# 2-descent for Second Order Linear Differential Equations 

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#### Abstract

Let $L$ be a second order linear ordinary differential equation with coefficients in $\mathbb{C}(x)$. The goal in this paper is to reduce $L$ to an equation that is easier to solve. The starting point is an irreducible $L$, of order two, and the goal is to decide if $L$ is projectively equivalent to another equation $\tilde{L}$ that is defined over a subfield $\mathbb{C}(f)$ of $\mathbb{C}(x)$.

This paper treats the case of 2 -descent, which means reduction to a subfield with index $[\mathbb{C}(x): \mathbb{C}(f)]=2$. Although the mathematics has already been treated in other papers, a complete implementation could not be given because it involved a step for which we do not have a complete implementation. The contribution of this paper is to give an approach that is fully implementable [5]. Examples illustrate that this algorithm is very useful for finding closed form solutions (2-descent, if it exists, reduces the number of true singularities from $n$ to at most $n / 2+2$ ).


## Categories and Subject Descriptors

G. 4 [Mathematical Software]: Algorithm design and analysis; I.1.2 [Symbolic and Algebraic Manipulation]: Al-gorithms-Algebraic algorithms

## General Terms

Algorithms

## 1. INTRODUCTION

Let $L=\sum_{i=0}^{n} a_{i} \partial^{i}$ be a differential operator with coefficients in a differential field $K=\mathbb{C}(x)$, where $\partial$ is the usual differentiation $\frac{d}{d x}$. The corresponding differential equation is $L(y)=0$, i.e. $a_{n} y^{(n)}+\cdots+a_{1} y^{\prime}+a_{0} y=0$. The problem of finding closed form solutions of $L$ becomes easier if we can factor $L$ as a product of lower order operators [2, 7, 1] or apply some other approach to reduce the order [9, 14].

A different type of reduction is called descent. Here, the goal is to reduce $L$ to an operator $\tilde{L}$ of the same order, but

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this time defined over a proper subfield $k=\mathbb{C}(f)$ of $K$. Here $\tilde{L}$ must be projectively equivalent to $L$. Informally, this means that $L$ can be solved in terms of the solutions of $\tilde{L}$ and vice versa (a precise definition will be given in Section 2.2).

In this paper, we treat the case of 2 -descent, meaning that $k$ is a subfield of $K$ with index 2 . For now, we treat only second order equations. After applying Kovacic' algorithm, we can assume that $L$ is irreducible (i.e. not a product of lower order factors), and that it has no Liouvillian solutions.

Descent reduces the number of true singularities (Definition 5) from $n$ to $\leq n / 2+2$, which helps to solve differential equations as illustrated in Section 7. In particular, if the number of true singularities ${ }^{1}$ drops to 3 , and if these are regular singularities ${ }^{2}$, then a ${ }_{2} F_{1}$-type solution can be obtained quickly. We can also stop reducing when we reach an operator with four true singularities, because 4 -singularity equations with ${ }_{2} F_{1}$-type solutions are currently being classified [6] by van Hoeij and Vidunas. Classifying equations with closed form solutions and $>4$ singularities would be hard to do, this is where 2-descent becomes crucial.

If $L \in \mathbb{C}(x)[\partial]$ then there is a finitely generated extension $\mathbb{Q} \subseteq C$ with $L \in C(x)[\partial]$, just take $C$ to be the extension of $\mathbb{Q}$ given by the coefficients of $L$. The main design goal for our algorithm is to introduce as few algebraic extensions of $C$ as possible. Without this design goal, Sections 3 and 5 would have been much shorter (if we simply compute the splitting field of the singularities then for Section 5 we can follow [3] and Section 3 becomes trivial. Sections 3 and 5 become non-trivial when we aim to minimize field extensions).

The main results in this paper are in Section 4. We know from [11] that if there is a gauge transformation $G$ from $L$ to $\sigma(L)$, then $L$ will allow descent with respect to $\sigma$. The question is, given $G$, how to find the descent? Is it necessary (as in the terminology in [11]) to trivialize a 2-cocycle, or to perform some equivalent complicated operation such as finding a point on a conic over $C(x)$ ? The answer is no; we give a short and efficient algorithm in Section 4, and we even show (Theorem 1) that it produces a result over an optimal extension of $C$.

### 1.1 Relation to prior work

It is shown in $[3,11]$ that the problem of computing 2 descent can be reduced to another problem (trivializing a 2-cocycle) although no step by step algorithm is given in these papers. The paper [9] does give an algorithm, and im-

[^1]plementation, that can be used to find 2-descent, as follows. If $\sigma$ is a Möbius transformation of order 2 , and $\mathbb{C}(f)$ is the fixed field of $\sigma$, and if $L$ is projectively equivalent to $\sigma(L)$, then we can compute the so-called symmetric product of $L, \sigma(L)$, then apply factorization (DFactorLCLM in Maple), take the 3 'rd order factor found that way, and run the algorithm from [9] to find a second order operator. All of these steps are implemented, and the end result is a 2 -descent.

The problem with the above methods is that they rely on an algorithm that can find a point on a conic defined over $K$ (or an algorithm that solves an equivalent problem). Although such a point must exist when $K=\mathbb{C}(x)$, the proof does not show how to find such a point over a field of constants that is optimal or close to optimal (recall that we wish to minimize the extension of $C$ that the algorithm introduces, where $C \subset \mathbb{C}$ ). There is only an implementation [10] for this step if $C$ is $\mathbb{Q}$ or a transcendental extension of $\mathbb{Q}$. If $L$ contains algebraic numbers, then there is no implementation for finding a point on a conic, and without that, it is not clear how to obtain from $[11,9,3]$ a complete implementation for finding 2-descent.

In this paper we describe a step by step algorithm for finding 2 -descent. The algorithm can be fully implemented [5] because it does not call a conic algorithm. Note: If $L \in$ $C(x)[\partial]$ with $C \subset \mathbb{C}$, and if one allows unnecessary algebraic extensions of $C$ (potentially exponentially large), then it is not hard to implement a conic algorithm, in which case one can consider 2-descent an already solved problem. But in practice our algorithm would be much preferable because it only extends $C$ when necessary (i.e. when there is no 2 descent defined over $C$ ).

## 2. PRELIMINARIES

### 2.1 Differential Operators and Singularities

Let $K=\mathbb{C}(x)$ denote the differential field and let $\mathcal{D}=K[\partial]$ be the ring of differential operators with coefficients in the differential field $K$. Here $\partial$ denotes the usual differentiation $\frac{d}{d x}$. Then elements $L \in \mathcal{D}$ are of the form $L=a_{n} \partial^{n}+\cdots+$ $a_{1} \partial+a_{0}$ with $a_{i} \in K$.

A point $p \in \mathbb{P}^{1}=\mathbb{C} \cup\{\infty\}$ is called a singularity of a differential operator $L \in K[\partial]$, if $p$ is a zero of the leading coefficient of $L$ or $p$ is a pole of one of the other coefficients of $L . p$ is called a regular point if it is not a singularity.

