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# BRANCHING OF SINGULARITIES FOR SOME SECOND OR THIRD ORDER MICROHYPERBOLIC OPERATORS MULTIPLY CHARACTERISTIC AT $x_{1}=0$ 

Hideshi Yamane<br>Department of Mathematical Sciences, University of Tokyo


#### Abstract

. We study microhyperbolic operators of the form $\left(D_{1}-\alpha_{1} x_{1} D_{n}\right)\left(D_{1}-\alpha_{2} x_{1} D_{n}\right)+$ lower, $\left(D_{1}-x_{1} D_{n}\right) D_{1}\left(D_{1}-x_{1} D_{n}\right)+$ lower. When the lower order terms take a certain form, we can obtain a very detailed information about the singularity of a solution. We look at this problem from the viewpoint of boundary value problems. General arguments for m -th order operators are given.


## INTRODUCTION

For (micro) differential operators whose characteristic variety has a non-involutory intersection, the propagation and branching of singularities is the most interesting problem. $[\mathrm{H}],[\mathrm{O}],[\mathrm{Al}]$ and $[\mathrm{T}-\mathrm{T}]$ are well-known results. The former two are about "Fuchsian" operators and the latter two are about second order hyperbolic operators.

In this paper, we study the branching of singularities for second and third order (micro) hyperbolic operators. The principal symbols are $\left(\xi_{1}-\alpha_{1} x_{1} \xi_{n}\right)\left(\xi_{1}-\alpha_{2} x_{1} D_{n}\right)$, $\left(\xi_{1}-x_{1} \xi_{n}\right) \xi_{1}\left(\xi_{1}+x_{1} \xi_{n}\right)$. General arguments (heuristic in some parts) about m-th order operators are given.

Our approach is based on the study of boundary value problems. Although this viewpoint was already taken in [Al], it is more apparent in the present paper.

Another feature of our approach is that we use ODEs of Fuchs type, while in [Al] and [T-T], ODEs with irregular singularities were used.

The plan of this paper is as follows. PART O gives a general background about an operator of arbitray order. PARTS 1 and 2 are about the second and the third order cases respectively. The main theorems are.found in PART 1 §1 and PART 2 §1.

This paper is an application of [Kat 2].
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## PART 0 GENERAL BACKGROUND

Let $P(x, D)=D_{1}^{m}+P_{1}\left(x, D^{\prime}\right) D_{1}^{m-1}+\cdots+P\left(x, D^{\prime}\right)$ be an m-th order microhyperbolic microdiferential operator defined in a neighborhood of $p \in\left\{(x ; i \xi) \in i T^{*} M ; x_{1}=0, \xi_{n}>\right.$ $0\}, M=\mathbb{R}^{n}$. Assume that its principal symbol is of the form

$$
\sigma(P)(x, \xi)=\left(\xi_{1}-x_{1}^{\lambda} \alpha_{1}\left(x, \xi^{\prime}\right)\right) \ldots\left(\xi_{1}-x_{1}^{\lambda} \alpha_{m}\left(x, \xi^{\prime}\right)\right)
$$

where $\alpha_{j}(a \leq j \leq m)$ is real analytic in $\left(x, \xi^{\prime}\right)$ and $\alpha_{j} \neq \alpha_{k}(j \neq k)$. In addition, we impose the Levi condition:

$$
\begin{aligned}
\operatorname{ord} \frac{\partial^{q} P_{l}}{\partial x_{1}^{q}}\left(0, x^{\prime}, D^{\prime}\right) & \leq \frac{q+l}{\lambda+1}(<l) \\
1 & \leq l \leq m, 0 \leq q<\lambda l .
\end{aligned}
$$

Branching of singularities for this kind of operator has been studied by Alinhac, Taniguchi-Tozaki, Amano and others. A typical method is to apply the partial Fourier transform and reduce the problem to that of an ordinary differential operator with an irregular singular point. Here in the present paper, we choose a different approach. Instead of the partial Fourier transform, we use a singular coordinate change and the quantized Legendre transform. Then the problem is reduced to that of an ordinary differential operator of Fuchs type. In PARTS 1 and 2, we treat the cases of $\lambda=1, m=2,3$. In these cases, we encounter Gauss and Jordan-Pochhammer hypergeometric equations respectively.

