)-\alpha(u+v)$ for solutions $u, v$ of $E_{0}, E_{1}$.
Repeating the study of the representation of $S_{2}$ into $E \otimes E$ as for Class (1) equations yields that

$$
\Sigma_{\tau}^{2}(E)=M_{1} \oplus M_{2} \subset E \otimes E
$$

and

$$
\Sigma_{-\tau}^{2}(E)=M_{3} \oplus M_{4} \subset E \otimes E
$$

with

$$
E \otimes E=\Sigma_{\tau}^{2}(E) \oplus \Sigma_{-\tau}^{2}(E)=M_{1} \oplus M_{2} \oplus M_{3} \oplus M_{4}
$$

The modules $M_{i}$ correspond to the following equations

$$
\begin{array}{lll}
M_{1}: & m_{1}^{\prime} & +2\left(\frac{\left(\beta-\alpha^{\prime}\right)}{(\alpha+1)}+\frac{\beta^{\prime}}{\beta}\right) m_{1}=0 \\
M_{2}: & m_{2}^{\prime}+2\left(\frac{\alpha\left(\beta-\alpha^{\prime}\right)}{\left(\alpha^{2}-1\right)}+\frac{\beta^{\prime}}{\beta}\right) m_{2}=0 \\
M_{3}: & m_{3}^{\prime}+2\left(\frac{\left(\beta-\alpha^{\prime}\right)}{(\alpha-1)}+\frac{\beta^{\prime}}{\beta}\right) m_{3}=0 \\
M_{4}: & m_{4}^{\prime}-8 \beta\left(\frac{\alpha\left(\beta-\alpha^{\prime}\right)}{\left(\alpha^{2}-1\right)}+\frac{\beta^{\prime}}{\beta}\right) m_{4}=0 \tag{3.94}
\end{array}
$$

## Chapter 4

## Symmetries and representations

In this chapter we study symmetries of equations, in particular through symmetry operators. Section 4.2 contains results for linear operator symmetries, most of which is discussed in detail in [12]. The most important addition to these results is the description on how this embeds into the category $\mathcal{L O D E}$, through Proposition 4.2.2 and Theorem 4.2.2.

Symmetry operators are convinient tools for calculations with symmetries in $\mathcal{D}$-modules, and are important to ensure full applicability of solving strategies developed in Chapters 5 and 6.

### 4.1 Symmetry algebras and representations

In the category $\mathcal{L O D E}$ a symmetry of an equation $(E, \delta)$ is an endomorphism of $E$ which is $\delta$-invariant, that is an element

$$
\begin{equation*}
X \in \operatorname{End}_{\mathcal{A}}(E)^{\#} \tag{4.1}
\end{equation*}
$$

Such an $X$ is simply a map of the module $E$ into itself such that it maps solutions to solutions,

$$
\begin{equation*}
X: E \rightarrow E \text { such that } X\left(E^{\#}\right) \subset E^{\#} \tag{4.2}
\end{equation*}
$$

We thus arrive at a natural way to introduce symmetry algebras of equations in our picture, in terms of representation theory.

Definition 4.1.1. A Lie algebra $\mathfrak{g}$ is a Lie algebra of linear symmetries of an equation $(E, \delta) \in \operatorname{Ob}(\mathcal{L O D E})$ if there is a representation

$$
\rho: \mathfrak{g} \longrightarrow \operatorname{End}_{\mathcal{A}}(E)
$$

such that

$$
\rho(g) \circ \delta=\delta \circ \rho(g), \quad \forall g \in \mathfrak{g}
$$

i.e. $\rho$ maps $\mathfrak{g}$ into $\delta$-invariant endomorphisms of $E$,

$$
\rho(\mathfrak{g}) \subset \operatorname{End}_{\mathcal{A}}(E)^{\#}
$$

Proposition 4.1.1. If $\mathfrak{g}$ is a symmetry algebra of $(E, \delta)$ with associated representation $\rho$, then

$$
\rho: \mathfrak{g} \rightarrow \operatorname{End}_{\mathbb{R}}\left(E^{\#}\right)
$$

is a representation into the $\mathbb{R}$-vector space $E^{\#}$.
We just need to recall that taking the kernel of $\delta$ commutes with the algebraic constructions in our category, $\operatorname{End}_{\mathcal{A}}(E)^{\#}=\operatorname{End}_{\mathbb{R}}\left(E^{\#}\right)$, ref. Proposition 2.2.2.

The consequences of combining Proposition 2.2.2 and Theorem 2.2.3 are immediate, this enables us to make use of results from the rich theory of representations of Lie algebras into vector spaces. In particular decomposition theorems from the theory concerning semisimple Lie algebras.

### 4.2 Symmetry operators

Most of the results in this section concerning linear differential operators are found in [12].
Recalling from Section 2.5 that we may define our $\mathcal{D}$-modules by using classes of linear operators, we shall in this section embed the action of symmetries in this view. Let $\mathcal{K}$ still denote the ring of linear differential operators over $\mathbb{R}$.

Definition 4.2.1. A symmetry operator of an equation

$$
\begin{equation*}
L(y)=y^{(k)}+c_{k}(x) y^{(k-1)}+\ldots+c_{1}(x) y=0 \tag{4.3}
\end{equation*}
$$

is a linear operator

$$
\Delta=b_{0}(x)+b_{1}(x) \partial+\ldots+b_{k-1}(x) \partial^{k-1} \in \mathcal{K}
$$

with the property that there exists $\nabla \in \mathcal{K}$ such that

$$
\begin{equation*}
L \circ \Delta=\nabla \circ L . \tag{4.4}
\end{equation*}
$$

Note that $\Delta(\operatorname{ker} L) \subset \operatorname{ker} L$, i. e. it maps solutions to solutions.
Remark: Recall the discussion in Section 2.4 on the jet space approach to ODEs. Studying an equation (4.3) we make the following connection. Associated to an operator $\Delta=b_{0}(x)+b_{1}(x) \partial+\ldots+b_{k-1}(x) \partial^{k-1}$ is the function

$$
f_{\Delta}=b_{0} p_{0}+b_{1} p_{1}+\ldots+b_{k-1} p_{k-1} \in C^{\infty}(\mathcal{E})
$$

where $\mathcal{E} \xrightarrow{\alpha} \mathbb{R}$ is the linear subbundle in $J^{k}(\mathbb{R})$ corresponding to (4.3).
The function $f$ solves the Lie equation for (4.3) and generates a shuffling symmetry of the Cartan distribution $\mathcal{C}_{\mathcal{E}}$ on $\mathcal{E} \subset J^{k} \mathbb{R}$ if and only if $\Delta$ is a symmetry in the above sense. See [3] for a more detailed discussion on symmetries of Cartan distributions and their generating functions.

We denote the set of symmetries $\Delta$ by $\operatorname{Sym}(L)$.
Consider more generally

$$
\begin{equation*}
\Sigma(L)=\{P \in \mathcal{K} \mid \exists Q \in \mathcal{K} \text { such that } L \circ P=Q \circ L\} \tag{4.5}
\end{equation*}
$$

Proposition 4.2.1. $\Sigma(L)$ is
(1) an associative $\mathbb{R}$-algebra with respect to composition of operators, and

$$
L \circ\left(P_{1} \circ P_{2}\right)=\left(Q_{1} \circ Q_{2}\right) \circ L .
$$

(2) a Lie algebra with respect to commutators of operators, and

$$
L \circ\left[P_{1}, P_{2}\right]=\left[Q_{1}, Q_{2}\right] \circ L
$$

Lemma 4.2.1. For any $P \in \mathcal{K}$ of order $l, L$ as in (4.3) there are uniquely determined operators $C_{P}$ and $R_{P}$ in $\mathcal{K}$, of order $\leq l-k$ and $\leq k$ respectively such that

$$
\begin{equation*}
P=C_{P} \circ L+R_{P} . \tag{4.6}
\end{equation*}
$$

This is proved by induction and order arguments, as done in [12].
An analougous argument shows that there are unique operators $c_{P}$ and $r_{P}$ such that

$$
\begin{equation*}
P=L \circ c_{P}+r_{P} . \tag{4.7}
\end{equation*}
$$

Theorem 4.2.1. For $P \in \Sigma(L)$ the remainder $R_{P}$ from right division by $L$ as in Lemma 4.2.1 is an element of Sym $(L)$. The map

$$
R: \Sigma(L) \rightarrow \operatorname{Sym}(L) \text { with } R(P)=R_{P}
$$

induces
(1) associative $\mathbb{R}$-algebra structure on $\operatorname{Sym}(L)$ by

$$
m\left(\Delta_{1}, \Delta_{2}\right)=R_{\Delta_{1} \circ \Delta_{2}}
$$

(2) Lie algebra structure

$$
\begin{gathered}
{[\cdot, \cdot]_{L}: \operatorname{Sym}(L) \times \operatorname{Sym}(L) \rightarrow \operatorname{Sym}(L)} \\
\quad \text { by } \quad\left[\Delta_{1}, \Delta_{2}\right]_{L}=R_{\left[\Delta_{1}, \Delta_{2}\right]}
\end{gathered}
$$

Proof. Let $\Delta \in \operatorname{Sym}(L)$ with associated operator $\nabla$ s.t. $L \circ \Delta=\nabla \circ L$. Decompose $\Delta=C_{\Delta} \circ L+R_{\Delta}$ according to Lemma 4.2.1, and likewise $\nabla=$ $L \circ c_{\nabla}+r_{\nabla}$ by left division by $L$. Then the symmetry property implies

$$
L \circ\left(c_{\nabla}-C_{\Delta}\right) \circ L=L \circ R_{\Delta}-r_{\nabla} \circ L .
$$

If $c_{\nabla}-C_{\Delta} \neq 0$ the left hand side is an operator of order $>2 k$, which is impossible since the operator on the right hand side is of maximum order $2 k$, whence $c_{\nabla}=C_{\Delta}$ and $L \circ R_{\Delta}=r_{\nabla} \circ L$, which proves that $R_{\Delta} \in \operatorname{Sym}(L)$. Statements (1) and (2) follow directly.

In the category $\mathcal{L O D \mathcal { E }}$ we viewed symmetries of an equation $E$ as $\delta$-invariant endomorphisms of $E$. We will make the link between symmetries viewed as endomorphisms and symmetry operators in the following way.

Proposition 4.2.2. Let $E=E_{L^{t}}$ be the factor $\mathcal{D}$-module corresponding to an equation

$$
L(y)=0
$$

A symmetry operator $\Delta \in \mathcal{K}$ of the equation induces a $\delta$-invariant endomorphism $\bar{\Delta}$ of $E$

$$
\begin{equation*}
\bar{\Delta}: E \rightarrow E \tag{4.8}
\end{equation*}
$$

defined by

$$
\begin{equation*}
[X] \mapsto\left[X \circ \nabla^{t}\right] \tag{4.9}
\end{equation*}
$$

where $\nabla \in \mathcal{K}$ such that $L \circ \Delta=\nabla \circ L$.
Proof. Note primarily that right composition by $\nabla^{t}$ is well defined with respect to choice of representative modulo $L^{t}$ :

$$
\begin{aligned}
\left(X+A \circ L^{t}\right) \circ & \nabla^{t} \\
& =X \circ \nabla^{t}+A \circ(\nabla \circ L)^{t}= \\
& =X \circ \nabla^{t}+A \circ(L \circ \Delta)^{t}=X \circ \nabla^{t}+\left(A \circ \Delta^{t}\right) \circ L^{t} .
\end{aligned}
$$

Moreover, $\bar{\Delta}$ is an $\mathcal{A}$ homorphism, and obviously commutes with $\delta$ :

$$
(\bar{\Delta} \circ \delta)[X]=\left[\partial \circ X \circ \nabla^{t}\right]=(\delta \circ \bar{\Delta})[X]
$$

Thus, $\bar{\Delta} \in E n d_{\mathcal{A}}(E)^{\#}$.