We denote the solution space of a differential operator as $V(L)=\{y \mid L(y)=0\}$ where the $y$ are taken in some universal extension [15] of $\mathbb{C}(x)$. If $p$ is a regular point of $L$, we can write all solutions of $L$ at $p$ as convergent power series $\sum_{i=0}^{\infty} a_{i} t_{p}^{i}$, where $t_{p}$ denotes the local parameter which is $t_{p}=\frac{1}{x}$ if $p=\infty$ and $t_{p}=x-p$, otherwise.

### 2.2 Transformations

There are three known types of transformations that send, for any second order $L_{1} \in K[\partial]$, the solution space of $L_{1}$ to the solution space of some $L_{2} \in K[\partial]$, again of order 2 . They are (notation as in [4]):
(i) change of variables: $y(x) \rightarrow y(f(x)), \quad f(x) \in K \backslash \mathbb{C}$.
(ii) exp-product: $y \rightarrow e^{\int r d x} \cdot y, \quad r \in K$.
(iii) gauge transformation: $y \rightarrow r_{0} y+r_{1} y^{\prime}, \quad r_{0}, r_{1} \in K$.

Definition 1. Let $L_{1}, L_{2} \in K[\partial]$ with order 2. They are called gauge equivalent (notation: $L_{1} \sim_{g} L_{2}$ ) if there exists a
so-called gauge transformation from $V\left(L_{1}\right)$ to $V\left(L_{2}\right)$, which means a bijection of the form (iii).

Remark 1. Let $L_{1}, L_{2} \in K[\partial]$. The $\mathcal{D}$-modules $\mathcal{D} / \mathcal{D} L_{i}$, $i=1,2$ are isomorphic if and only if $L_{1} \sim_{g} L_{2}$. In particular, $\sim_{g}$ is an equivalence relation (see [1]).

Definition 2. Let $L_{1}, L_{2} \in K[\partial]$ with order 2. They are called projectively equivalent (notation: $L_{1} \sim_{p} L_{2}$ ) if there exists a bijection $V\left(L_{1}\right) \rightarrow V\left(L_{2}\right)$ of the form

$$
\begin{equation*}
y \longrightarrow e^{\int r} \cdot\left(r_{0} y+r_{1} y^{\prime}\right) \tag{1}
\end{equation*}
$$

for some $r, r_{0}, r_{1} \in K$.
Projective equivalence is also an equivalence relation, see [1]. An implementation (for order 2) is given in [8] to decide if $L_{1} \sim_{p} L_{2}$, and if so, to find the projective equivalence (the $r, r_{0}, r_{1}$ in (1)). An algorithm for arbitrary order $n$ was given in [1] (implemented in ISOLDE).

### 2.3 2-descent

Definition 3. Let $f=\frac{A}{B}$ with $A, B \in \mathbb{C}[x]$ coprime, then the degree of $f$ is defined as

$$
\operatorname{deg}(f)=\max (\operatorname{deg}(A), \operatorname{deg}(B))=[\mathbb{C}(x): \mathbb{C}(f)]
$$

Remark 2. If $\sigma \in \operatorname{Aut}(\mathbb{C}(x) / \mathbb{C})$ has order 2, then the fixed field of $\sigma$ is a subfield of $\mathbb{C}(x)$ of index 2, and by Lüroth's theorem this subfield is of the form $\mathbb{C}(f)$, for some $f \in \mathbb{C}(x)$ of degree 2 (note: we can find such $f$ in $\{x+$ $\sigma(x), x \sigma(x)\} \backslash C)$. Any subfield $\mathbb{C}(f) \subset \mathbb{C}(x)$ of index 2 is the fixed field of some $\sigma \in \operatorname{Aut}(\mathbb{C}(x) / \mathbb{C})$ of order 2 (after all, every extension of degree 2 is Galois). The automorphisms of $\mathbb{C}(x)$ over $\mathbb{C}$ are Möbius transformations:

$$
\begin{equation*}
x \mapsto \frac{a x+b}{c x+d} \tag{2}
\end{equation*}
$$

This paper treats 2 -descent, so we only consider $\sigma$ of order 2 , which is equivalent to having $d=-a$ in (2).

Remark 3. Any $\sigma \in \operatorname{Aut}(\mathbb{C}(x) / \mathbb{C})$ extends to an automorphism of $\mathbb{C}(x)[\partial]$. If $\sigma$ has finite order, and if $\mathbb{C}(f)$ is the fixed field of $\sigma$, and if $L \in \mathbb{C}(x)[\partial]$, then

$$
\begin{equation*}
L=\sigma(L) \Longleftrightarrow L \in \mathbb{C}(f)\left[\partial_{f}\right], \tag{3}
\end{equation*}
$$

in other words, $\mathbb{C}(f)\left[\partial_{f}\right]$ is the fixed ring of $\sigma$. Here $\partial_{f}:=$ $\frac{d}{d f}=\frac{1}{f^{\prime}} \partial$, where ' is differentiation w.r.t. $x$.

Definition 4. Let $L \in \mathbb{C}(x)[\partial]$. We say that $L$ has 2 descent if $\exists f \in \mathbb{C}(x)$ with $\operatorname{deg}(f)=2$ and $\exists \tilde{L} \in \mathbb{C}(f)\left[\partial_{f}\right]$ such that $L \sim_{p} \tilde{L}$.
One could instead use the term "projective 2-descent" for this (because we use projective equivalence $\sim_{p}$ ) but we opted to use the shorter term.

Main goal: Let $L \in K[\partial]$ be irreducible and of order 2 . The goal of this paper is to give an explicit algorithm that can decide if $L$ has 2-descent, and if so, find it (i.e. find $\tilde{L} \in \mathbb{C}(f)\left[\partial_{f}\right]$ with $L \sim_{p} \tilde{L}$ for some $f$ of degree 2). Moreover, if $L$ is defined over some field $C \subset \mathbb{C}$, we should only introduce algebraic extensions of $C$ when necessary.

We will divide our algorithm into several steps. The first step is to find candidates for $\mathbb{C}(f)$ with $\operatorname{deg}(f)=2$. Such a field is the fixed field of a Möbius transformation of order 2.

## 3. MÖBIUS TRANSFORMATIONS

Proposition 1. A Möbius transformation has order 2 if it is of the form $\sigma(x)=\frac{a x+b}{c x-a}$. Such $\sigma$ has 2 fixed points in $\mathbb{C} \cup\{\infty\}$.
One could apply a transformation that moves the fixed points of $\sigma$ to $0, \infty$, which reduces $\sigma$ to the notationally convenient $x \mapsto-x$. Our algorithm does not do this because it can introduce an unnecessary algebraic extension of the constants.

### 3.1 The singularity structure

Definition 5. Let $L \in \mathcal{D}$ have order $n$. Assume $p$ is a singularity of $L$. If there exists a basis of $V(L)$ of the form $e^{\int r^{r}} f_{1}, \ldots, e^{\int r} f_{n}$ where $r \in \mathbb{C}(x)$ and $f_{1}, \ldots, f_{n}$ are analytic at $x=p$, then $p$ is called a removable singularity (also called false singularity). Otherwise $p$ is called $a$ true singularity.

Suppose $p$ is a singularity of $L$. If there exists a projectively equivalent $\tilde{L}$ for which $p$ is a regular point, then $p$ is a removable singularity. The true singularities of $L$ are precisely those $p$ that stay singular when $L$ is replaced by any projectively equivalent operator.