We rely on the theory of a coordinate change of fractional order in [Kat 2]. Here we don't give a technical detail.

Let $b_{j}^{ \pm} \subset\left\{ \pm x_{1}>0\right\}$ be the half-bicharacteristic strip of $\xi_{1}=x_{1}^{\lambda} \alpha_{j}$, issuing from $p$. We say that $u \in\left(\Gamma_{\left\{ \pm x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p}$ is a $j$-null solution if $u=0$ on $b_{j}^{ \pm}$, where $\mathcal{C}_{M}^{P}$ is the kernel sheaf of $P: \mathcal{C}_{M}^{P} \rightarrow \mathcal{C}_{M}^{P}$. Let $N u l l(j, \pm)$ be the totality of j -null solutions, that is,

$$
\operatorname{Null}(j, \pm)=\left\{u \in\left(\Gamma_{\left\{ \pm x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p} ; u=0 \text { on } b_{j}^{ \pm}\right\} .
$$

We explain how to construct j-null solutions. For convenience, we only consider the case of + . If $u$ is a solution in $x_{1}>0$, it is mild and its canonical extention $\tilde{u}$ is defined. It satisfies

$$
x_{1}^{m} P \tilde{u}=0 .
$$

Put $t=\frac{1}{\lambda+1} x_{1}^{\lambda+1}$. We denote by $\tau$ the dual variable of $t$. Then $\xi_{1}-x_{1}^{\lambda} \alpha_{j}=x_{1}^{\lambda}\left(\tau-\alpha_{j}\right)$. Since $\alpha_{j}$ 's are mutually distinct (while $x_{1}^{\lambda} \alpha_{j}$ 's are not), the problem has become easier at least from the geometric point of view. Moreover

$$
\begin{aligned}
& \frac{1}{(\lambda+1)^{m}} x_{1}^{m} P(x, D) \\
= & \prod_{j=0}^{m-1}\left(t D_{t}-\frac{j}{\lambda+1}\right)+\sum_{l=1}^{m} \sum_{q=0}^{\infty}(\lambda+1)^{\frac{q-\lambda l}{\lambda+1}} \frac{1}{q!} t^{\frac{q+1}{\lambda+1}} \frac{\partial^{q} P_{l}}{\partial x_{1}^{q}}\left(0, x^{\prime}, D^{\prime}\right) \prod_{j=0}^{m-l-1}\left(t D_{t}-\frac{j}{\lambda+1}\right) .
\end{aligned}
$$

Here assume that there's no contribution by the terms corresponding to
$\frac{q+1}{\lambda+1} \neq 0,1,2,3, \ldots$. Now apply the quantized Legendre transform $\beta_{n}^{+}$with respect to $\left(t, x^{\prime}\right)$. (See [Kat 1] for the definition.) Then we obtain the following operator:

$$
\begin{aligned}
& \tilde{P}= \prod_{j=0}^{m-1}\left(-D_{\zeta} \zeta-\frac{j}{\lambda+1}\right) \\
&+\sum_{l=1}^{m} \sum_{q=0}^{\infty} \frac{(\lambda+1)^{\frac{q-\lambda t}{\lambda+1}}}{q!} \frac{\partial^{q} P_{l}}{\partial x_{1}^{\natural}}\left(0, x^{\prime \prime}, x_{n}+D_{\zeta} \zeta D_{n}^{-1}, D^{\prime}\right) \\
& \times\left(-i D_{\zeta} D_{n}^{-1}\right)^{\frac{q+1}{\lambda+1}} \prod_{j=0}^{m-l-1}\left(-D_{\zeta} \zeta-\frac{j}{\lambda+1}\right), \\
& x^{\prime \prime}=\left(x_{2}, \ldots, x_{n-1}\right) .
\end{aligned}
$$