Theorem 4.2.2. Given a symmetry operator $\Delta$ of the equation $L(y)=0$ the corresponding $\delta$-invariant endomorphism $\bar{\Delta}: E \rightarrow E$ acts as follows when restricted to $\operatorname{ker} \delta=E^{\#} \subset E$ :

$$
\bar{\Delta}:\left[X_{y}\right] \mapsto\left[X_{\Delta(y)}\right]
$$

where $\left[X_{y}\right] \in E^{\#}$ is generated by a solution $y$ of $L(y)=0$.
Proof. We start by noting that for a representative $X$ of a class $[X] \in \operatorname{ker} \delta$ there is an associated operator $A$ such that

$$
\begin{equation*}
\partial \circ X=A \circ L^{t} \tag{4.10}
\end{equation*}
$$

Recall from Section 2.3 on primitive element bases that a solution $y$ of $L(y)=0$, generates an element $\left[X_{y}\right] \in E^{\#}$ on the form

$$
\left[X_{y}\right]=H_{k}(y) e+H_{k-1}(y) \delta e+\ldots+H_{2}(y) \delta^{k-2} e+y \delta^{k-1} e,
$$

that is, with a representative

$$
X_{y}=H_{k}(y)+H_{k-1}(y) \partial+\ldots+H_{2}(y) \partial^{k-2}+y \partial^{k-1} \in \mathcal{K}
$$

with primitive element $e=[1], \delta e=[\partial]$ etc. For choice of representative $X_{y}$ the left hand side in (4.10) is of degree $k$ with highest degree coefficient equal to $y$, whence the operator $A$ is of degree zero, and equal to $y$, that is $\partial \circ X_{y}=y L^{t}$. Further we know that

$$
\begin{equation*}
X_{y} \circ \nabla^{t}=X_{u}+B \circ L^{t} \tag{4.11}
\end{equation*}
$$

for some solution $u$ of $L(y)=0$, and some operator $B$. Again, applying $\delta$ to this representative yields

$$
\partial \circ\left(X_{u}+B \circ L^{t}\right)=(u+\partial \circ B) \circ L^{t} .
$$

But $\partial \circ\left(X_{y} \circ \nabla^{t}\right)=\left(\partial \circ X_{y}\right) \circ \nabla^{t}=\left(y L^{t}\right) \circ \nabla^{t}=y \Delta^{t} \circ L^{t}$, whence

$$
u+\partial \circ B=y \Delta^{t}
$$

i.e. $u+B \circ \partial=\Delta \circ y$. Collecting terms of degree zero yields precisely that

$$
u=\Delta(y)
$$

Remark: Certainly we wish to be able to calculate the action of $\bar{\Delta}$ directly for a known symmetry $\Delta$, whence $\nabla^{t}$ should be directly retrieved from $\Delta$. If we look to the condition (4.4) of $\Delta$ being a symmetry of the equation $L(y)=0$, we see the following: Recall that $L=\partial^{k}+c_{k} \partial^{k-1}+\ldots+c_{1}$. Let

$$
\begin{aligned}
\Delta & =A_{1}+A_{2} \partial+\ldots+A_{k} \partial^{k-1} \\
\nabla & =B_{1}+B_{2} \partial+\ldots+B_{k} \partial^{k-1}
\end{aligned}
$$

On the one hand we get

$$
L \circ \Delta=\left[A_{k}\right] \partial^{2 k-1}+\sum_{l=1}^{k-1}\left[A_{k-l}+\alpha_{k-l}(A, c)\right] \partial^{2 k-l-1}+\sum_{l=1}^{k} \phi_{k-l}(A, c) \partial^{k-l}
$$

where $\alpha_{k-l}$ depends on the coefficient functions $c_{i}(x)$ of $L$, and $A_{j}$-s and their derivatives for $j>k-l$. The functions $\phi_{k-l}$ depend on $A_{i}$-s and $c_{i}$-s. Likewise,

$$
\nabla \circ L=\left[B_{k}\right] \partial^{2 k-1}+\sum_{l=1}^{k-1}\left[B_{k-l}+\beta_{k-l}(B, c)\right] \partial^{2 k-l-1}+\sum_{l=1}^{k} \psi_{k-l}(B, c) \partial^{k-l}
$$

where, similarly $\beta_{k-l}$ depends on the coefficient functions $c_{i}(x)$ of $L$, and $B_{j}$-s and their derivatives for $j>k-l$. The functions $\psi_{k-l}$ depend on $B_{i}$-s and $c_{i}$-s. Thus, by setting $L \circ \Delta=\nabla \circ L$ and collecting terms of the same order in $\partial$ we arrive at $2 k$ equations. The first $k$ equations determine the $B_{i}$-s in terms of $A_{j}$-S and $c_{i}$-S:
(i) $B_{k}=A_{k}$
(ii) $\quad B_{k-1}=A_{k-1}+\alpha_{k-1}\left(A_{k}, c\right)-\beta_{k}\left(B_{k}, c\right)$

$$
(l+1) B_{k-l}=A_{k-l}+\alpha_{k-l}\left(A_{k-l+1}, \ldots, A_{k}, c\right)-\beta_{k-l}\left(B_{k-l+1}, \ldots, B_{k}, c\right)
$$

$$
(k) \quad B_{1} \quad=A_{1}+\alpha_{1}\left(A_{2}, \ldots, A_{k}, c\right)-\beta_{1}\left(B_{2}, \ldots, B_{k}, c\right)
$$

Starting with equation $(i)$ and succesively substituting into the following equations we find the $B_{i}$-s in terms of the coefficients $A_{i}$ of $\Delta$, and the coefficients $c_{i}$ of $L$. Thus $\nabla$, and subsequently $\bar{\Delta}$ is derived directly from $\Delta$. The last $k$ equations are the differential equations

$$
(k+1) \quad \psi_{k-1}(B, c)=\phi_{k-1}(A, c)
$$

$$
(2 k) \quad \psi_{0}(B, c)=\phi_{0}(A, c)
$$

that determine conditions on $A_{j}$-s for $\Delta$ to be a symmetry, Lie equations for $L$.

### 4.2.1 Skew- and self- adjoint equations.

Note that the map

$$
\phi: \operatorname{Sym}(L) \longrightarrow \operatorname{Sym}\left(L^{t}\right)
$$

$$
\begin{equation*}
\Delta \mapsto \nabla^{t} \tag{4.12}
\end{equation*}
$$

is an isomorphism.
Whenever $L$ is skew- or self-adjoint, i.e $L^{t}= \pm L$, we note that

$$
\begin{equation*}
\phi: \Sigma(L) \longrightarrow \Sigma(L) \tag{4.13}
\end{equation*}
$$

and likewise

$$
\begin{equation*}
\phi: \operatorname{Sym}(L) \longrightarrow \operatorname{Sym}(L) \tag{4.14}
\end{equation*}
$$

gives us an involution on symmetries, $\phi^{2}=I d$. Whence the $\Sigma(L)$ and the symmetry space decompose into

$$
\begin{aligned}
\Sigma(L) & =\Sigma_{0}(L) \oplus \Sigma_{1}(L) \\
\operatorname{Sym}(L) & =\operatorname{Sym}_{0}(L) \oplus \operatorname{Sym}_{1}(L)
\end{aligned}
$$

where

$$
\begin{align*}
\Sigma_{0}(L) & =\left\{\Delta \mid L \circ \Delta=\Delta^{t} \circ L\right\}  \tag{4.15}\\
\Sigma_{1}(L) & =\left\{\Delta \mid L \circ \Delta=-\Delta^{t} \circ L\right\} \tag{4.16}
\end{align*}
$$

and $\operatorname{Sym}_{a}(L)=\Sigma_{a}(L) \cap \operatorname{Sym}(L)$.
Theorem 4.2.3. Let $L$ be skew- or self-adjoint. Then

$$
\operatorname{Sym}(L)=\operatorname{Sym}_{0}(L) \oplus \operatorname{Sym}_{1}(L)
$$

is a $\mathbb{Z}_{2}$-graded Lie algebra, i.e.

$$
\left[\operatorname{Sym}_{a}(L), \operatorname{Sym}_{b}(L)\right] \subset \operatorname{Sym}_{a+b}(L)
$$

$a, b \in \mathbb{Z}_{2}$.

### 4.3 Second order equations

We will investigate in detail the symmetry equation of a second order equations, using both the operator approach and direct calculation in the module of endomorphisms of the equation, and thus illustrate both methods. Consider the equation

$$
\begin{equation*}
L(y)=y^{\prime \prime}+\left(a_{2} y\right)^{\prime}-a_{1} y=0 \tag{4.17}
\end{equation*}
$$

The corresponding module $(E, \delta)$ has primitive element basis $\left\{e_{1}=e, e_{2}=\delta e\right\}$ with $\delta^{2} e=a_{1} e+a_{2} \delta e$. The module of endomorphisms $\left(E n d_{\mathcal{A}}(E), \delta\right)$ may be identified with $\left(E^{*} \otimes E, \delta\right)$, and we may write a general endomorphism

$$
\begin{equation*}
F=F_{1}(x) e_{1}^{*} \otimes e_{1}+F_{2}(x) e_{1}^{*} \otimes e_{2}+F_{3}(x) e_{2}^{*} \otimes e_{1}+F_{4}(x) e_{2}^{*} \otimes e_{2} \tag{4.18}
\end{equation*}
$$

Thus

$$
\delta F=0
$$

if and only if the coefficient functions $F_{i}$ satisfy the system

$$
\begin{array}{ccrrrr}
\text { (i) } & F_{1}^{\prime} & & +a_{1} F_{2} & -F_{3} & =0  \tag{4.19}\\
\text { (ii) } & F_{2}^{\prime} & +F_{1} & +a_{2} F_{2} & & -F_{4}
\end{array}=0
$$

Adding (i) and (iv) yields

$$
F_{4}^{\prime}=-F_{1}^{\prime}
$$

Integrating, we get that

$$
F_{4}=-F_{1}+c
$$

for some constant $c \in \mathbb{R}$. Denoting

$$
\begin{equation*}
F_{2}=p(x) \tag{4.20}
\end{equation*}
$$

equation (ii) implies that

$$
\begin{equation*}
F_{1}=\frac{1}{2}\left(c-p^{\prime}-a_{2} p\right), \tag{4.21}
\end{equation*}
$$

and thus,

$$
\begin{equation*}
F_{4}=\frac{1}{2}\left(c+p^{\prime}+a_{2} p\right) \tag{4.22}
\end{equation*}
$$

Finally, from ( $i$ ) we get that

$$
\begin{equation*}
F_{3}=\frac{1}{2}\left(-p^{\prime \prime}-\left(a_{2} p\right)^{\prime}+2 a_{1} p\right) \tag{4.23}
\end{equation*}
$$

and (iii) becomes

$$
\begin{equation*}
p^{\prime \prime \prime}+\left(2 a_{2}^{\prime}-a_{2}^{2}-4 a_{1}\right) p^{\prime}+\left(a_{2}^{\prime \prime}-a_{2} a_{2}^{\prime}-2 a_{1}^{\prime}\right) p=0 \tag{4.24}
\end{equation*}
$$

This may be summed up as follows. In matrix form, $F$ is given by

$$
M_{F}=\left[\begin{array}{cc}
\frac{1}{2}\left(c-p^{\prime}-a_{2} p\right) & a_{1} p-\frac{1}{2}\left(\left(a_{2} p\right)^{\prime}+p^{\prime \prime}\right)  \tag{4.25}\\
p & \frac{1}{2}\left(c+p^{\prime}+a_{2} p\right)
\end{array}\right]
$$

where $p=p(x)$ solves equation (4.24) and $c \in \mathbb{R}$.
Note that this tells us that, as a $\mathcal{D}$-module,

$$
\operatorname{End}_{A}(E) \cong E_{0} \oplus E_{1}
$$

where $\left(E_{0}, \delta_{0}\right)$ and $\left(E_{1}, \delta_{1}\right)$ are the $\mathcal{D}$-modules corresponding to equations

$$
s^{\prime}=0
$$

and

$$
p^{\prime \prime \prime}+\left(2 a_{2}^{\prime}-a_{2}^{2}-4 a_{1}\right) p^{\prime}+\left(a_{2}^{\prime \prime}-a_{2} a_{2}^{\prime}-2 a_{1}^{\prime}\right) p=0
$$

respectively. $\left(E_{0}, \delta_{0}\right)$ contributes with the trivial part of our invariant $F$, the constant $c$. This part $\frac{c}{2}\left(e_{1}^{*} \otimes e_{1}+e_{2}^{*} \otimes e_{2}\right)$ only acts by multiplying an element in $E$ by $\frac{c}{2}$.