Denote (as in [12, 4]) the (generalized) exponent-difference as $\Delta(L, p)$.

Definition 6. For any true singularity p, denote

$$
\operatorname{type}(L, p):=\left\{\begin{array}{cll}
" \text { irreg }^{\prime \prime} & \text { if } & \Delta(L, p) \notin \mathbb{C} \\
" \text { irrat }^{\prime \prime} & \text { if } & \Delta(L, p) \in \mathbb{C} \backslash \mathbb{Q} \\
e \in\left[0, \frac{1}{2}\right] & \text { if } & \Delta(L, p) \in \mathbb{Q}
\end{array}\right.
$$

Here, $e \in\left[0, \frac{1}{2}\right]$ such that $\Delta(L, p) \in(e+\mathbb{Z}) \cup(-e+\mathbb{Z})$. Then we write the singularity structure of $L$ as

$$
S^{\mathrm{type}}:=\{(p, \operatorname{type}(L, p)) \mid p \text { true sing }\}
$$

Let $\pi_{i}$ project on the $i^{\prime}$ th entry of $S^{\text {type }}$, then $S:=\pi_{1}\left(S^{\text {type }}\right) \subseteq$ $\mathbb{P}^{1}(\mathbb{C})$ denotes the set of true singularities of $L$.

Lemma 1. [12, 4] If $L \sim_{p} \tilde{L} \in \mathcal{D}$ then $L$ and $\tilde{L}$ have the same singularity structure $S^{\text {type }}$.
If $L \in C(x)[\partial]$ for some field $C \subset \mathbb{C}$, we denote:

$$
\begin{aligned}
M_{\mathbb{C}} & :=\left\{\left.\sigma=\frac{a x+b}{c x-a} \right\rvert\, a, b, c \in \mathbb{C} \text { and } \sigma(S)=S\right\} \\
M_{C} & :=\left\{\left.\sigma=\frac{a x+b}{c x-a} \right\rvert\, a, b, c \in C \text { and } \sigma(S)=S\right\} \\
M_{\mathbb{C}}^{\text {type }} & :=\left\{\sigma \in M_{\mathbb{C}} \mid \sigma\left(S^{\text {type }}\right)=S^{\text {type }}\right\} \\
M_{C}^{\text {type }} & :=\left\{\sigma \in M_{C} \mid \sigma\left(S^{\text {type }}\right)=S^{\text {type }}\right\}
\end{aligned}
$$

places $(C):=\{f \in C[x] \mid f$ is monic and irreducible $\} \bigcup\{\infty\}$.
Remark 4. places $(\mathbb{C}) \cong \mathbb{P}^{1}(\mathbb{C})=\mathbb{C} \bigcup\{\infty\}$
If $\sigma \in \operatorname{Aut}(C(x) / C)$ then $\sigma$ acts on places $(C)$ in a natural way, preserving degrees, which are defined as:

$$
\operatorname{deg}(p)=\left\{\begin{aligned}
1 & \text { if } \quad p=\infty ; \\
\operatorname{deg}(p) & \text { if } \quad p \text { is a polynomial } .
\end{aligned}\right.
$$

If $L=a_{n} \partial^{n}+\cdots+a_{0} \partial^{0}$ with $a_{0}, \ldots, a_{n} \in C[x]$, then computing the singularities as a subset of $\mathbb{P}^{1}(\bar{C}) \subset \mathbb{P}^{1}(\mathbb{C})$ would mean computing all roots (the splitting field) of $a_{n}$. The algorithm does not compute this splitting field because
it could have exponentially high degree over $C$. Instead, it uses irreducible factors of $a_{n}$ in $C[x]$ (and the point $\infty$ ) to represent the singularities, then we have the notation $S_{C}^{\text {type }}$ and

$$
M_{C}^{\text {type }}:=\left\{\sigma \in M_{C} \mid \sigma\left(S_{C}^{\text {type }}\right)=S_{C}^{\text {type }}\right\}
$$

To ensure that $S$ is invariant under $\sim_{p}$ it is essential to discard all removable singularities.

Example 1. Let $C=\mathbb{Q}$, and

$$
L:=\partial^{2}+\frac{12 x^{4}+1}{x\left(2 x^{2}-1\right)\left(2 x^{2}+1\right)} \partial-\frac{8}{\left(2 x^{2}-1\right)^{2}}
$$

For this example we find

$$
S^{\text {type }}:=\left\{(\infty, 0),(0,0),\left(\frac{-1}{\sqrt{2}}, 0\right),\left(\frac{1}{\sqrt{2}}, 0\right),\left(\frac{-1}{\sqrt{-2}}, 0\right),\left(\frac{1}{\sqrt{-2}}, 0\right)\right\}
$$

The set of true singularities is

$$
S=\pi_{1}\left(S^{\text {type }}\right)=\left\{\infty, 0, \frac{1}{\sqrt{2}}, \frac{-1}{\sqrt{2}}, \frac{1}{\sqrt{-2}}, \frac{-1}{\sqrt{-2}}\right\}
$$

Written in terms of places $(\mathbb{Q})$ it becomes

$$
\begin{aligned}
S_{C} & :=\left\{\infty, x, x^{2}+\frac{1}{2}, x^{2}-\frac{1}{2}\right\} \subset \operatorname{places}(\mathbb{Q}) \\
S_{C}^{\mathrm{type}} & :=\left\{(\infty, 0),(x, 0),\left(x^{2}+\frac{1}{2}, 0\right),\left(x^{2}-\frac{1}{2}, 0\right)\right\}
\end{aligned}
$$

and

$$
M_{C}^{\text {type }}=\left\{-x, \frac{1}{2 x}, \frac{-1}{2 x}\right\}
$$

This example was quite easy because it has obvious 2 descent. Moreover, all singularities were true singularities with type $(L, p)=0$. Removable singularities are common in larger examples, such as Example 3 in Section 7. Using $S$ instead of $S_{C}$ would have introduced an extension of $C=\mathbb{Q}$ of degree 4 in this example, however, such an extension could have been much larger (e.g. if $x^{5}-x-1$ had appeared in the denominator of $L$, which has a splitting field of degree 120).

### 3.2 Finding candidates for $\sigma$

For $i=1,2, \ldots$, let $S_{i}$ denote the set of all $p \in S_{C}$ with $\operatorname{deg}(p)=i$.

Algorithm: Compute Möbius transformations.
Input: The singularity structure $S_{C}^{\text {type }}$.
Output: The set $M_{C}^{\text {type }}$, i.e., the set of all $\sigma \in \operatorname{Aut}(C(x) / C)$ of order 2 that fix $S_{C}^{\text {type }}$. (In this paper we omit 2-descent for $\sigma$ 's that are not defined over $C$ because in that case is better to compute a larger descent, of type $C_{2} \times C_{2}, D_{n}, A_{4}$, $S_{4}$, or $A_{5}$ ).
Step 1: Compute $S_{i}$ from $S_{C}^{\text {type }}$ and let $n_{i}$ denote the number of elements of $S_{i}$.

Step 2: Let $n_{\text {sing }}:=\sum i n_{i}$ (the total number of true singularities when counted in $\mathbb{P}^{1}(\bar{C})$ ).