We will see that $\tilde{P}$ is an ordinary differential operator of Fuchs type modulo perturbation. Set

$$
\begin{aligned}
Q= & \prod_{j=0}^{m-1}\left(-D_{\zeta} \zeta-\frac{j}{\lambda+1}\right) \\
+ & \sum_{l=1}^{m} \sum_{q=0}^{\lambda l} \frac{(\lambda+1)^{\frac{q-\lambda l}{\lambda+1}}}{q!} \sigma_{0}\left(\frac{\partial^{q} P_{l}}{\partial x_{1}^{q}}\left(0, x^{\prime}, \xi^{\prime}\right) \xi_{n}^{-\frac{s+l}{\lambda+1}}\right) \\
& \quad \times\left(-i D_{\zeta}\right)^{\frac{s+1}{\lambda+1}} \prod_{j=0}^{m-l-1}\left(-D_{\zeta} \zeta-\frac{j}{\lambda+1}\right) .
\end{aligned}
$$

## Proposition 1.

$Q$ is of Fuchs type in $\zeta$ if we freeze the parameters $\left(x^{\prime}, \xi^{\prime}\right)$. Its regular singular points are $\infty$ and $i \alpha_{j}\left(0, x^{\prime}, \xi^{\prime}\right) \xi_{n}^{-1}(1 \leq j \leq m)$. The characteristic exponents at $\infty$ are $1,1+\frac{1}{\lambda+1}, \ldots, 1+\frac{m-1}{\lambda+1}$ and it is a non-logarithmic singularity. $0,1,2,3, \ldots, m-2$ are characteristic exponents at $i \alpha_{j}\left(0, x^{\prime}, \xi^{\prime}\right) \xi_{n}^{-1}$. If the remaining characteristic exponent is a non-integer, then this is a non-logarithmic singularity.

## PROOF.

The coefficient of $D_{\zeta}^{m}$ is

$$
\begin{aligned}
& (-\zeta)^{m}+\sum_{l=1}^{m} \frac{1}{(\lambda l)!} \sigma_{0}\left(\frac{\partial^{\lambda l} P_{l}}{\partial x_{1}^{\lambda l}}\left(0, x^{\prime}, \xi^{\prime}\right) \xi_{n}^{-l}\right)(-i)^{l}(-\zeta)^{m-l} \\
= & (-1)^{m}\left\{\zeta^{m}+\sum_{l=1}^{m} \frac{1}{(\lambda l)!} \sigma_{l}\left(\frac{\partial^{\lambda l} P_{l}}{\partial x_{1}^{\lambda l}}\left(0, x^{\prime}, \xi^{\prime}\right)\right)\left(\frac{i}{\xi_{n}}\right)^{l} \zeta^{m-l}\right\} .
\end{aligned}
$$

We want to determine its zeroes. Recall that

$$
\begin{aligned}
& \sigma_{m}\left(\xi_{1}^{m}+P_{1}\left(x, \xi^{\prime}\right) \dot{\xi}_{1}^{m-1}+\cdots+P_{m}\left(x, \xi^{\prime}\right)\right) \\
= & \left(\xi_{1}-x_{1}^{\lambda} \alpha_{1}\left(x, \xi^{\prime}\right)\right) \ldots\left(\xi_{1}-x_{1}^{\lambda} \alpha_{m}\left(x, \xi^{\prime}\right)\right) .
\end{aligned}
$$

Hence

$$
\begin{aligned}
& \left(x_{1}^{\lambda} \tau\right)^{m}+\sigma_{1}\left(P_{1}\right)\left(x, \xi^{\prime}\right)\left(x_{1}^{\lambda} \tau\right)^{m-1}+\cdots+\sigma_{m}\left(P_{m}\right)\left(x, \xi^{\prime}\right) \\
= & \left(x_{1}^{\lambda} \tau-x_{1}^{\lambda} \alpha_{1}\right) \ldots\left(x_{1}^{\lambda} \tau-x_{1}^{\lambda} \alpha_{m}\right) .
\end{aligned}
$$

and $\sigma_{l}\left(P_{l}\right)\left(x, \xi^{\prime}\right)$ is divisible by $\left(x_{1}^{\lambda}\right)^{l}$ (Levi condition). Comparing the coefficients of $\left(x_{1}^{\lambda}\right)^{m}$ in the above equality, we obtain

$$
\begin{aligned}
& \tau^{m}+\sum_{l=1}^{m} \frac{1}{(\lambda l)!} \frac{\partial^{\lambda l}}{\partial x_{1}^{\lambda l}} \sigma_{l}\left(P_{l}\right)\left(0, x^{\prime}, \xi^{\prime}\right) \tau^{m-l} \\
& =\left(\tau-\alpha_{1}\left(0, x^{\prime}, \xi^{\prime}\right)\right) \ldots\left(\tau-\alpha_{m}\left(0, x^{\prime}, \xi^{\prime}\right)\right)
\end{aligned}
$$