We will go through calculations once more, but this time by the operator approach. Our equation is given by

$$
L=\partial^{2}+\partial a_{2}-a_{1}
$$

Recall that a first order linear operator

$$
P=P_{1}+P_{2} \partial
$$

is a symmetry of the equation $L y=0$ if there is an operator

$$
Q=Q_{1}+Q_{2} \partial
$$

such that

$$
\begin{equation*}
L \circ P=Q \circ L \tag{4.26}
\end{equation*}
$$

Setting $L \circ P=Q \circ L$ and collecting terms of the same order gives four equations, of which the two first determine the functions $Q_{1}$ and $Q_{2}$ in terms of $P_{1}$ and $P_{2}$, as promised in Section 4.2 :

$$
\begin{align*}
& Q_{2}=P_{2},  \tag{4.27}\\
& Q_{1}=P_{1}-a_{2} Q_{2}+2 P_{2}^{\prime}+a_{2} P_{2} \tag{4.28}
\end{align*}
$$

thus

$$
\begin{equation*}
Q_{1}=P_{1}+2 P_{2}^{\prime} \tag{4.29}
\end{equation*}
$$

The two last equations are

$$
\begin{array}{ccc}
P_{2}^{\prime \prime}+\left(a_{2} P_{2}\right)^{\prime}+a_{1} P_{2}+2 P_{1}^{\prime}+a_{2} P_{1}-a_{2} Q_{1}-a_{1} Q_{2}-2 a_{2}^{\prime} Q_{2} & =0 \\
P_{1}^{\prime \prime} & +\left(a_{2} P_{1}\right)^{\prime}+a_{1} P_{1}-a_{2}^{\prime \prime} Q_{2}-a_{1}^{\prime} Q_{2}-a_{2}^{\prime} Q_{1}-a_{1} Q_{1} & =0 . \tag{4.31}
\end{array}
$$

They become

$$
\begin{align*}
P_{2}^{\prime \prime}-\left(a_{2} P_{2}\right)^{\prime}+2 P_{1}^{\prime} & =0,  \tag{4.32}\\
P_{1}^{\prime \prime}+a_{2} P_{1}^{\prime}-2 a_{1} P_{2}^{\prime}-2 a_{2}^{\prime} P_{2}^{\prime}-\left(a_{1}^{\prime}+a_{2}^{\prime \prime}\right) & =0, \tag{4.33}
\end{align*}
$$

of which the first may be integrated to give us

$$
\begin{equation*}
P_{1}=\frac{1}{2}\left(a_{2} P_{2}-P_{2}^{\prime}\right)+c \tag{4.34}
\end{equation*}
$$

for $c \in \mathbb{R}$. Setting $P_{2}=p(x)$ we arrive at the following. Any symmetry operator of equation (4.17) is on the form

$$
\begin{equation*}
P=c+\frac{1}{2}\left(a_{2} p-p^{\prime}\right)+p \partial \tag{4.35}
\end{equation*}
$$

where $p=p(x)$ solves

$$
\begin{equation*}
p^{\prime \prime \prime}+\left(2 a_{2}^{\prime}-a_{2}^{2}-4 a_{1}\right) p^{\prime}+\left(a_{2}^{\prime \prime}-a_{2} a_{2}^{\prime}-2 a_{1}^{\prime}\right) p=0 . \tag{4.36}
\end{equation*}
$$

This equation is precisely equation (4.24), which we arrived at when considering endomorphisms of our equation.

Again,

$$
\operatorname{Sym}(L) \cong S y m_{t r} \oplus S_{t y m}^{e q}
$$

where $S_{y m} m_{t r}$ is the trivial part, i.e $P=c \in \mathbb{R}$, a solution of the equation $s^{\prime}=0$.

There is another property to note from the non-trivial symmetry equation (4.24). It is the symmetric 2-power, $S^{2}\left(E_{0}\right)$, of the equation

$$
\begin{equation*}
L_{0}=\partial^{2}+W(x), \tag{4.37}
\end{equation*}
$$

where

$$
\begin{equation*}
W(x)=\frac{1}{2} a_{2}^{\prime}-\frac{1}{4} a_{2}^{2}-a_{1} . \tag{4.38}
\end{equation*}
$$

From Theorem 2.3.1 we know that if $\{u, v\}$ is a set of fundamental solution $\{u, v\}$ of the equation

$$
\begin{equation*}
L_{0}(u)=0 \tag{4.39}
\end{equation*}
$$

then $\left\{u^{2}, u v, v^{2}\right\}$ are fundamental solutions of

$$
\begin{equation*}
S^{2}\left(L_{0}\right)(p)=0 \tag{4.40}
\end{equation*}
$$

Thus, we may produce symmetries from solutions of

$$
\begin{equation*}
u^{\prime \prime}+\left(\frac{1}{2} a_{2}^{\prime}-\frac{1}{4} a_{2}^{2}-a_{1}\right) u=0 \tag{4.41}
\end{equation*}
$$

by operators $P_{u^{2}}, P_{u v}$ and $P_{v^{2}}$.

## Chapter 5

## Solvable symmetry algebras and quadratures

It is a general opinion that to solve an ODE by quadratures with the use of symmetries we will need a solvable algebra of symmetries of dimension equal to the order of the equation. In [3] it is shown that knowing a solvable $k$-dimensional transversal Lie algebra of symmetries of a $k$ th order ODE one can find the general solution by quadratures. In this geometric approach the method consists of finding a complete set of first integrals of the Cartan distribution of the equation by integrating closed 1 -forms and solving functional equations, where the solvability of the algebra is crucial to recover the appropriate 1-forms.

Our approach here is somewhat different, and we shall see that whether we are able to solve an ODE directly by quadratures is not dependent on the order of the equation or the dimension of its symmetry algebra, but rather on eigenvalues of symmetries viewed as endomorphisms of the corresponding $\mathcal{D}$-module. Thus it may happen that a single symmetry is sufficient to solve an equation, conditions for this are stated in in Theorem 5.1.1.

In Theorem 5.2.5 we have stated a condition for an equation with a solvable symmetry algebra to be solved directly by quadratures, with no requirement on the dimension of the algebra.

### 5.1 Decomposition of equations by eigenspaces of symmetries

We begin this chapter with a result that should be kept in mind whenever working with symmetries of equations. It is not limited to any particular type of symmetry algebra, and even states that a single symmetry may be enough to solve an equation by quadratures, regardless of order of the equation, the only factor being eigenspaces of the action of the symmetry.

Given an equation $(E, \delta)$ with a symmetry $X \in \operatorname{End}_{\mathcal{A}}(E)^{\#}$, for $\lambda \in \mathbb{R}$ denote

$$
\begin{equation*}
E_{\lambda}=\{h \in E \mid X(h)=\lambda h\}, \tag{5.1}
\end{equation*}
$$

For a non-empty $E_{\lambda}$ we call $\lambda$ a eigenvalue of $X$.
Proposition 5.1.1. $E_{\lambda}$ is a sub-D-module of $E$.
Proof. $E_{\lambda}$ is obviously a sub-module of $E$. For $h \in E$,

$$
X(\delta h)=\delta(X(h))=\delta(\lambda h)=\lambda \cdot \delta h
$$

since $X$ commutes with $\delta$, thus $\delta\left(E_{\lambda}\right) \subset E_{\lambda}$, and $E_{\lambda}$ is a $\mathcal{D}$-module.
The rank of the module $E_{\lambda}$ over $\mathcal{A}$ we will call the multiplicity of $\lambda$.
Theorem 5.1.1. Let $(E, \delta)$ be a $\mathcal{D}$-module of rank $n$, and $X$ a symmetry of $E$. If $X$ has $n$ distinct eigenvalues $\lambda_{1}, \ldots, \lambda_{n}$, all of multiplicity 1 , then

$$
\begin{equation*}
E=E_{\lambda_{1}} \oplus \ldots \oplus E_{\lambda_{n}}, \tag{5.2}
\end{equation*}
$$

and $E$ is solvable by quadratures.
Proof. The $E_{\lambda_{i}}$ are all non-empty by definition, and are sub- $\mathcal{D}$-modules of $E$. Since the eigenvalues are distinct, the $E_{i}$-s don't intersect, and they span the whole of $E$, since there are $n$. Each $E_{\lambda_{i}}$ corresponds to a first order equation which is identified by applying $\delta$ to an arbitrary element $w_{i} \in E_{\lambda_{i}}$. If $\delta w_{i}=\alpha_{i} w_{i}$, then the corresponding equation is

$$
u_{i}^{\prime}+\alpha_{i} u_{i}=0
$$

and elements $e_{i}=u_{i} w_{i}$, where the $u_{i}$ solve the $E_{\lambda_{i}}$-equations span the solution space $E^{\#} \subset E$.

The Theorem describes a situation where we get a maximal decomposition of the module $E$ by pure algebraic calculations, due to the fact that multiplicities equal 1 for all eigenvalues. We may encounter situations where we have eigen-module decomposition of $E$ with multiplicities of eigenvalues larger than 1.

Theorem 5.1.2. Let $(E, \delta)$ be a $\mathcal{D}$-module of rank $n$, and $X$ a symmetry of $E$. If $X$ has $k$ distinct eigenvalues $\lambda_{1}, \ldots, \lambda_{k}$, of multiplicities $m_{i}$ respectively, with $\sum_{i=1}^{k} m_{i}=n$, then

$$
\begin{equation*}
E=E_{\lambda_{1}} \oplus \ldots \oplus E_{\lambda_{k}} \tag{5.3}
\end{equation*}
$$

Proof. The $E_{\lambda_{i}}$ are non-empty, non-intersecting sub- $\mathcal{D}$-modules of $E$. The sum of their ranks over $\mathcal{A}$ equals the rank of $E$, thus they span the whole of $E$ and their direct sum equals $E$.

Note: Knowing a decomposition of $E$ as in Theorem 5.1.2 corresponds to knowing a set of equations of lower order whose solution spaces span the solution space of the original equation.

Assume that we know a decomposition of $E$ into modules $E_{\lambda_{j}}$, i.e. we know a set of "eigenvectors" $\left\{w_{1}^{j}, \ldots, w_{m_{j}}^{j}\right\}$ in $E$ that span $E_{\lambda_{j}}$. To identify the equation corresponding to $\left(E_{\lambda_{i}}, \delta\right)$, simply apply $\delta$ to the $w_{i}^{j}$. If

$$
\delta \underline{w}^{j}=A_{j}(x) \underline{w}^{j}
$$

for some matrix $A_{j}(x)$, then the corresponding ODE is the $m_{j} \times m_{j}$ system

$$
\begin{equation*}
\underline{f}^{\prime}+A_{j}^{T} \underline{f}=0 \tag{5.4}
\end{equation*}
$$

If $\underline{f}_{i}^{j}, i=1, \ldots, m_{j}$ are solutions of the respective $\left(E_{\lambda_{j}}, \delta\right)$ systems, then $E^{\#}$ is spanned by

$$
E^{\#}=\operatorname{Span}_{\mathbb{R}}\left\{\underline{f}_{i}^{j} \cdot \underline{w}^{j T}, i=1, . ., m_{j}, j=1, . ., k\right\} \subset E .
$$

### 5.2 Solving and decomposition procedures for solvable algebras

We turn to study $\mathcal{D}$-modules $(E, \delta)$ with a solvable symmetry algebra. Following the nature of solvable Lie algebras, we shall describe a procedure
whose aim is to identify a chain $E_{i} \subset E_{i-1}$ of sub $\mathcal{D}$-modules in $E$ of codimension 1 at each step so as to solve the total equation by combined algebraic operations and quadratures.