Step 3: If $n_{\text {sing }}<3$ then return "With $<3$ singularities, descent is not necessary nor implemented" and stop.
Step 4: Now $n_{\text {sing }} \geq 3$.
(i) If $n_{1} \geq 3$, then call Case 1
(ii) If $n_{1}=1, n_{2}=1$, then call Case2
(iii) If $n_{1}=2, n_{2}=1$, then call Case3
(iv) If $n_{2} \geq 2$, then call Case4
(v) If $n_{i} \geq 1$ for some $i \geq 3$, then call Case5

Algorithm: Case1.
Input: $S_{C}^{\text {type }}$ with $S_{1}$ having $\geq 3$ elements.
Output: The set $M_{C}^{\text {type }}$.
Before describing Algorithm Case1, first some remarks. In general $\sigma=\frac{a x+b}{c x+d}$ is determined by the image of three points $\sigma\left(p_{1}\right), \sigma\left(p_{2}\right), \sigma\left(p_{3}\right)$. Since we assume $|\sigma|=2$, we can write $\sigma=\frac{a x+b}{c x-a}$. In general, such $\sigma$ is determined by two points $\sigma\left(p_{1}\right), \sigma\left(p_{2}\right)$ except in one case: when $\sigma\left(p_{1}\right)=p_{2}, \sigma\left(p_{2}\right)=$ $p_{1}$. In that case one more point is needed to determine $\sigma=\frac{a x+b}{c x-a}$.

Algorithm Case1 will choose a pair $p_{1}, p_{2} \in S_{1}\left(p_{1} \neq p_{2}\right)$ and loops over all $n(n-1)$ pairs $q_{1}, q_{2} \in S_{1}\left(q_{1} \neq q_{2}\right)$. If the types of $q_{1}, q_{2}$ match those of $p_{1}, p_{2}$, the algorithm will compute the $\sigma$ that maps $p_{1}, p_{2}$ to $q_{1}, q_{2}$. In the one case that $q_{1}, q_{2}=p_{2}, p_{1}$, a third point $p_{3}$ is used to determine $\sigma$. There are $n-2$ choices for $\sigma\left(p_{3}\right)$, namely from $S_{1}-\left\{p_{1}, p_{2}\right\}$. The number of computed $\sigma$ 's is then $\leq n(n-1)-1+(n-2)$ (equality if they all have the same type). Then we remove those $\sigma$ for which $S_{C}^{\text {type }}$ is not $\sigma$-invariant (That means remove all $\sigma$ 's that send a true singularity to a non-singular point or to a false singularity (Definition 5), and, remove all $\sigma$ 's that send a singularity to a singularity of a different type).

## Algorithm: Case2

Input: $S_{C}^{\text {type }}$ with $S_{1}$ having 1 element and $S_{2}$ having 1 element.
Output: The set $M_{C}^{\text {type }}$.
Step 1: Let the polynomial in $S_{2}$ be $x^{2}+c_{1} x+c_{0}$.
Step 2: Write $\sigma_{1}=-\frac{c_{1} x+2 c_{0}}{2 x+c_{1}}$ and $\sigma_{2}=\frac{a x+c_{0} c+c_{1} a}{c x-a}$.
Remark 5. $\sigma_{1}$ is the unique Möbius transformation of order 2 that fixes the roots of $x^{2}+c_{1} x+c_{0} ; \sigma_{2}$ is the parameterized family of all $\sigma$ of order 2 that swap the roots of $x^{2}+c_{1} x+c_{0}$.

Step 3: Let $p_{1}$ be the one element of $S_{1}$. Equating $\sigma\left(p_{1}\right)$ to $p_{1}$ gives a linear equation that determines the values of the homogeneous parameters $a, c$ in $\sigma_{2}$.
Step 4: Check which (if any) of $\sigma_{1}, \sigma_{2}$ fix $S_{C}^{\text {type }}$ and return those.

Algorithm Case3 is similar to Algorithm Case2.
Algorithm: Case4
Input: $S_{C}^{\text {type }}$ with $S_{2}$ having $\geq 2$ elements.
Output: The set $M_{C}^{\text {type }}$.
Step 1: Choose one polynomial from $S_{2}$. Denote it as $f_{1}=$ $x^{2}+c_{1} x+c_{0}$.

Step 2: Do the following substeps $1-4$ to get the set $T_{1}$ :

1. Write $\sigma_{1}=-\frac{c_{1} x+2 c_{0}}{2 x+c_{1}}$ and $\sigma_{2}=\frac{a x+c_{0} c+c_{1} a}{c x-a}$ (See the Remark in Algorithm Case2).
2. Choose another polynomial in $S_{2}$, and denote it as $f_{2}=x^{2}+d_{1} x+d_{0}$.
3. Write $\sigma_{3}=-\frac{d_{1} x+2 d_{0}}{2 x+d_{1}}$ and $\sigma_{4}=\frac{a x+d_{0} c+d_{1} a}{c x-a}$.
4. Let $a:=d_{0}-c_{0}, c:=c_{1}-d_{1}$, then $\sigma_{2}=\sigma_{4}$ swaps the roots of $f_{1}$ as well as the roots of $f_{2}$. $T_{1}:=\left\{\sigma \in\left\{\sigma_{1}, \sigma_{2}, \sigma_{3}\right\} \mid \sigma\right.$ fixes $\left.S_{C}^{\text {type }}\right\}$.

Step 3: Denote the polynomials in $S_{2}$ as $f_{i}$, then $T_{2}:=$ $\bigcup_{i=2}^{n_{2}} \operatorname{FindMaps}\left(f_{1}, f_{i}\right)$
(See below for the subalgorithm FindMaps)
Step 4: $T_{3}:=\bigcup_{i=3}^{n_{2}} \operatorname{FindMaps}\left(f_{2}, f_{i}\right)$.
Step 5: $T_{1} \bigcup T_{2} \bigcup T_{3}$.
Remark. Taking a set union means removing duplicates. The duplicates are the elements of $T_{3}$ that do not swap the roots of $f_{1}$, and $\sigma_{3}$ might also be duplicate (it could be in $T_{2}$ if $n_{2}>2$ ).

Subalgorithm: FindMaps
Input: Two irreducible polynomials $f, g \in C[x]$ of degree 2 .
Output: All $\sigma \in M_{C}^{\text {type }}$ that map roots of $f$ to roots of $g$.

1. Compute the roots of $g$ in $C(\alpha) \cong C[x] /(f)$. (Note: there are either 0 or 2 roots $\beta_{j}$ )
2. For each root $\beta_{j}$, compute $a, b, c \in C$ (not all 0 ) with $\frac{a \alpha+b}{c \alpha-a}=\beta_{j}$.
This is done by computing coefficients (w.r.t $\alpha$ ) of $a \alpha+$ $b-\beta_{j}(c \alpha-a)$ and equating them to 0 .
3. For each $\frac{a x+b}{c x-a}$ found in step 2 check if it fixes $S_{C}^{\text {type }}$, if so, include it in the output.