Set $\tau=\frac{\xi_{n}}{i} \zeta$, then

$$
\begin{aligned}
& \left(\frac{\xi_{n}}{i} \zeta\right)^{m}+\sum_{l=1}^{m} \frac{\partial^{\lambda l}}{\partial x_{1}^{\lambda l}} \sigma_{l}\left(P_{l}\right)\left(0, x^{\prime}, \xi^{\prime}\right)\left(\frac{\xi_{n}}{i} \zeta\right)^{m-l} \\
= & \left(\frac{\xi_{n}}{i} \zeta-\alpha_{1}\left(0, x^{\prime}, \xi^{\prime}\right)\right) \ldots\left(\frac{\xi_{n}}{i} \zeta-\alpha_{m}\left(0, x^{\prime}, \xi^{\prime}\right)\right) .
\end{aligned}
$$

Multiplication by $\left(\frac{i}{\xi_{n}}\right)^{m}$ yields

$$
\begin{aligned}
& \zeta^{m}+\sum_{l=1}^{m} \frac{1}{(\lambda l)!} \frac{\partial^{\lambda l}}{\partial x_{1}^{\lambda l}} \sigma_{l}\left(P_{l}\right)\left(0, x^{\prime}, \xi^{\prime}\right)\left(\frac{i}{\xi_{n}}\right)^{l} \zeta^{m-l} \\
= & \left(\zeta-\frac{i}{\xi_{n}} \alpha_{1}\left(0, x^{\prime}, \xi^{\prime}\right)\right) \ldots\left(\zeta-\frac{i}{\xi_{n}} \alpha_{m}\left(0, x^{\prime}, \xi^{\prime}\right)\right) .
\end{aligned}
$$

Hence we know the location of the singularities of $Q$. At $i \alpha_{j} \xi_{n}^{-1}, Q$ has the form

$$
Q=(\text { nonzero function }) \times\left(\left(\zeta-i \alpha_{j} \xi_{n}^{-1}\right) D_{\zeta}^{m}+\ldots\right) .
$$

So the assertion about $i \alpha_{j} \xi_{n}^{-1}$ is obvious.
To study $Q$ at $\infty$, we need the following lemma.

## Lemma 2.

Set

$$
L=\prod_{j=0}^{m-1}\left(-\zeta D_{\zeta}-\frac{j+\lambda+1}{\lambda+1}\right)+\sum_{l=1}^{m} \sum_{q} a_{l q} D_{\zeta}^{\frac{s+1}{\lambda+1}} \prod_{j=0}^{m-l-1}\left(-\zeta D_{\zeta}-\frac{j+\lambda+1}{\lambda+1}\right) .
$$

Here $\sum_{q}$ is taken with respect to $q=0,1, \ldots, \lambda l, q \equiv \lambda l \bmod \lambda+1$. Then $\infty$ is an regular singular point of $Q$ and its characteristic exponents are $1,1+\frac{1}{\lambda+1}, \ldots, 1+\frac{m-1}{\lambda+1}$.
PROOF. Set $\zeta^{-\frac{1}{\lambda+1}}=x, \theta=x D_{x}$. Then $\zeta^{-1}=x^{\lambda+1}, \zeta D_{\zeta}=-\frac{1}{\lambda+1} \theta$ and $D_{\zeta}=$ $-\frac{1}{\lambda+1} x^{\lambda+1} \theta$. We have
$L=\prod_{j=0}^{m-1}\left(\frac{1}{\lambda+1} \theta-\frac{j+\lambda+1}{\lambda+1}\right)+\sum_{l=1}^{m} \sum_{q} a_{l q}\left(-\frac{1}{\lambda+1} x^{\lambda+1} \theta\right)^{\frac{2+1}{\lambda+1}} \prod_{j=0}^{m-l-1}\left(\frac{1}{\lambda+1} \theta-\frac{j+\lambda+1}{\lambda+1}\right)$.