We begin by recalling the Lie Theorem for representations of solvable Lie algebras.

Theorem 5.2.1. (Lie) Let $\mathfrak{g}$ be a solvable Lie algebra over a base field $F$, $\operatorname{char}(F)=0$ and $F$ algebraically closed. Given a representation $\mathfrak{g} \rightarrow \operatorname{End}_{F}(V)$, $V \neq 0$ a finite dimensional vector space over $F$, there exists a non-zero $v \in V$ such that it is a common eigenvector for the whole action of $\mathfrak{g}$. I. e.

$$
\rho(g)(v)=\lambda(g) v, \quad \forall g \in \mathfrak{g}
$$

for a weight $\lambda \in \mathfrak{g}^{*}$.
As seen earlier, structures on the vector space level $V=E^{\#}$ can be lifted to the $\mathcal{D}$-module $E$, due to Theorem 2.2 .3 . To apply Lie's Theorem in full generality we need the base field, $F$, to be algebraically closed, so if needed we may assume that we work with $\mathbb{C}$-valued smooth real functions, i.e with $A_{C}=C^{\infty}(\mathbb{R}, \mathbb{C})=\mathbb{C} \otimes_{\mathbb{R}} A, \delta_{A_{C}}$ being the usual derivative in $x$. The analogue of the Lie Theorem for $\mathcal{D}$-modules then reads as follows.

Theorem 5.2.2. ( $\mathcal{D}$-Lie) Let $\mathfrak{g}$ be a solvable algebra under the conditions in Theorem 5.2.1, and a symmetry algebra of $(E, \delta)$, a $\mathcal{D}$-module over $\left(A_{C}, \delta_{A_{C}}\right)$. Then there exists $\beta \in E$ and $\lambda \in \mathfrak{g}^{*} \otimes \mathbb{C}$ such that

$$
\rho(g)(\beta)=\lambda(g) \beta, \quad \forall g \in \mathfrak{g} .
$$

Proof. Applying Lie's Theorem to the representation

$$
\rho: \mathfrak{g} \rightarrow \operatorname{End}_{\mathbb{C}}(V)
$$

with $V=E^{\#}$, ensures that there exists an element $v \in V$ which is a common eigenvector of the dual representation, with a corresponding $\lambda \in \mathfrak{g}^{*} \otimes \mathbb{C}$. Any multiple $\beta=f(x) v \in E, f(x) \in A_{\mathbb{C}}$ satisfies $g(\beta)=\lambda(g) \beta$, the action of $\mathfrak{g}$ being linear in functions.
Definition 5.2.1. Given a symmetry algebra $\mathfrak{g}$ of an equation $(E, \delta)$ with a corresponding representation $\rho: \mathfrak{g} \rightarrow \operatorname{End}_{A_{C}}(E)$. An element $\lambda \in \mathfrak{g}^{*} \otimes \mathbb{C}$ such that the associated sub-module

$$
\begin{equation*}
E_{\lambda}=\{\beta \in E \mid \rho(g)(\beta)=\lambda(g) \beta \forall g \in \mathfrak{g}\} \subset E \tag{5.5}
\end{equation*}
$$

is non-empty, is called a weight of the representation. The rank of $E_{\lambda}$ over $A_{C}$ is called the multiplicity of $\lambda$, and $E_{\lambda}$ the associated eigen-sub module.

Proposition 5.2.1. For a weight $\lambda$ of a representation of a symmetry algebra $\mathfrak{g}$ of an equation $(E, \delta), E_{\lambda}$ is a sub-D-module of $E$.
Proof. The proof is almost identical to the proof of Proposition 5.1.1. $E_{\lambda}$ is obviously a sub-module, and $\delta\left(E_{\lambda}\right) \subset E_{\lambda}$ since

$$
\rho(g)(\delta h)=\delta(\rho(g)(h))=\delta(\lambda(g) h)=\lambda(g) \cdot \delta h
$$

for any $g \in \mathfrak{g}, h \in E_{\lambda}$. Thus $E_{\lambda}$ is a $\mathcal{D}$-module.
The next step will be to recognize that the existence of a common eigenvector of the dual representation of $\mathfrak{g}$ into $E^{*}$ ensures that there exists a sub-module $E_{1} \subset E$ of codimension 1 .

Corollary 5.2.1. There exists $\alpha \in E^{*}$ and $\lambda \in \mathfrak{g}^{*} \otimes \mathbb{C}$ such that

$$
\rho^{*}(g)(\alpha)=\lambda(g) \cdot \alpha \quad \forall g \in \mathfrak{g}
$$

If the multiplicity of $\lambda, i$. e. $\operatorname{rank}_{\mathcal{A}_{C}} E_{\lambda}^{*}$, is 1 , then

$$
\begin{equation*}
E^{\prime}=\operatorname{ker} \alpha \subset E \tag{5.6}
\end{equation*}
$$

is a codimension 1 sub-D-module of $E$, which is stable under the representation of $\mathfrak{g}$.
Proof. The existence of $\alpha$ follows directly from Theorem 5.2.2. $\operatorname{ker} \alpha=$ $\operatorname{Ann}\left(E_{\lambda}^{*}\right)$ is a sub-module in $E$, due to the $A_{C}$-linearity of $\alpha: E \rightarrow A_{C}$. Also,

$$
\begin{equation*}
<\alpha, \rho(g) h>=<\rho^{*}(g) \alpha, h>=\lambda(g)<\alpha, h> \tag{5.7}
\end{equation*}
$$

so $\rho(g)(\operatorname{ker} \alpha) \subset \operatorname{ker} \alpha$ for all $g \in \mathfrak{g}$. $E_{\lambda}^{*}$ is of rank 1 , so $\alpha$ spans $E_{\lambda}^{*}$, with

$$
\delta(\alpha)=f(x) \alpha
$$

for some $f(x) \in A_{C}$. By definition

$$
\delta \alpha=\alpha \circ \delta-\delta \circ \alpha
$$

so, combining the two yields

$$
\alpha(\delta h)=(\delta \alpha)(h)+\delta(\alpha(h))=(f(x)+\delta)(\alpha(h))
$$

Whence, $h \in \operatorname{ker} \alpha$ implies that $\delta(h) \in \operatorname{ker} \alpha$, thus $\operatorname{ker} \alpha$ is a $\mathcal{D}$-module.

Theorem 5.2.3. Given a $\mathcal{D}$-module $(E, \delta)$ with a sub-D-module

$$
E^{\prime} \subset E
$$

of codimension 1 for which we know a full set of solutions. Then we can solve the whole of $E$ by quadratures.

Proof. Let $h_{1}, \ldots h_{n-1} \in E^{\#}$ be a basis of $E^{\prime}$. Pick any element $h_{n} \in E \backslash E^{\prime}$, $\delta h_{n}=\sum_{i=1}^{n} f_{i}(x) h_{i}$. For a general element $s=\sum_{i=1}^{n} s_{i}(x) h_{i} \in E$ we get

$$
\delta s=\sum_{i=1}^{n}\left(s_{i}^{\prime}+f_{i} s_{n}\right) h_{i}
$$

Thus $\delta s=0$ if and only if

$$
\begin{array}{cl}
s_{1}^{\prime}+f_{1} s_{n} & =0 \\
s_{2}^{\prime}+f_{2} s_{n} & =0 \\
\vdots & \\
s_{n}^{\prime}+f_{n} s_{n} & =0
\end{array}
$$

i. e. $\quad s_{i}^{\prime}=-f_{i} v$ where $v$ solves the last equation, $v^{\prime}+f_{n} v=0$, and $\left\{s, h_{1}, \ldots, h_{n-1}\right\} \subset E^{\#}$ is a basis of the whole module $E$ over $A$.

So, if we have a procedure for stepwise identifying sub-modules of codimension 1 , starting with $E_{1} \subset E$, we can solve the equation by quadratures.

Theorem 5.2.4. Let $(E, \delta)$ be a $\mathcal{D}$-module with a solvable Lie algebra of symmetries. Then there exists a filtration of $E$ by sub-D-modules

$$
0 \subset E_{n} \subset \ldots \subset E_{2} \subset E_{1} \subset E
$$

where $E_{i+1} \subset E_{i}$ is of codimension 1 at each step.
Proof. Applying Lie's Theorem for vector spaces to the dual representation

$$
\rho^{*}: \mathfrak{g} \rightarrow E n d_{\mathbb{C}}\left(V^{*}\right)
$$

ensures that there exists an element $v^{*} \in V^{*}, V^{*}=\left(E^{*}\right)^{\#}$, which is a common eigenvector of the dual representation, with a corresponding $\lambda \in \mathfrak{g}^{*} \otimes \mathbb{C}$. It is clear that

$$
V_{1}=\operatorname{ker} v^{*} \subset V
$$

is a sub-space of codimension 1 of $V$. Denote by $E_{1}$ the sub-module of $E$ generated by $\operatorname{kerv}^{*} \subset V^{*}$ over $\mathcal{A}_{C}$,

$$
E_{1}=\left\{\sum f_{i}(x) v_{i} \mid f_{i} \in \mathcal{A}_{C}, v_{i} \in \operatorname{ker} v^{*}\right\}
$$

The module $E_{1}$ is in fact a $\mathcal{D}$-module, since $\delta\left(E_{1}\right) \subset E_{1}$, by

$$
\delta\left(\sum f_{i} v_{i}\right)=\sum\left(f_{i}^{\prime} v_{i}+f_{i} \delta v_{i}\right)=\sum f_{i}^{\prime} v_{i} \in E_{1}
$$

for a general element $\sum f_{i} v_{i}$ in $E_{1}$.
Moreover, $E_{1}$, as well as $\operatorname{ker} v^{*}$, is stable under the action of $\mathfrak{g}$ :

$$
<v^{*}, \rho(g)(w)>=<\rho^{*}(g) v^{*}, w>=\lambda(g)<v^{*}, w>,
$$

so $w \in \operatorname{ker} v^{*} \subset V$ implies that $\rho(g)(w) \in \operatorname{ker} v^{*}$, for any $g \in \mathfrak{g}$. The $\mathcal{A}_{C}$-linearity of $v^{*}$ gives the same result for $E_{1}$, thus the representation of $\mathfrak{g}$ restricts to $E_{1}$. Repeating the procedure $n=\operatorname{rank}(E)$ times, starting with $E_{1}$, proves that the desired filtration exists.

Note to Theorem 5.2.4: Theorem 5.2.4 merely says something about the existence of such a filtration, it uses the underlying vectorspace $V^{*}$, which is of course, in general unavailable to us, since knowing it corresponds to having solved the equation in the first place. In practice we will always work with the representation into the module $E$. The main obstruction in this algorithm to reduce the problem to quadratures, is to get codimensions 1 for the desired sub- $\mathcal{D}$-modules with pure algebraic tools.

Theorem 5.2.5. Let $(E, \delta)$ be a $\mathcal{D}$-module of rank $n$ with a solvable symmetry algebra $\mathfrak{g}$. If there are $n$ distinct weights $\lambda_{1}, \ldots, \lambda_{n}$ of multiplicity 1 of the dual representation of $\mathfrak{g}$ into $E^{*}$, then a filtration

$$
0 \subset E_{n} \subset \ldots \subset E_{2} \subset E_{1} \subset E, \quad \operatorname{codim}\left(E_{i} \subset E_{i-1}\right)=1
$$

can be found directly, whence $E$ can be solved directly by quadratures.
Proof. Given a filtration of sub- $\mathcal{D}$-modules as above, Theorem 5.2.3 explains how to stepwise solve $E$ by quadratures, starting with the first order equation $E_{n}$. To find the filtration, start with an arbitrary eigenvalue $\lambda_{1}$ of the dual
representation, find $E_{\lambda_{1}}^{*}=\left\{\alpha \in E^{*} \mid \rho^{*}(g)(\alpha)=\lambda_{1}(g)(\alpha)\right\}$, and take $E_{1}$ to be the annihilator of $E_{\lambda_{1}}^{*}$ in $E$,

$$
E_{1}=\operatorname{Ann}\left(E_{\lambda_{1}}^{*}\right)=\left\{e \in E \mid \alpha(e)=0, \forall \alpha \in E_{\lambda_{1}}^{*}\right\}
$$

The $\mathcal{D}$-module $E_{1}$ is of codimension 1 in $E$. An arbitrary remaining $\lambda_{i}$ will produce a sub-D-module $E_{\lambda_{i}}^{*}$ with $E_{\lambda_{i}}^{*} \cap E_{\lambda_{1}}^{*}=0$. We may choose $\lambda_{2}$, and take

$$
E_{2}=\operatorname{Ann}\left(E_{\lambda_{2}}^{*}\right) \subset E_{1} \subset E,
$$

which is again necessarily of codimension 1 in $E_{1}$. Repeat for the remaining eigenvalues, and get the whole filtration.