Algorithm: Case5
Input: $S_{C}^{\text {type }}$ with $S_{i}$ having $\geq 1$ elements and $i \geq 3$.
Output: The set $M_{C}^{\text {type }}$.
Step 1: Find $S_{i}$ for an $i \geq 3$ with $n_{i}>0$.
Step 2: Choose a polynomial $f$ in $S_{i}$. Denote $C(\alpha) \cong$ $C[x] /(f)$, with $f(\alpha)=0$.

Step 3: For each polynomial $g \in S_{i}$, call $\operatorname{FindMaps}(f, g)$. Then $M_{C}^{\text {type }}$ would be $\bigcup_{g \in S_{i}}$ FindMaps $(f, g)$.

## 4. COMPUTING 2-DESCENT, CASE A

Notations: Let $L \in C(x)[\partial]$ have order 2, and be irreducible (even in $\mathbb{C}(x)[\partial]$ ). Let $\sigma \in \operatorname{Aut}(C(x) / C$ ) have order 2 and fixed field $C(f) \subset C(x)$.

Lemma 2. If $\exists \tilde{L} \in \mathbb{C}(f)\left[\partial_{f}\right]$ with $L \sim_{p} \tilde{L}$, then $L \sim_{p}$ $\sigma(L)$.

Proof. $L \sim_{p} \tilde{L}=\sigma(\tilde{L}) \sim_{p} \sigma(L)$.
So if not $L \sim_{p} \sigma(L)$ then $L \in C(x)[\partial] \subset \mathbb{C}(x)[\partial]$ does not descend to $\mathbb{C}(f)$. If $L \sim_{p} \sigma(L)$ then we will consider two cases:

Notation 1. Case $\boldsymbol{A}$ is when there exists $G=r_{0}+$ $r_{1} \partial \in \mathbb{C}(x)[\partial]$ such that $G(V(L))=V(\sigma(L))$, i.e. $L \sim_{g}$ $\sigma(L)$.

$$
\text { Case B is when there exists } G=e^{\int^{r}} .\left(r_{0}+\right.
$$ $\left.r_{1} \partial\right)$ such that $\overline{G(V(L))}=V(\sigma(L))$, i.e. $L \sim_{p} \sigma(L)$.

(Note: Case $A \Rightarrow$ Case B.)
This section treats only Case A. Section 5 will reduce Case B to Case A.

In Case A, when $L \sim_{g} \sigma(L)$, it is known [11] that there exists $\tilde{L} \in \mathbb{C}(f)\left[\partial_{f}\right]$ with $\tilde{L} \sim_{g} L$. Then we have the following diagram:
Diagram 1


Here, $A, \sigma(A)$, and $\tilde{L}$ are unknown. Whether or not such a diagram commutes is studied in Theorem 1 below.

Remark 6. A gauge transformation is a bijective map $A: V(L) \rightarrow V(\tilde{L})$ that can be represented by a differential operator in $\mathbb{C}(x)[\partial]$. So we can define $\sigma(A)$ simply by applying $\sigma$ to the operator that represents the map $A$.

Theorem 1. Let $L$ and $\sigma$ be as before, and $G: V(L) \rightarrow$ $V(\sigma(L))$ be a gauge transformation. Suppose $\tilde{L_{1}}, \tilde{L_{2}} \in \mathbb{C}(f)\left[\partial_{f}\right]$ and $A_{i}: V(L) \rightarrow V\left(\tilde{L_{i}}\right)$ are gauge transformations. Then:

1. For each $i=1,2$, there is exactly one $\lambda_{i} \in \mathbb{C}^{*}$ such that the following diagram commutes.
Diagram 2

2. If $\tilde{L_{1}} \sim_{g} \tilde{L_{2}}$ over $\mathbb{C}(f)$, then $\lambda_{1}=\lambda_{2} ;$ Otherwise, $\lambda_{1}=$ $-\lambda_{2}$.
3. In particular, $\left\{\lambda_{1},-\lambda_{1}\right\}$ depends only on $(L, \sigma, G)$.

Proof. First consider the diagram without $\lambda_{i}$ in it. In it we find two gauge transformations $V(L) \rightarrow V\left(\tilde{L}_{i}\right)$, namely $A_{i}$ and $\sigma\left(A_{i}\right) G$. After choosing bases of $V(L)$ and $V\left(\tilde{L_{i}}\right)$, we can view these gauge transformations as bijections: $\mathbb{C}^{2} \rightarrow$ $\mathbb{C}^{2}$. Then by linear algebra, there is a constant $\lambda_{i} \in \mathbb{C}^{*}$ such that the map:

$$
\begin{equation*}
A_{i}-\lambda_{i} \sigma\left(A_{i}\right) G: V(L) \rightarrow V\left(\tilde{L}_{i}\right) \tag{4}
\end{equation*}
$$

has a non-zero kernel. The kernel of (4) corresponds to a right hand factor of $L$, namely, the GCRD of $L$ and the operator in (4). However, $L$ is irreducible so this kernel must be $V(L)$ itself. That means Diagram 2 commutes. That $\lambda_{i}$ is unique follows from linear algebra: there can be at most one $\lambda_{i}$ for which (4) is the zero map. Item 1 follows.

For item 2, since $\tilde{L_{1}} \sim_{g} L \sim_{g} \tilde{L_{2}}$, there exists a gauge transformation $B: V\left(\tilde{L_{1}}\right) \rightarrow V\left(\tilde{L_{2}}\right)$. This $B$ is unique up
to multiplying by a constant that we choose in such a way that the composition $B A_{1}: V(L) \rightarrow V\left(\tilde{L_{2}}\right)$ coincides with $A_{2}$. Since $\sigma\left(\tilde{L_{1}}\right)=\tilde{L_{1}}, \sigma\left(\tilde{L_{2}}\right)=\tilde{L_{2}}$ one sees that $\sigma(B)$ maps $V\left(\tilde{L_{1}}\right)$ to $V\left(\tilde{L_{2}}\right)$ as well. So $\sigma(B)$ must be $c \cdot B$ for some $c \in \mathbb{C}^{*}$. Then $|\sigma|=2$ implies that $c= \pm 1$. Now $c=1 \mathrm{iff}$ $\sigma(B)=B$ iff $B \in \mathbb{C}(f)\left[\partial_{f}\right]$ iff $\tilde{L_{1}}, \tilde{L_{2}}$ are gauge-equivalent over $\mathbb{C}(f)$. Otherwise, if $c=-1$, then $B \notin \mathbb{C}(f)\left[\partial_{f}\right]$ and $\tilde{L_{1}}, \tilde{L_{2}}$ are gauge-equivalent over $\mathbb{C}(x)$ but not over $\mathbb{C}(f)$. To prove item 2 we now have to show that $\lambda_{2}=c \lambda_{1}$.

If $\lambda_{i}$ is such that Diagram 2 commutes (for $i=1,2$ ) then the following diagram commutes:

## Diagram 3



The composed map $B A_{1}$ at the left of Diagram 3 coincides with the map $A_{2}$ in Diagram 2 for $i=2$. Applying $\sigma$ to $B A_{1}$ and $A_{2}$, we see that the composed map at the right of Diagram 3 coincides with the map $\sigma\left(A_{2}\right)$ in Diagram 2 for $i=2$. Then the maps at the top of Diagram 3 and Diagram 2 for $i=2$ must coincide as well, i.e., $\lambda_{2} G=c \lambda_{1} G$. Hence $\lambda_{2}=c \lambda_{1}$. Item 2 (and hence item 3) follow.

### 4.1 Algorithm for finding 2-descent in Case A

Notations $L, C, G, \sigma, A$ are as in Section 4. Our goal is to compute 2-descent: $L \sim_{p} \tilde{L} \in \mathbb{C}(f)\left[\partial_{f}\right]$. Here $f$ is determined from $\sigma$ as in Remark 2. We will compute $A: V(L) \rightarrow V(\tilde{L})$ first, then use $A$ to find $\tilde{L}$.

Algorithm: Case A for computing a 2-descent $\tilde{L}$ for $L$.
Input: $L, G, \sigma$ and $C$.
Output: $\tilde{L}$ and $A$, defined over an optimal extension of $C$.
Step 1: Write $A=\left(a_{00}+a_{01} x\right) \partial+\left(a_{10}+a_{11} x\right)$, with $a_{00}$, $a_{01}, a_{10}, a_{11}$ unknowns (which will take values in $\mathbb{C}(f)$ ).

Step 2: The operator $A-\lambda \sigma(A) G$ in (4) should vanish on $V(L)$, so the remainder of $A-\sigma(A) \lambda G$ right divided by $L$ must be 0 . This remainder is of the form $\left(R_{00}+\right.$ $\left.R_{01} x\right) \partial^{0}+\left(R_{10}+R_{11} x\right) \partial$, where the $R_{i j}$ are $C(\lambda, f)$ linear combinations of $a_{i j}$. This produces a system of 4 equations $R_{i j}=0$ in 4 unknowns $a_{i j}$.

Step 3: To have a nontrivial solution, the corresponding $4 \times 4$ matrix $M$ must have determinant 0 . Equating $\operatorname{det}(M)$ to 0 gives a degree 4 equation for $\lambda$. Solve for $\lambda$.
Remark. The equation for $\lambda$ is of the form $\left(\lambda^{2}-a\right)^{2}=$ 0 , where $a=\lambda_{1}^{2}=\lambda_{2}^{2}$ with $\lambda_{1}, \lambda_{2}$ as in Theorem 1. If $L$ and $\sigma$ are defined over a field $C \subseteq \mathbb{C}$ then $\tilde{L}$ and $A$ are defined over $C(\sqrt{a})$.
If $\sqrt{a} \notin C$ then it follows from Theorem 1 that the extension by $\lambda_{i}= \pm \sqrt{a}$ is necessary.

Step 4: Plug in one value for $\lambda$ in $M$, then solve $M$ to find values for $a_{00}, a_{01}, a_{10}, a_{11}$ in $C(\sqrt{a}, f)$.

Step 5: Compute $\operatorname{LCLM}(A, L)$ to obtain $\tilde{L} A$. Right divide by $A$ to find $\tilde{L} \in C(\sqrt{a}, f)\left[\partial_{f}\right]$.

Step 6: (optional) Introduce a new variable, say $x_{1}$, and compute an operator $L_{x_{1}} \in C\left(\sqrt{a}, x_{1}\right)\left[\partial_{x_{1}}\right]$ that corresponds to $\tilde{L}$ under the change of variables $x_{1} \mapsto f$.

## 5. COMPUTING 2-DESCENT, CASE B

Definition 7. Let $L_{1}, L_{2} \in \mathcal{D}=K[\partial]$. The symmetric product $L_{1}\left(L_{2}\right.$ is defined as the monic differential operator in $\mathcal{D}$ with minimal order for which $y_{1} y_{2} \in V\left(L_{1}\left(L_{2}\right)\right.$ for all $y_{1} \in V\left(L_{1}\right), y_{2} \in V\left(L_{2}\right)$.

Lemma 3. If $L=\partial^{2}+c_{0} \in C(x)[\partial]$, and $G:=e^{\int r} \cdot\left(r_{0}+\right.$ $r_{1} \partial$ ) is a bijection from $V(L)$ to $V(\sigma(L))$, then $\left(e^{\int r}\right)^{2}$ is a rational function.
If $L:=\partial^{2}+a_{1} \partial+a_{0} \in \mathbb{C}(x)[\partial]$, then $L_{1}:=L(S)\left(\partial-\frac{1}{2} a_{1}\right)$ is of the form $\partial^{2}+c_{0}$ (with $c_{0}=a_{0}-\frac{1}{4} a_{1}^{2}-\frac{1}{2} a_{1}^{\prime}$ ).

The proof of the lemma follows by computing the effect of $G$ on the Wronskian, and the fact that the Wronskians of $\partial^{2}+c_{0}$ and $\sigma\left(\partial^{2}+c_{0}\right)$ are rational functions (1 and $\sigma(x)^{\prime}$ respectively).

Let $L \in C(x)[\partial]$ irreducible (even over $\mathbb{C}$ ) and of order 2 , and $\sigma \in \operatorname{Aut}(C(x) / C)$ of order 2. The implementation equiv [8] can check if $L \sim_{p} \sigma(L)$, and if so, find $r, r_{0}, r_{1} \in C(x)$ for which $G:=e^{\int r} \cdot\left(r_{0}+r_{1} \partial\right)$ is a bijection from $V(L)$ to $V(\sigma(L))$. Assume that such $\sigma$ and $G$ are given. After the simple transformation in the lemma above, we may assume that $\left(e^{\int r}\right)^{2}$ is a rational function.

If $e^{\int r}$ itself is a rational function, then we are in Case A. Otherwise, we can write $e^{\int r}=p(x) \sqrt{f(x)}$ for some squarefree polynomial $f(x)$, and some $p(x) \in C(x)$.

Definition 8. The branch points of $G$ are the roots of $f(x)$, and $\infty$ if $f(x)$ has odd degree.
To reduce Case B to Case A, we have to eliminate the branch points. Our algorithm will first eliminate all branch points that can be eliminated without a field extension of $C$. It will only extend $C$ if there is no descent w.r.t. $\sigma$ defined over $C$.

### 5.1 Branch points

It is convenient to view the set of branch points as a subset of $\mathbb{P}^{1}(\bar{C})$. However, to avoid splitting fields, the algorithm represents the branch points with a set $B \subset \operatorname{places}(C)$ instead. This $B$ is the set of irreducible factors of $f(x)$ in $C[x]$, as well as $\infty$ if $f(x)$ has odd degree. The goal is to eliminate branch points until we reach $B=\emptyset$, i.e., Case A.

Definition 9. If $\sigma(\infty)=\infty$, then denote $\operatorname{Inf}:=\{\infty\}$, otherwise $\operatorname{Inf}:=\{\infty, x-\sigma(\infty)\}$. Denote $B_{I}=B \bigcap \operatorname{Inf}$ and $B_{N}=B \backslash B_{I}$.
Let $f_{1}(x), f_{2}(x) \in B_{N}$. We say that $f_{1}(x)$ matches $f_{2}(x)$ when the roots of $f_{2}(x)$ are the same as the roots of $f_{1}(\sigma(x))$ (i.e. the numerator of $f_{1}(\sigma(x))$ is $f_{2}$ ).