## Chapter 6

## Equations with semisimple symmetry algebras

Semisimple algebras are popular in representation theory, as there is a general theory on how to decompose representations of semisimple Lie algebras into irreducible representations, and up to isomorphisms more or less everything is known about irreducible representations for the classical (semi)simple Lie algebras.

By using Theorem 2.2.3 we are now ready to transfer results on representations of Lie algebras into vector spaces, to $\mathcal{D}$-modules and ODEs.

We find that for a number of equations with semisimple symmetry algebras we obtain solvability by algebraic methods. An algorithm to decompose and solve equations is provided.

### 6.1 General results for semisimple symmetry algebras

Some results from representation theory of Lie algebras into vector spaces depend on having an algebraically closed base field, as seen in the chapter on solvable algebras. This problem occurs whenever we encounter eigenvalue calculations; to be able to say something about roots of characteristic polynomials in general, we need algebraic closure of the coefficient field. And this is certainly a crucial part of studying representations of semisimple algebras,
where calculating roots and weights is more or less the whole trick. Thus, we may, as in Chapter 5, choose to work with modules over complex valued functions, $\mathcal{A}_{C}=C^{\infty}(\mathbb{R}, \mathbb{C})=\mathbb{C} \otimes_{\mathbb{R}} \mathcal{A}$, with $\delta_{A_{C}}$ being the usual derivative in real variable $x$.

Let $\mathfrak{g}$ be a semisimple algebra (over $\mathbb{C}$ whenever algebraic closure of the base field is necessary), and $R(\mathfrak{g})$ the associated Grothendieck ring of isomorphism classes of finite dimensional vector space representations of $\mathfrak{g}$.

Definition 6.1.1. Denote by $\mathcal{D}(\mathfrak{g})$ the ring of isomorphism classes of $\mathcal{D}$ modules with a semisimple symmetry algebra $\mathfrak{g}$. We shall refer to $\mathcal{D}(\mathfrak{g})$ as the symmetry ring of $\mathfrak{g}$.

Let $\omega_{1}, \ldots, \omega_{n}$ be a set of fundamental weights for $\mathfrak{g}$, and let $\Gamma_{1}, \ldots, \Gamma_{n}$ denote the corresponding isomorphism classes in $R(\mathfrak{g})$ with highest weights $\omega_{1}, \ldots, \omega_{n}$. Recall the following result from the theory of representations of semisimple Lie-algebras, see e.g. [4].

Theorem 6.1.1. The representation ring $R(\mathfrak{g})$ is a polynomial ring in the variables $\Gamma_{1}, \ldots, \Gamma_{n}$.

Combining Theorems 6.1 .1 and 2.2.3 we immediately get the following.
Theorem 6.1.2. For a semisimple Lie algebra $\mathfrak{g}$ the symmetry ring $\mathcal{D}(\mathfrak{g})$ is a polynomial ring in classes of $\mathcal{D}$-modules $E_{1}, \ldots, E_{n}$ such that for each $i=1, \ldots n E_{i}^{\#}$ is isomorphic to $\Gamma_{i}$ as a $\mathfrak{g}$-space.

We will call $\mathcal{D}$-modules $E_{1}, \ldots, E_{n}$ generators of $\mathcal{D}(\mathfrak{g})$.
Corollary 6.1.1. Let $(E, \delta)$ be a $\mathcal{D}$-module with a semisimple algebra of symmetries $\mathfrak{g}$. Then $E$ is isomorphic, as a $\mathfrak{g}$-module, to a polynomial $P\left(E_{i}\right)$ in generators $E_{1}, \ldots, E_{n}$ of $\mathcal{D}(\mathfrak{g})$.

An equation that correponds to an irreducible representation $\Gamma_{i}$, associated to a highest weight $\omega_{i}$ of $\mathfrak{g}$ as above we will call a model equation for this symmetry algebra.

## $6.2 \quad \mathfrak{s l}_{2}$ equations

Representations of the Lie algebra $\mathfrak{s l}_{2}(\mathbb{R})$ has a special place in the view of symmetric powers of second order equations, given that all irreducible representations of $\mathfrak{s l}_{2}(\mathbb{R})$ are isomorphic to symmetric powers of the standard two
dimensional representation.

### 6.2.1 Schrödinger equations.

Equations on the following type we will denote as being of Schrödinger type, with potential $W(x)$.

$$
\begin{equation*}
y^{\prime \prime}+W(x) y=0 \tag{6.1}
\end{equation*}
$$

It is self adjoint, hence the corresponding module $E$ is isomorphic to $E^{*}$. Conditions for

$$
\Delta=a_{1}+a_{2} \partial
$$

to be a symmetry operator of (6.1) are by direct calculation found to be that $a_{2}$ solves

$$
\begin{equation*}
z^{\prime \prime \prime}+4 W z^{\prime}+2 W^{\prime} z=0 \tag{6.2}
\end{equation*}
$$

and that $a_{1}=\frac{-a_{2}^{\prime}}{2}$. Thus symmetries are given by generating functions that solve (6.2), i.e.

$$
\Delta_{a}=-\frac{a^{\prime}}{2}+a \partial, \quad a \in \operatorname{Sol}(6.2)
$$

In Section 4.3 we studied the symmetry equation of general second order equations in detail. Recall that the non-trivial symmetries of an equation $L y=0$ were generated by solutions of the equation $S^{2}\left(L_{0}\right)$ where

$$
L_{0}=\partial^{2}+\left(\frac{1}{2} a_{2}^{\prime}-\frac{1}{4} a_{2}^{2}-a_{1}\right)
$$

For a Schrödinger equation (6.1) $a_{2}=0, a_{1}=-W(x)$, so

$$
\begin{equation*}
L_{0}=\partial^{2}-a_{1}=\partial^{2}+W=L \tag{6.3}
\end{equation*}
$$

Theorem 6.2.1. For a Schrödinger equation

$$
\begin{equation*}
L(y)=y^{\prime \prime}+W y=0 \tag{6.4}
\end{equation*}
$$

the symmetry equation is its second symmetric power, i. e.

$$
\operatorname{Sym}(L)=\left\{\left.\Delta_{a}=-\frac{a^{\prime}}{2}+a \partial \right\rvert\, a \text { solves } S^{2}(L)(z)=0\right\}
$$

where

$$
\begin{equation*}
S^{2}(L)=\partial^{3}+4 W \partial+2 W^{\prime} \tag{6.5}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
\operatorname{Sym}(L) \cong \mathfrak{s l}_{2} \tag{6.6}
\end{equation*}
$$

as Lie algebras, where the Lie bracket operation in $\operatorname{Sym}(L)$ is the commutator of symmetry operators.

Proof. The calculations in Section 4.3 and discussions above prove that $\operatorname{Sym}(L) \cong S^{2}(E)^{\#} \cong \operatorname{Sol}(6.2) . \operatorname{Sym}(L)$ is a Lie algebra with respect to commutators of operators, as discussed in Section 4.2, on symmetry operators. To prove that it is isomorphic to $\mathfrak{s l}_{2}(\mathbb{R})$ one may calculate symbolically with a set of fundamental solutions of (6.2), using differential relations, and get the desired result. However, the proposition that follows enables us to prove this in yet another way.

Proposition 6.2.1. The commutator of symmetry operators in Theorem 6.2.1 corresponds to the following Lie bracket $\langle\cdot, \cdot\rangle$ on the solution space of the symmetry equation (6.2), given by

$$
\begin{equation*}
\left[\Delta_{a}, \Delta_{b}\right]=\Delta_{\langle a, b\rangle} \tag{6.7}
\end{equation*}
$$

which yields simply

$$
\langle a, b\rangle=a b^{\prime}-a^{\prime} b
$$

for $a, b \in \operatorname{Sol}(6.2)$. Denote the equivalent Lie bracket on $S^{2}(E)^{\#}$

$$
\begin{equation*}
[\cdot, \cdot]: S^{2}(E)^{\#} \times S^{2}(E)^{\#} \rightarrow S^{2}(E)^{\#} \tag{6.8}
\end{equation*}
$$

where

$$
\begin{equation*}
\left[\theta_{a}, \theta_{b}\right]=\theta_{\langle a, b\rangle} \tag{6.9}
\end{equation*}
$$

for elements $\theta_{a}, \theta_{b} \in S^{2}(E)^{\#}$ generated by solutions $a, b$ of (6.2). The solution space of (6.2) is isomorphic to $\mathfrak{s l}_{2}(\mathbb{R})$ with respect to this bracket.

Proof. Calculating the commutator $\left[\Delta_{a}, \Delta_{b}\right]$ directly gives precisely the formula $\langle a, b\rangle=a b^{\prime}-a^{\prime} b$. Theorem 2.3.1 asserts that solutions $a, b$ of the $S^{2}(E)$-equation are linear combinations of solutions $u^{2}, u v, v^{2}$ where $u, v$ are linear independent solutions of the Schrödinger equation. Whence, we may calculate all brackets of elements from a basis $\left\{u^{2}, u v, v^{2}\right\}$ of the solution space of the symmetry equation (6.2), and will find that as a Lie algebra it is isomorphic to $\mathfrak{s l}_{2}(\mathbb{R})$.

Corollary 6.2.1. The Lie bracket [., •] on $S^{2}(E)^{\#}$ in Proposition 6.2.1 extends by $\mathcal{A}$-linearity to a bracket

$$
\begin{equation*}
[\cdot, \cdot]: S^{2}(E) \times S^{2}(E) \rightarrow S^{2}(E) \tag{6.10}
\end{equation*}
$$

with respect to which $\left(S^{2}(E), \delta\right)$ is a $\mathcal{D}$-Lie-algebra.
Proof. Theorem 2.2.3 states that $S^{2}(E)^{\#}$ spans $S^{2}(E)$ over $\mathcal{A}$, thus the bracket extends in a well-defined way by $\mathcal{A}$-linearity to $S^{2}(E)$, with the Liebracket properties intact. We need only check that $[\cdot, \cdot]$ commutes with $\delta$ in accordance with Definition 2.2.3. Let $X=f_{1} \theta_{1}, Y=f_{2} \theta_{2}$, with $f_{i} \in \mathcal{A}$, $\theta_{i} \in S^{2}(E)^{\#}$. Then

$$
\begin{aligned}
\delta[X, Y] & =\delta\left[f_{1} \theta_{1}, f_{2} \theta_{2}\right]=\delta\left(f_{1} f_{2}\left[\theta_{1}, \theta_{2}\right]\right) \\
& =\left(f_{1} f_{2}\right)^{\prime}\left[\theta_{1}, \theta_{2}\right]+f_{1} f_{2} \delta\left[\theta_{1}, \theta_{2}\right] \\
& =\left(f_{1}^{\prime} f_{2}+f_{1} f_{2}^{\prime}\right)\left[\theta_{1}, \theta_{2}\right]=\left[f_{1}^{\prime} \theta_{1}, f_{2} \theta_{2}\right]+\left[f_{1} \theta_{1}, f_{2}^{\prime} \theta_{2}\right] \\
& =\left[\delta\left(f_{1} \theta_{1}\right), f_{2} \theta_{2}\right]+\left[f_{1} \theta_{1}, \delta\left(f_{2} \theta_{2}\right)\right] \\
& =[\delta X, Y]+[X, \delta Y]
\end{aligned}
$$

whence $S^{2}(E)$ is a $\mathcal{D}$-Lie algebra.