If $\sigma(\infty) \neq \infty$, then we say that the polynomial $x-\sigma(\infty)$ matches $\infty$.

Lemma 4. If $f_{1}(x) \neq f_{2}(x) \in B_{N}$ and $f_{1}(x)$ matches $f_{2}(x)$, then $B_{N}$ turns into $B_{N} \backslash\left\{f_{1}, f_{2}\right\}$ when we replace $L$ by $L_{\text {new }}:=L(S)\left(\partial-\frac{1}{2} \cdot \frac{f_{1}(x) \prime}{f_{1}(x)}\right)$.

Proof. The composed transformation

$$
V\left(L_{\text {new }}\right) \rightarrow V(L) \rightarrow V(\sigma(L)) \rightarrow V\left(\sigma\left(L_{\text {new }}\right)\right)
$$

is

$$
\sqrt{\sigma\left(f_{1}\right)} \cdot G \cdot \frac{1}{\sqrt{f_{1}}}
$$

The polynomial $f$ equals $f_{1} f_{2} \cdots$ where the $\cdots$ refer to the other factors of $f$ in $B \backslash\{\infty\}$. The transformation $G$ is of the form $\sqrt{f_{1} f_{2} \cdots} \cdot\left(r_{0}+r_{1} \partial\right)$. Factors can be removed from the square-root in $G$ either by division or by multiplication by a square-root (factors in $C(x)$ can be moved to $r_{0}, r_{1}$ ). So in the composed transformation, the factors $f_{1}$ and $f_{2}$ will disappear from the square-root in $G$ (note: this uses the assumption $f_{1} \neq f_{2}$ (which implies that their gcd is 1 since they are monic irreducible polynomials)).
A subtlety is that if $\sigma(\infty) \neq \infty$, then $\sigma\left(f_{1}\right)$ is not $f_{2}$ but $c f_{2} /(x-\sigma(\infty))^{d}$, for some $c \in C$, where $d$ is the degree of $f_{1}$ and $f_{2}$. This means that if $\sigma(\infty) \neq \infty$ and $d$ is odd, then the set $B_{I}$ will change when we replace $L$ by $L_{\text {new }}$ ( $B_{I}=\emptyset$ will change to Inf, and $B_{I}=\operatorname{Inf}$ will change to $\left.\emptyset\right)$.

Lemma 5. If $\sigma(\infty) \neq \infty$, and $B_{I}=\left\{\infty, f_{1}\right\}$ (here $f_{1}=$ $x-\sigma(\infty)$ ) then the factor $f_{1}$ inside the square root in $G$ will cancel out (i.e. $B_{I}$ will become $\emptyset$ ) if we replace $L$ by $L_{\text {new }}:=L(S)\left(\partial-\frac{1}{4} \cdot \frac{1}{f_{1}}\right)$.

Proof. The solutions of $L_{\text {new }}$ differ a factor $\sqrt[4]{f_{1}}$ from the solutions of $L$. The lemma follows from a similar computation as the proof of Lemma 4, except that this time $\sigma\left(f_{1}\right)$ is of the form $c / f_{1}$ for some constant $c$. Thus, the composed map is of the form $\sqrt[4]{c / f_{1}} \cdot G \cdot 1 / \sqrt[4]{f_{1}}$, and $\sqrt{f_{1}}$ is cancelled from the square root in $G$.

In the following algorithm, $L$ and $\sigma$ are as in Section 4, and $G=e^{\int r} \cdot\left(r_{0}+r_{1} \partial\right)$ with $r, r_{0}, r_{1} \in C(x)$.

Algorithm: Case B for computing a 2-descent $\tilde{L}$ for $L$.
Input: $L, G, \sigma$ and $C$.
Output: $\tilde{L}$ and $A$ (defined over $C$ whenever possible).
Step 1 Initialization: If $\left(e^{\int r}\right)^{2}$ is not a rational function, then replace $L$ by $L \subseteq\left(\partial-\frac{1}{2} \cdot \frac{a_{1}}{a_{2}}\right)$ as in Lemma 3 and update $G$ accordingly.
Rewrite $G$ as $\sqrt{f(x)}\left(r_{0}+r_{1} \partial\right)$ with $f(x)$ monic and square-free (updating $r_{0}, r_{1} \in C(x)$ to move any rational factor from $e^{\int r}$ to $\left.r_{0}, r_{1}\right)$. If $f(x)=1$ then call Case $\mathbf{A}$ and stop.

Step 2: Factor $f(x)$ in $C[x]$ to find $B, B_{I}, B_{N} \subset \operatorname{places}(C)$.
Step 3: $g:=\boldsymbol{\operatorname { F i n d g }}\left(B_{N}, \sigma, C\right)$.
(See below for the subalgorithm Findg)
Step 4: Let $h:=\frac{1}{2} \cdot \frac{g^{\prime}}{g}$. Replace $L$ by $L(S)(\partial-h)$ and update $G, B, B_{I}, B_{N}$ accordingly. Now $B_{N}$ should be $\emptyset$.
Step 5: If $B_{I} \neq \emptyset$ then let $h:=\frac{1}{4} \cdot \frac{1}{f_{1}}$ with $f_{1}$ as in Lemma 5 . Replace $L$ by $L S(\partial-h)$ and update $G, B$ accordingly. Now $B$ should be $\emptyset$.

## Step 6: Call Case A.

Subalgorithm: Findg.
Input: $B_{N}, \sigma, C$.
Output: $g$.

Step 1: If $B_{N}=\emptyset$, return 1 and stop.
Step 2: Else, for each $P_{i} \in B_{N}$,

1. Find its matched (Def. 9) element $P_{j} \in B_{N}$.
2. If $P_{i} \neq P_{j}$ then $g:=\boldsymbol{F i n d g}\left(B_{N} \backslash\left\{P_{i}, P_{j}\right\}, \sigma, C\right)$, return $g \cdot P_{i}$ and stop.

Step 3: Now each $P \in B_{N}$ matches itself, and hence has even degree. Choose $P \in B_{N}$ with minimal degree, and let $\alpha \in \bar{C}$ be one root of $P$, so $C(\alpha) \cong C[x] /(P)$. Let $B_{N}^{\alpha}$ be the set of all irreducible factors in $C(\alpha)[x]$ of all elements of $B_{N}$. Return Findg $\left(B_{N}^{\alpha}, \sigma, C(\alpha)\right)$.

## 6. MAIN ALGORITHM

## Algorithm 2-descent.

Input: A second order irreducible differential operator $L \in$ $C(x)[\partial]$ and the field $C$.
Output: descent, if it exists for some $\sigma \in \operatorname{Aut}(C(x) / C)$ of order 2.

Step 1: Compute the set of true singularities, and the singularity structure $S_{C}^{\text {type }}$.

Step 2: Call Compute Möbius transformations in Section 3.2 to compute the set $M_{C}^{\text {type }}$.
Step 3: For each $\sigma \in M_{C}^{\text {type }}$, call [8] to check if $L \sim_{p} \sigma(L)$, and if so, to find $G: V(L) \rightarrow V(\sigma(L))$.
If we find $\sigma$ with $L \sim_{p} \sigma(L)$, then call algorithm Case B in Section 5.1 and stop.