### 6.2.2 Symmetric powers of a Schrödinger equation.

We shall see that the Schrödinger equations have special properties. From our basic equation we can derive a whole hierarchy of new equations $S^{k}(E)$. Throughout this section we will work with $S^{k}\left(E^{*}\right)$, choosing to work with the dual module $E^{*}$ merely simplifies calculations, and generates exactly the same equations as the module $E$.

In Section 2.3 symmetric powers of second order equations were discussed in some detail. Let $(E, \delta)$ be the $\mathcal{D}$-module corresponding to

$$
\begin{equation*}
y^{\prime \prime}+W y=0 \tag{6.11}
\end{equation*}
$$

with basis $\left\{e_{1}^{*}, e_{2}^{*}\right\}$ dual to the primitive element basis of $E$. As before, denote the induced basis of $S^{k}\left(E^{*}\right)$

$$
\begin{equation*}
\left\{\alpha_{l}=\left(e_{1}^{*}\right)^{k-l+1} \cdot\left(e_{2}^{*}\right)^{l-1}, l=1, \ldots, k+1\right\} \tag{6.12}
\end{equation*}
$$

For Schrödinger equations with $a_{1}=-W(x), a_{2}=0$ we get that

$$
\begin{equation*}
\delta: \alpha_{l} \mapsto(k-l+1) W \alpha_{l+1}-(l-1) \alpha_{l-1} \tag{6.13}
\end{equation*}
$$

for $l=1, \ldots, k+1$. For a general element

$$
\begin{equation*}
\theta=\sum_{l=1}^{k+1} g_{l}(x) \alpha_{l} \tag{6.14}
\end{equation*}
$$

in $S^{k}\left(E^{*}\right)$, the requirement $\delta \theta=0$ results in the system of $k+1$ equations

$$
\begin{equation*}
g_{s}^{\prime}+(k-s+2) W g_{s-1}-s g_{s+1}=0 \tag{6.15}
\end{equation*}
$$

$s=1, \ldots, k+1$. Thus, for Schrödinger equations Proposition 2.3.2 has the following form.

Proposition 6.2.2. For each $k$ the kernel $S^{k}\left(E^{*}\right)^{\#}$ consists of elements

$$
\begin{equation*}
\theta_{y}=y \alpha_{1}+y^{\prime} \alpha_{2}+\sum_{l=3}^{k+1} g_{l}(y) \alpha_{l}, \quad y \in \operatorname{Sol}(k) \tag{6.16}
\end{equation*}
$$

where

$$
\begin{equation*}
g_{l}=\frac{1}{l-1}\left[(k-l+3) W \cdot g_{l-2}+g_{l-1}^{\prime}\right] \tag{6.17}
\end{equation*}
$$

for $l=2, . ., k+1$, where $g_{1}=y$ solves the $S^{k}\left(E^{*}\right)$ equation, i.e. the equation in $y$ we get from setting

$$
\delta \theta_{y}=0
$$

for $\theta_{y}$ on the form (6.16), with $g_{l}-s$ expressed in derivatives of $y$.
Fix $(k)$ below to denote the equation $S^{k}\left(E^{*}\right)$. They are on the form

$$
\begin{array}{ll}
\text { (i) } \quad y^{\prime \prime}+W y & =0 \\
\text { (ii) } y^{\prime \prime \prime}+4 W y^{\prime}+2 W^{\prime} y & =0 \\
\text { (iii) } y^{(4)}+10 W y^{\prime \prime}+10 W^{\prime} y^{\prime}+\left(9 W^{2}+3 W^{\prime \prime}\right) y & =0 \\
\text { (iv) } y^{(5)}+20 W y^{\prime \prime \prime}+30 W^{\prime} y^{\prime \prime}+\left[64 W^{2}+18 W^{\prime \prime}\right] y^{\prime}+\left[64 W W^{\prime}+4 W^{\prime \prime \prime}\right] y & =0
\end{array}
$$

and so on.

Focusing on the hierarchy of symmetric powers of Schrödinger equations we shall now see that the bracket operation on $S^{2}\left(E^{*}\right)$ extends to the whole hierarchy.

The bracket in Corollary 6.2.1 can be obtained in a different way.
Proposition 6.2.3. Let $(E, \delta)$ be the $\mathcal{D}$-module corresponding to the Schrödinger equation (6.11). The equation corresponding to its second exterior power, $\left(\bigwedge^{2}\left(E^{*}\right), \delta\right)$, is

$$
u^{\prime}=0 .
$$

that is, $\Omega=e_{1} \wedge e_{2}$ is $\delta$-invariant. Here $\left\{e_{1}, e_{2}\right\}$ is the standard primitive element basis of $E$. Moreover, $\Omega$ determines $a \delta$-invariant skew-symmetric, $\mathcal{A}$-linear bracket operation on $E^{*}$ defined by

$$
\left[\theta_{1}, \theta_{2}\right]=\left\langle\theta_{1} \wedge \theta_{2}, \Omega\right\rangle
$$

for $\theta_{1}, \theta_{2} \in E^{*}$.
Proof. The skew-symmetric form $\Omega=e_{1} \wedge e_{2} \in \Lambda^{2}(E)$, is in the kernel of $\delta$ :

$$
\delta \Omega=e_{2} \wedge e_{2}+e_{1} \wedge\left(a_{1} e_{1}+a_{2} e_{2}\right)=a_{2} e_{1} \wedge e_{2}=0
$$

since $a_{1}=-W, a_{2}=0$. The bracket is thus $\delta$-invariant and obviously skewsymmetric and $\mathcal{A}$-linear, due to the properties of $\Omega$.

Proposition 6.2.4. Given a $\mathcal{D}$-module $(E, \delta)$ corresponding to an equation of Schrödinger type as in Theorem 6.2.1, there is a unique skew-symmetric bracket

$$
[\cdot, \cdot]: S^{m}\left(E^{*}\right) \times S^{n}\left(E^{*}\right) \rightarrow S^{m+n-2}\left(E^{*}\right)
$$

for all $m, n \geq 1$ which is
(i) $\mathcal{A}$-linear.
(ii) $[f \cdot g, h]=f \cdot[g, h]+g \cdot[f, h] \quad \forall f, g, h \in S \cdot\left(E^{*}\right)$
(iii) For $n=m=1$ the bracket coincides with the bracket in Proposition 6.2.3.

Proof. Given the bracket operation in Proposition 6.2.3 the properties $(i)-$ (ii) determine its extension to symmetric powers $S^{m}\left(E^{*}\right) \times S^{n}\left(E^{*}\right)$.

We immediately observe that $\Omega$ being $\delta$-invariant implies that so is the extended bracket $[\cdot, \cdot]$. Thus, it restricts to kernels of $\delta$-s as follows.

Proposition 6.2.5. The bracket operation in Proposition 6.2.4 restricts to kernels $\left(S^{k}\left(E^{*}\right)\right)^{\#}=S^{k}(V)$, where $V=\left(E^{*}\right)^{\#}$. The bracket

$$
[\cdot, \cdot]: S^{m}(V) \times S^{n}(V) \longrightarrow S^{m+n-2}(V)
$$

has the properties
(i) $\mathbb{R}$-linearity
(ii) $[f \cdot g, h]=f \cdot[g, h]+g \cdot[f, h], \quad \forall f, g, h \in S \cdot(V)$

This is obviously equivalent to a bracket on the solution spaces of the symmetic power equations,

$$
[\cdot, \cdot]: \operatorname{Sol}(m) \times \operatorname{Sol}(n) \rightarrow \operatorname{Sol}(m+n-2)
$$

with

$$
\left[\theta_{y}, \theta_{z}\right]=\theta_{[y, z]}
$$

for solutions $y, z$ of $S^{m}\left(E^{*}\right)$ and $S^{n}\left(E^{*}\right)$ equations repectively. This means that solutions of the $S^{2}\left(E^{*}\right)$ equation (6.2) produce symmetries of all equations $S^{k}\left(E^{*}\right)$, and not only $E^{*}$.

Theorem 6.2.2. Any solution $a \in \operatorname{Sol}\left(S^{2}\left(E^{*}\right)\right)$ produces a symmetry

$$
\mathcal{O}_{\theta_{a}}^{m} \stackrel{\text { def }}{=}\left[\theta_{a}, \cdot\right]: S^{m}\left(E^{*}\right) \longrightarrow S^{m}\left(E^{*}\right)
$$

The corresponding symmetry operator is

$$
\begin{equation*}
\mathcal{O}=\mathcal{O}_{a}^{m}: \operatorname{Sol}(m) \longrightarrow \operatorname{Sol}(m) \tag{6.18}
\end{equation*}
$$

with the correspondence

$$
\mathcal{O}_{\theta_{a}}^{m}\left(\beta_{y}\right)=\beta_{\mathcal{O}_{a}^{m}(y)}
$$

The precise expression is

$$
\begin{equation*}
\mathcal{O}_{a}^{m}=\frac{1}{2}\left(-m a^{\prime}+2 a \partial\right) \tag{6.19}
\end{equation*}
$$

for any $m \geq 1, a \in \operatorname{Sol}\left(S^{2}\left(E^{*}\right)\right)$.
Due to Theorem 2.2.3 we know that two linearly independent solutions

$$
\begin{equation*}
u, v \in \operatorname{Sol}\left(E^{*}\right) \Rightarrow \text { basis }\left\{\theta_{u}, \theta_{v}\right\} \subset V \text { of } E^{*} \text { over } \mathcal{A} \tag{6.20}
\end{equation*}
$$

Hence, for any $k \geq 1$,

$$
\begin{equation*}
\left\{\theta_{u^{k-l} v^{l}}=\theta_{u}^{k-l} \cdot \theta_{v}^{l}\right\} \subset S^{k}(V), \quad 0 \leq l \leq k \tag{6.21}
\end{equation*}
$$

is the basis of $S^{k}\left(E^{*}\right)$ over $\mathcal{A}$ corresponding to the fundamental set of solutions

$$
u^{k-l} v^{l} \in \operatorname{Sol}(k), \quad 0 \leq l \leq k
$$

It is now easy to calculate the action of the symmetries

$$
\begin{equation*}
\mathcal{O}_{1}^{k}=\mathcal{O}_{u^{2}}^{k}, \quad \mathcal{O}_{2}^{k}=\mathcal{O}_{u v}^{k}, \quad \mathcal{O}_{3}^{k}=\mathcal{O}_{v^{2}}^{k} \tag{6.22}
\end{equation*}
$$

on basis elements $\theta_{u^{k-l} v^{l}}$ just in terms of brackets of the generating functions.
Theorem 6.2.3. For any $k \geq 1$ the symmetries of $S^{k}\left(E^{*}\right)$

$$
X_{+}=-\frac{1}{2 c} \mathcal{O}_{3}^{k}, \quad X_{-}=\frac{1}{2 c} \mathcal{O}_{1}^{k} \text { and } H=\frac{1}{c} \mathcal{O}_{2}^{k}
$$

where $c=\langle u, v\rangle \in \mathbb{R}$ constitute a basis of the $\mathfrak{s l}_{2}(\mathbb{R})$-algebra of symmetries $\cong \operatorname{Sol}(i i)$ with commutators

$$
\left[X_{+}, X_{-}\right]=H, \quad\left[H, X_{+}\right]=2 X_{+}, \quad\left[H, X_{-}\right]=-2 X_{-}
$$

Hence, $S^{k}\left(E^{*}\right)$ decomposes into rank 1 sub-D-modules corresponding to different eigenvalues of $H$

$$
S^{k}\left(E^{*}\right)=\left\langle\theta_{u^{k}}\right\rangle_{\mathcal{A}} \oplus\left\langle\theta_{u^{k-1} v}\right\rangle_{\mathcal{A}} \oplus \ldots \oplus\left\langle\theta_{v^{k}}\right\rangle_{\mathcal{A}}
$$

$-k,-k+2, \ldots, k-2, k$

Proof. The commutator relations are calculated directly in terms of the brackets

$$
\left[v^{2}, u^{2}\right]=-4 c u v, \quad\left[u v, v^{2}\right]=2 c v^{2}, \quad\left[u v, u^{2}\right]=-2 c u^{2}
$$

and knowing the form of the operators $\mathcal{O}_{i}^{k}$ from (6.19). Furthermore,

$$
H\left(\theta_{u^{k-l} v^{l}}\right)=(2 l-k) \theta_{u^{k-l} v^{l}}
$$

for $0 \leq l \leq k$. Certainly

$$
\begin{aligned}
& X_{+}: \operatorname{Eig}_{\lambda}(H) \rightarrow \operatorname{Eig}_{\lambda+2}(H) \\
& X_{-}: \operatorname{Eig}_{\lambda}(H) \rightarrow \operatorname{Eig}_{\lambda-2}(H)
\end{aligned}
$$

where $E i g_{\lambda}(H)$ denotes the eigen-submodule of $H$ corresponding to the eigenvalue $\lambda$. It is generated over $\mathcal{A}$ by the eigenspace in $S^{k}(V)$ of $H$.

Theorem 6.2.4. $A \mathcal{D}$-module $(E, \delta)$ with a representation of symmetries

$$
\mathfrak{s l}_{2}(\mathbb{R}) \rightarrow E n d_{\mathcal{A}}(E)
$$

is decomposable into a direct sum of $\mathcal{D}$-modules

$$
E=\bigoplus_{i=1}^{m} E_{i}
$$

where each $E_{i}$ is an irreducible subrepresentation of $E$. Moreover, each $E_{i}$ is isomorphic to $S^{n_{i}}\left(M_{i}\right)$ as a $\mathcal{D}$-module and as an $\mathfrak{s l}_{2}(\mathbb{R})$-module, for a rank 2 $\mathcal{D}$-module $M_{i}$.

Proof. The representation of symmetries into $E$ restricts to a representation of $\mathfrak{s l}_{2}(\mathbb{R})$ in the $\mathbb{R}$-vector space $V=E^{\#} \subset E$, hence it decomposes into a direct sum of representations

$$
V=\oplus_{i=1}^{m} V_{i}
$$

where the $V_{i}$ are subspaces of $V$ such that restricted to $V_{i}$ the representation is irreducible. But any irreducible representation of $\mathfrak{s l}_{2}(\mathbb{R})$ into a vector space of dimension $n_{i}(<\infty)$ is isomorphic to the $\left(n_{i}-1\right)$ th symmetric power of
the standard two dimensional representation, i.e. each $V_{i} \cong S^{n_{i}-1}\left(W_{i}\right)$, where $\operatorname{dim}_{\mathbb{R}}\left(W_{i}\right)=2$ for all $i$. Recall that the symmetries commute with $\delta$ in $E$, hence $E_{i}=\mathcal{A} \cdot V_{i}$, the $\mathcal{A}$-module generated of $V_{i}$ over $\mathcal{A}$. Due to Theorem 2.2.3 the vector space isomorphism $V_{i} \cong S^{n_{i}-1}\left(W_{i}\right)$ lifts to an isomorphism of $\mathcal{D}$-modules $E_{i} \cong S^{n_{i}}\left(M_{i}\right)$, where $S^{n_{i}}\left(M_{i}\right)$ is generated of $S^{n_{i}}\left(W_{i}\right)$ over $\mathcal{A}$.

Corollary 6.2.2. Let $(E, \delta)$ be a $\mathcal{D}$-module with an $\mathfrak{s l}_{2}(\mathbb{R})$ algebra of symmetries as in Theorem 6.2.4. If the irreducible $\mathfrak{s l}_{2}(\mathbb{R})$-modules $E_{i}$ in its decomposition are of distinct ranks, then the equation corresponding to $E$ can be solved by algebraic operations and quadrature.

Corollary 6.2.3. Let $(E, \delta)$ be a $\mathcal{D}$-module with an $\mathfrak{s l}_{2}(\mathbb{R})$ algebra of symmetries as in Theorem 6.2.4. If there are irreducible $\mathfrak{s l}_{2}(\mathbb{R})$-modules $E_{i}$ with ranks of multiplicity $m_{1}, \ldots, m_{k}>1$ in its decomposition, then the obstruction to solve the equation corresponding to $E$ by algebraic operations and quadrature is $k$ first order systems of ODEs, with size $m_{1} \times m_{1}, \ldots, m_{k} \times m_{k}$ respectively.

See Subsection 6.2.3 for a detailed account on how to decompose and solve $\mathfrak{s l}_{2}(\mathbb{R})$-equations.

Example 6.2.1. A Schrödinger equation

$$
y^{\prime \prime}+W(x) y=0
$$

is a model equation of $\mathfrak{s l}_{2}(\mathbb{R})$.

### 6.2.3 Algorithm to solve $\mathfrak{s l}_{2}(\mathbb{R})$-equations.

The calculations preceeding Theorem 6.2.4 tell us how we should approach $\mathfrak{s l}_{2}(\mathbb{R})$-modules in order to find its complete reduction, identify sub- $\mathcal{D}$-modules $M_{i}$ as in Theorem 6.2.4, and eventually solve the original equation.

Theorem 2.3.1 tells us that

$$
\text { Solutions of } M_{i} \Rightarrow \text { Solutions of } S^{l_{i}}\left(M_{i}\right) \Rightarrow \text { Solutions of } E
$$

An outline of the algorithm is as follows.

Step 1 Given an $\mathfrak{s l}_{2}(\mathbb{R})$ - $\mathcal{D}$-module $(E, \delta)$ as in Theorem 6.2.4, find a basis $\left\{X_{+}, X_{-}, H\right\}$ of its $\mathfrak{s l}_{2}(\mathbb{R})$-algebra of symmetries that satisfies the commutator relations in Theorem 6.2.3. Calculate the eigen-sub-modules $E_{\lambda_{i}}$ in $E$ corresponding to weights $\left\{\lambda_{1}, \ldots, \lambda_{k}\right\}$ of the diagonal element $H$. This yields a decomposition

$$
E=\oplus_{i=1}^{k} E_{\lambda_{i}}
$$

such that

$$
\begin{aligned}
& X_{+}: E_{\lambda_{i}} \rightarrow E_{\lambda_{i+2}} \quad \text { and } \\
& X_{-}: E_{\lambda_{i}} \rightarrow E_{\lambda_{i-2}}
\end{aligned}
$$

The rank of $E_{\lambda_{i}}$ over $\mathcal{A}$ is the multiplicity of the weight $\lambda_{i}$, which we denote $m_{i}$. Given a decomposition of $E$ as in 6.2.4, then the values of the weights are precisely integers $\lambda_{j}=2 j-n_{i}, 0 \leq j \leq n_{i}, i=1, \ldots, m$.

Step 2 Identify all $\lambda_{i}$-s of multiplicity $m_{i}=1$. For each weight $\lambda_{i}$ of multiplicity 1 any non-zero $w \in \mathrm{E}_{\lambda_{i}}$ with $X_{+}(w)=0$ generates an irreducible $\mathfrak{s l}_{2}(\mathbb{R})$-module in $E$

$$
\begin{equation*}
\left\langle X_{-}^{n_{i}}(w)\right\rangle \oplus \ldots \oplus\left\langle X_{-}(w)\right\rangle \oplus\langle w\rangle \tag{6.23}
\end{equation*}
$$

where $n_{i}$ is the smallest integer such that $X_{-}^{n_{i}}(w)=0$.
Moreover, this is a sub- $\mathcal{D}$-module in $E$, and it is isomorphic to $S^{n_{i}}\left(M_{i}\right)$, for a rank 2 "Schrödinger" module $M_{i}$ as a $\mathcal{D}$-module and a $\mathfrak{s l}_{2}(\mathbb{R})$-module.

Recall from Theorem 6.2.3 the structure of symmetric powers of a Schrödinger equations:

$$
S^{n_{i}}\left(M_{i}\right)=\left\langle\theta_{u^{n_{i}}}\right\rangle_{\mathcal{A}} \oplus\left\langle\theta_{u^{n_{i}-1} v}\right\rangle_{\mathcal{A}} \oplus \ldots \oplus\left\langle\theta_{v^{n_{i}}}\right\rangle_{\mathcal{A}}
$$

for a fundamental set of solutions $u, v$ of the Schrödinger equation corresponding to $M_{i}$.
To identify $M_{i}$ take the "fraction" $w / X_{-}(w)$. Its last coefficient is a fraction

$$
\begin{equation*}
\alpha=\frac{v}{u} \tag{6.24}
\end{equation*}
$$

of fundamental solutions of the $M_{i}$-equation. Denote the potential of that equation by $W(x)$. Differentiating and using the differential relations

$$
\begin{equation*}
u^{\prime \prime}+W(x) u=0, \quad v^{\prime \prime}+W(x) v=0 \tag{6.25}
\end{equation*}
$$

yields the following expression for $W$

$$
\begin{equation*}
W=\gamma^{2}+\gamma^{\prime} \tag{6.26}
\end{equation*}
$$

where

$$
\begin{equation*}
\gamma=\frac{\ln (\alpha)^{\prime \prime}-\left(\ln \alpha^{\prime}\right)^{2}}{2 \ln \alpha^{\prime}} \tag{6.27}
\end{equation*}
$$

For $w=f(x) \theta_{v^{k}}, \delta w=f^{\prime}(x) w$, so from $\delta w / w$ we get $\eta=(\ln f)^{\prime}$. Integrating, we get

$$
\begin{equation*}
f=e^{\int \eta d x} \tag{6.28}
\end{equation*}
$$

and the last coefficient of $w / f$ is $v^{k}$, from which we deduce $v$. Then, $u=v / \alpha$.
Step 3 "Remove" the irreducible $\mathfrak{s l}_{2}(\mathbb{R})$-modules in Step 2 from the module $E$, i.e. work in their complement in $E$. For each weight $\lambda$ denote the complement in $E_{\lambda}$ of these sub-modules by $\widetilde{E}_{\lambda}$. If there are weights of "remaining" multiplicity 1, i.e. $\operatorname{rank}\left(\widetilde{E}_{\lambda}\right)=1$, repeat Step 2 for those weights.

Identify the weight with highest integer value, denote it $\lambda$, and its remaining multiplicity $m>1$. We still have that a non-zero $w \in \widetilde{E}_{\lambda}$ with $X_{+}(w)=0$ generates an irreducible $\mathfrak{s l}_{2}(\mathbb{R})$-module in $E$, but we are no longer guaranteed that this is also a sub- $\mathcal{D}$-module. If $w$ in addition satisfies the condition

$$
\begin{equation*}
\delta w=f(x) w \tag{6.29}
\end{equation*}
$$

then it generates a sub- $\mathcal{D}$-module in $E$ as in Step 2 .
Given a basis $\left\{w_{1}, \ldots, w_{m}\right\}$ of $\widetilde{E}_{\lambda}$, applying $\delta$ yields

$$
\begin{equation*}
\underline{w}=A_{i} \underline{w} \tag{6.30}
\end{equation*}
$$

for some matrix $A_{i}$ with coefficients in $\mathcal{A}$. Solving the corresponding first order system

$$
\begin{equation*}
\underline{h}^{\prime}+A_{i}^{T} \underline{h}=0 \tag{6.31}
\end{equation*}
$$

is the obstruction to identify the sub- $\mathcal{D}$-modules that are irreducible $\mathfrak{s l}_{2}(\mathbb{R})$ modules, as obtained in Step 2 above.

Repeat for highest value weights succesively to get the decomposition of $E$.