## 7. EXAMPLES

We give two examples. The first example is easy (it has $G=r_{0}+r_{1} \partial$ with $r_{1}=0$ ). The second one is less trivial ${ }^{3}$. The first example is in Case A as in Section 4, the second example involves both Case A and Case B.

Example 2. Let

$$
L=\partial^{2}+\frac{28 x-5}{x(4 x-1)} \partial+\frac{144 x^{2}+20 x-3}{x^{2}(4 x-1)(4 x+1)}
$$

Step 1: Compute the singularity structure of $L$

$$
S_{C}^{\mathrm{type}}:=\left\{(x, 0),(\infty, 0),\left(x-\frac{1}{4}, 0\right),\left(x+\frac{1}{4}, 0\right)\right\}
$$

Step 2: Compute Möbius transformations. Since $S_{1}$ has $n_{1}=4$ elements, we end up in algorithm Case1 of Section 3.2 which produces:

$$
\left\{-x, \frac{-1}{16 x}, \frac{1}{16 x}, \frac{-1}{4} \frac{4 x-1}{4 x+1}, \frac{1}{4} \frac{4 x+1}{4 x-1}\right\}
$$

Step 3: There are 5 choices for $\sigma$. The first one is $x \mapsto-x$ corresponding to the subfield $C(f)=C\left(x^{2}\right)$. The equiv [8] program finds $G=\frac{4 x-1}{4 x+1}$. Next we compute $A:=-4 x^{2}+x$, and then $\tilde{L}$. After applying a change of variable $x \mapsto \sqrt{x_{1}}$ the result reads

$$
L_{x_{1}}:=\left(16 x_{1}-1\right) x_{1} \partial^{2}+\left(32 x_{1}-2\right) \partial+4
$$

which has 3 true singularities and is easy to solve.

[^2]Example 3. Consider the operator:

$$
\begin{aligned}
L & :=\partial^{2}+\frac{4\left(1296 x^{5}+576 x^{4}-144 x^{3}-72 x^{2}+x+1\right)}{x(6 x-1)(2 x+1)(6 x+1)\left(12 x^{2}-1\right)} \partial+ \\
& \frac{2\left(5184 x^{6}-864 x^{5}-1656 x^{4}+48 x^{3}+162 x^{2}+6 x-1\right)}{(-1+2 x) x^{2}(6 x-1)(2 x+1)(6 x+1)\left(12 x^{2}-1\right)}
\end{aligned}
$$

Step 1: Compute the singularity structure of $L$
$S_{C}^{\text {type }}:=\left\{(x, 0),(\infty, 0),\left(x-\frac{1}{2}, 0\right),\left(x+\frac{1}{2}, 0\right),\left(x-\frac{1}{6}, 0\right),\left(x+\frac{1}{6}, 0\right)\right\}$
( $12 x^{2}-1$ is a removable singularity, Definition 5).
Step 2: Compute Möbius transformations. Since $S_{1}$ has $n_{1}=6$ elements, we are again in Case1, and find:

$$
\left\{-x, \frac{-1}{12 x}, \frac{1}{12 x}, \frac{-1}{2} \frac{2 x-1}{6 x+1}, \frac{1}{2} \frac{2 x+1}{6 x-1}, \frac{-1}{6} \frac{6 x-1}{2 x+1}, \frac{1}{6} \frac{6 x+1}{2 x-1}\right\}
$$

Step 3: The first $\sigma$ we try is $x \mapsto-x$. The equiv program finds

$$
G:=\frac{x\left(12 x^{2}+4 x-1\right)}{12 x^{2}-1} \partial+\frac{3}{2} \frac{(2 x+1)(10 x-1)}{12 x^{2}-1}
$$

so $G(V(L))=V(\sigma(L))$. Then compute a 4 by 4 matrix from the linear equations for the $a_{i j}$, equate the determinant to 0 and find $\lambda= \pm 2$. We choose $\lambda=2$ and find

$$
A:=\left(-36 x^{4}-\frac{1}{4}+10 x^{2}\right) \partial+1-\frac{1}{4} \frac{\left(288 x^{4}+1-84 x^{2}\right)}{x} .
$$

We get

$$
\begin{aligned}
L_{x_{1}}:= & 4 x_{1}^{2}\left(-1+36 x_{1}\right)\left(4 x_{1}-1\right)\left(12 x_{1}-1\right)^{2} \partial^{2}+ \\
& 8 x_{1}\left(12 x_{1}-1\right)\left(4 x_{1}-1\right)\left(216 x_{1}^{2}-54 x_{1}+1\right) \partial \\
& -3-2544 x_{1}^{2}+10368 x_{1}^{3}+48 x_{1}
\end{aligned}
$$

which is $\tilde{L} \in C\left(x^{2}\right)\left[\partial_{x^{2}}\right]$ rewritten with $x \mapsto \sqrt{x_{1}}$. This $L_{x_{1}}$ has 4 true singularities, and allows a further 2-descent. Applying steps (1)(2)(3) to $L_{x_{1}}$ again, we are actually in Case B as in Section 5, applying the algorithm (details are given in a Maple worksheet [5]) we find a new operator $\tilde{L_{1}} \sim_{p}$ $L_{x_{1}}$ defined over the subfield $\mathbb{C}\left(f_{1}\right)$ where $f_{1}:=x_{1}+\frac{1}{144 x_{1}}$. Replacing $f_{1}$ by a new variable $x_{2}$ we get:

$$
\begin{aligned}
L_{x_{2}}:= & 4\left(36 x_{2}+11\right)\left(18 x_{2}-5\right)\left(6 x_{2}+1\right)\left(6 x_{2}-1\right)^{2} \partial^{2}+ \\
& 36\left(6 x_{2}-1\right)\left(1296 x_{2}^{3}+1620 x_{2}^{2}+20 x_{2}-9\right) \partial \\
& +34992 x_{2}^{3}-207036 x_{2}^{2}-2331+3456 x_{2}
\end{aligned}
$$

which has 3 true regular singularities (as well as a few removable singularities). That means that $L_{x_{2}}$ (and hence L) has closed form solutions (see [5]) in terms of hypergeometric ${ }_{2} F_{1}$ functions.

## 8. FUTURE WORK

At the moment, we only consider $\sigma$ 's that are defined over the same field of constants $C$ over which $L$ is defined. We can modify the Compute Möbius transformations algorithm to also find $\sigma$ 's defined over an extension of $C$. However, for such $\sigma$ we do not plan to compute 2-descent because if there exists descent w.r.t. a $\sigma$ that is not defined over $C$, then a larger descent should exist as well.

We plan to work on finding (if it exists) descent to subfields of index 3. Degree 3 extensions need not be Galois, and so in general, to find 3-descent it is not enough to try all Möbius transformations that fix the singularity structure.

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[^0]:    *Supported by NSF grant 1017880 .

[^1]:    ${ }^{1}$ the number of removable singularities (Def. 5) is irrelevant
    ${ }^{2}$ for the irregular singular case, finding closed form solutions if they exist can be done with $[12,4]$

[^2]:    ${ }^{3}$ it was e-mailed to one of us to find its closed form solutions. There have been many such requests, which motivates us to develop these algorithms.