### 6.2.4 Schrödinger equations with shared symmetries.

Returning to a base equation with potential $W(x)$

$$
\begin{equation*}
y^{\prime \prime}+W(x) y=0 \tag{6.32}
\end{equation*}
$$

we may discuss, as is done in [12], which potentials share a symmetry $\Delta_{z}=$ $-\frac{z^{\prime}}{2}+z \partial, z$ being a solution of the symmetry equation (6.2). Let $W_{0}(x)$ be a potential of an equation with symmetry $\Delta_{z}$, fixed $z$. Then $W_{o}$ is a particular solution of the symmetry equation viewed as a first order equation for $W(x)$ :

$$
\begin{equation*}
z^{\prime \prime \prime}+4 W z^{\prime}+2 W^{\prime} z=0 \tag{6.33}
\end{equation*}
$$

Integrating the separable homogeneous equation yields that equations with potentials on the form

$$
\begin{equation*}
W(x)=W_{0}(x)+\frac{c}{z^{2}}, \quad c \in \mathbb{R} \tag{6.34}
\end{equation*}
$$

share the symmetry $\Delta_{z}$. Recall also that a fundamental set of solutions $\{u, v\}$ of the base equation (6.32) with potential $W_{0}(x)$ generate a fundamental set of solutions $\left\{u^{2}, u v, v^{2}\right\}$ of the symmetry equation (6.33), hence all equations with potentials

$$
\begin{equation*}
W(x)=W_{0}(x)+\frac{c}{\left(c_{1} u^{2}+c_{2} u v+c_{3} v^{2}\right)^{2}} \tag{6.35}
\end{equation*}
$$

where $c, c_{i} \in \mathbb{R}$ such that the denominator in the fraction is non-zero, are integrable (by quadratures).

Example 6.2.2. Let $W_{0}=0,\left\{u=a_{1}+a_{2} x, v=b_{1}+b_{2} x\right\}$. Then equations with potentials

$$
W=\frac{c}{\left(c_{1} u^{2}+c_{2} u v+c_{3} v^{2}\right)^{2}}
$$

are integrable in quadratures, symmetries $\Delta_{u^{2}}, \Delta_{u v}$ and $\Delta_{v^{2}}$.
Example 6.2.3. Let $W_{0}=\omega^{2}$, with $\{u=\cos (\omega x), v=\sin (\omega x)\}$. Then equations with potentials

$$
W=\omega^{2}+\frac{c}{\left(c_{1} \cos ^{2}(\omega x)+c_{2} \cos (\omega x) \sin (\omega x)+c_{3} \sin ^{2}(\omega x)\right)^{2}}
$$

are integrable in quadratures, symmetries

$$
\Delta_{\cos ^{2}(\omega x)}, \Delta_{\cos (\omega x) \sin (\omega x)} \text { and } \Delta_{\sin ^{2}(\omega x)}
$$

or, equivalently

$$
\Delta_{1}, \Delta_{\cos (2 \omega x)} \text { and } \Delta_{\sin (2 \omega x)}
$$

This is a way to generate new integrable base Schrödinger equations from simpler ones, with shared symmetries, which we may in turn take symmetric products and direct sums of and arrive at new solvable higher order equations.

### 6.3 Model equations for $\mathfrak{s l}_{3}$

Recall that Schrödinger equations are precisely the equations of order two that have the standard $\delta$-invariant "volume form" $\Omega=e_{1} \wedge e_{2} \in \wedge^{2}(E)^{\#}$, which in turn corresponds nicely to the fact that $\mathfrak{s l}_{2}(\mathbb{R})$ is connected to the preservation of a volume form on a two dimensional space.

We may expect that the geometric properties of the classic Lie algebras are reflected in the associated model equations. A search for model equations for $\mathfrak{s l}_{3}$ should thus point us towards third order equations with a $\delta$-invariant standard volume form.

A third order equation

$$
y^{\prime \prime \prime}+f_{3}(x) y^{\prime \prime}+f_{2}(x) y^{\prime}+f_{1}(x) y=0
$$

with associated $\mathcal{D}$-module $(E, \delta)$ has invariant volume forms on the form $g(x) \Omega \in \wedge^{3}\left(E^{*}\right), \Omega=e_{1}^{*} \wedge e_{2}^{*} \wedge e_{3}^{*}$, for any $g(x) \in \mathcal{A}$ that solves

$$
g^{\prime}-f_{1}(x) g=0
$$

Hence any third order equation on the form

$$
y^{\prime \prime \prime}+f(x) y^{\prime}+g(x) y=0
$$

has an invariant standard volume form $\Omega=e_{1}^{*} \wedge e_{2}^{*} \wedge e_{3}^{*} \in \wedge^{3}\left(E^{*}\right)$.

A sub-example of the above third order equations is a general skew-adjoint equation

$$
\begin{equation*}
y^{\prime \prime \prime}-f(x) y^{\prime}-\frac{1}{2} f^{\prime}(x) y=0 \tag{6.36}
\end{equation*}
$$

For $L=\partial^{3}-f(x) \partial-\frac{1}{2} f^{\prime}(x) \operatorname{Sym}(L)$ decomposes into

$$
\operatorname{Sym}(L)=\operatorname{Sym}_{0} \oplus \operatorname{Sym}_{1}
$$

thus we may split calculation of symmetries to consider $S y m_{0}$ and $S y m_{1}$ separately. The symmetry equations separate into

$$
\begin{gather*}
p^{\prime \prime \prime}-f p-\frac{f^{\prime}}{2} p=0  \tag{6.37}\\
s^{(5)}-5 f s^{\prime \prime \prime}-\frac{15}{2} f^{\prime} s^{\prime \prime}+\left(4 f^{2}-\frac{9}{2} f^{\prime \prime}\right) s^{\prime}+\left(4 f f^{\prime}-f^{\prime \prime \prime}\right) s=0 \tag{6.38}
\end{gather*}
$$

in the sense that

$$
\begin{align*}
& \operatorname{Sym}_{0}(L)=\left\{\Delta_{p}^{0}=-p^{\prime}+p \partial \mid p \text { solves (6.37) }\right\}, \quad \text { and }  \tag{6.39}\\
& \operatorname{Sym}_{1}(L)=\left\{\left.\Delta_{s}^{1}=\left(\frac{s^{\prime \prime}}{6}-\frac{2 f s}{3}\right)-\frac{s^{\prime}}{2} \partial+s \partial^{2} \right\rvert\, s \text { solves }(6.38)\right\} \tag{6.40}
\end{align*}
$$

We recall the graded structure of $\operatorname{Sym}(L)$ from Section 4.2, and note that the commutators

$$
\begin{align*}
{\left[\Delta_{p}^{0}, \Delta_{q}^{0}\right] } & =\Delta_{\{p, q\}_{00}}^{0}  \tag{6.41}\\
{\left[\Delta_{s}^{1}, \Delta_{p}^{0}\right] } & =\Delta_{\{s, p\}_{01}}^{1}  \tag{6.42}\\
{\left[\Delta_{s}^{1}, \Delta_{w}^{1}\right] } & =\Delta_{\{s, w\}_{11}}^{0} \tag{6.43}
\end{align*}
$$

induce the following structure on the solution spaces of the symmetry equations:

$$
\begin{align*}
\{p, q\}_{00} & =p q^{\prime}-p^{\prime} q  \tag{6.44}\\
\{s, p\}_{01} & =2 s p^{\prime}-s^{\prime} p  \tag{6.45}\\
\{s, w\}_{11} & =\frac{1}{6}\left(s^{\prime \prime \prime} w-w^{\prime \prime \prime} s\right)+\frac{1}{4}\left(s^{\prime} w^{\prime \prime}-s^{\prime \prime} w^{\prime}\right)+\frac{2}{3} f\left(s w^{\prime}-s^{\prime} w\right) \tag{6.46}
\end{align*}
$$

where

$$
\{\cdot, \cdot\}_{a b}: S o l_{a} \times \operatorname{Sol}_{b} \rightarrow \operatorname{Sol}_{a+b}
$$

The $S y m_{0}$-equation (6.37) is equal to the original equation $L y=0$, which is again the second symmetric power of the Schrödinger equation (6.47)

$$
\begin{equation*}
B(y)=y^{\prime \prime}-\frac{1}{4} f y=0 \tag{6.47}
\end{equation*}
$$

whereas the Sym $_{1}$-equation (6.38) is both the second symmetric power of (6.37) and the fourth symmetric power of (6.47). This observation makes calculations of the symmetry algebra easier, as we may do them in terms of powers of solutions to the basic Schrödinger equation (6.47).

For any set $\{u, v\}$ of independent solutions of (6.47)

$$
\left\{u^{2}, u v, v^{2}\right\} \quad \text { and } \quad\left\{u^{4}, u^{3} v, u^{2} v^{2}, u v^{3}, v^{4}\right\}
$$

are independent solutions of (6.37) and (6.38) respectively, and generate a full basis of $\operatorname{Sym}(L)$. Considering the commutators of functions asserts that

$$
\operatorname{Sym}(L) \cong \mathfrak{s l}_{3}
$$

and that the subalgebra $\operatorname{Sym}_{0}(L) \subset \operatorname{Sym}(L)$ is isomorphic to $\mathfrak{s l}_{2}(\mathbb{R})$. Its action on $S y m_{1}$ is precisely as described in section 6.2 on $\mathfrak{s l}_{2}(\mathbb{R})$-equations, of $S^{2}(B)$ into $S^{4}(B)$.

For $\mathfrak{g}=\mathfrak{s l}_{3}$ the representation ring $R(\mathfrak{g})$ is generated by $V$ and $\bigwedge^{2}(V) \cong V^{*}$, where $V=F^{3}, F=\mathbb{R}$ or $\mathbb{C}$ denotes the standard (matrix) representation.

Theorem 6.3.1. Equation (6.36) with corresponding $\mathcal{D}$-module $(E, \delta)$ is a model equation for the standard representation of $\mathfrak{s l}_{3}$. $\left(E^{*}, \delta\right) \cong\left(\bigwedge^{2}(E), \delta\right)$ also correspond to equation (6.36).

### 6.4 Invariant volume forms

Let the $\mathcal{D}$-module $(E, \delta)$ correspond to an $n$-th order equation

$$
y^{n}+\left(a_{n} y\right)^{(n-1)}-\left(a_{n-1} y\right)^{(n-2)}+\ldots+(-1)^{n-1} a_{1} y=0
$$

and let $\left\{e_{1}=e, \ldots, e_{n}=\delta^{n-1} e\right\}$ be the standard corresponding primitive element basis. Any element in $\bigwedge^{n}\left(E^{*}\right)$ is a functional multiple of

$$
\begin{equation*}
\Omega=e_{1}^{*} \wedge \ldots \wedge e_{n}^{*} \tag{6.48}
\end{equation*}
$$

By direct calculation we get

$$
\begin{align*}
\delta \Omega & =\sum_{i=1}^{n} e_{1}^{*} \wedge \ldots \wedge \delta\left(e_{i}^{*}\right) \wedge \ldots \wedge e_{n}^{*} \\
& =\left(-a_{1} e_{n}^{*}\right) \wedge e_{2}^{*} \wedge \ldots \wedge e_{n}^{*}  \tag{6.49}\\
& +\sum_{i=2}^{n} e_{1}^{*} \wedge \ldots \wedge e_{i-1}^{*} \wedge\left(e_{i-1}^{*}-a_{i} e_{n}^{*}\right) \wedge e_{i+1}^{*} \wedge \ldots \wedge e_{n}^{*} \\
& =-a_{n} \Omega
\end{align*}
$$

The equation corresponding to the $n$th exterior power of $E^{*},\left(\bigwedge^{n}\left(E^{*}\right), \delta\right)$, is

$$
\begin{equation*}
u^{\prime}-a_{n} u=0 \tag{6.50}
\end{equation*}
$$

Thus, any equation with zero highest coefficient,

$$
y^{n}-\left(a_{n-1} y\right)^{(n-2)}+\ldots+(-1)^{n-1} a_{1} y=0
$$

has an inherent $\delta$-invariant $n$-form $\Omega=e_{1}^{*} \wedge \ldots \wedge e_{n}^{*}$.
As comparison, in $\left(\bigwedge^{n}(E), \delta\right)$, applying $\delta$ to the $n$-form

$$
\Omega=e_{1} \wedge \ldots \wedge e_{n}
$$

yields

$$
\delta \Omega=a_{n} \Omega
$$

corresponding to the dual equation of (6.50),

$$
w^{\prime}+a_{n} w=0 .
$$

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