By choosing suitable constants $a_{l q}^{\prime}$,

$$
M:=(\lambda+1)^{m} L=\prod_{j=0}^{m-1}\{\theta-(j+\lambda+1)\}+\sum_{l=1}^{m} \sum_{q} a_{l q}^{l}\left(x^{\lambda+1} \theta\right)^{\frac{2+1}{\lambda+1}} \prod_{j=0}^{m-l-1}\{\theta-(j+\lambda+1)\} .
$$

Then this lemma follows from the one below.

## Lemma 3.

$$
\begin{aligned}
& \quad x^{-(\lambda+1)} M(x, D) x^{\lambda+1} \\
& = \\
& =x^{m} f(x)\left\{D_{x}^{m}+M_{1}(x) D_{x}^{m-1}+\cdots+M_{m}(x)\right\}
\end{aligned}
$$

where $f, M_{1}, \ldots, M_{m}$ are holomorphic functions defined in a neighborhood of $x=0$ and $f(0) \neq 0$.
PROOF. By using $x^{-1} \theta x=\theta+1$ repeatedly, we obtain

$$
x^{-(\lambda+1)} \theta x^{\lambda+1}=\theta+\lambda+1 \text {. }
$$

So we have

$$
\begin{aligned}
& x^{-(\lambda+1)} M(x, D) x^{\lambda+1} \\
= & \prod_{j=0}^{m-1}(\theta-j)+\sum_{l=1}^{m} \sum_{q} a_{l q}^{\prime}\left\{x^{\lambda+1}(\theta+\lambda+1)\right\}^{\frac{2+1}{x+1}} \prod_{j=0}^{m-l-1}(\theta-j) \\
= & x^{m} D_{x}^{m}+\sum_{l=1}^{m} \sum_{q} a_{l q}^{\prime}\left\{x^{\lambda+1}(\theta+\lambda+1)\right\}^{\frac{s+l}{x+1}} x^{m-l} D_{x}^{m-l} .
\end{aligned}
$$

Obviously, the coefficient of $D_{x}^{m}$ has the form $x^{m}+\mathcal{O}\left(x^{m+1}\right)$. The proof is finished as soon as we prove

$$
x^{-(\lambda+1)} M x^{\lambda+1} \in x^{m} \mathcal{D}_{x=0} .
$$

We have only to prove that

$$
\left\{x^{\lambda+1}(\theta+\lambda+1)\right\}^{\frac{3+i}{x+1}} x^{m-l} D_{x}^{m-l} \in x^{m} \mathcal{D}_{x=0}
$$

This inclusion follows from the sublemma below.
Sublemma 4. $\left\{x^{\lambda+1}(\theta+\lambda+1)\right\}^{a} x^{b} \in x^{(\lambda+1) a+b} \mathcal{D}_{x=0}$.
PROOF. Induction on $a$. The case $a=0$ is obvious. We have

$$
\begin{aligned}
& \left\{x^{\lambda+1}(\theta+\lambda+1)\right\} x^{(\lambda+1) a+b} \\
= & x^{\lambda+1}\left[x^{(\lambda+1) a+b} \theta+\{(\lambda+1) a+b+\lambda+1\} x^{(\lambda+1) a+b}\right] \\
\in & x^{(\lambda+1)(a+1)+b} \mathcal{D}_{x=0}
\end{aligned}
$$

and induction proceeds. We have finally proved Proposition 1.
In PARTS 1 and 2, the following lemma will be convenient.

## Lemma 5.

$$
L: \zeta^{-1} \mathbb{C}\left\{\zeta^{-\frac{1}{\lambda+1}}\right\} \rightarrow \zeta^{-\left(\frac{m}{\lambda+1}+1\right)} \mathbb{C}\left\{\zeta^{-\frac{1}{\lambda+1}}\right\}
$$

is surjective. Here $\mathbb{C}\{\cdot\}$ is the set of convergent power series.
PROOF. We have the following commutative diagram:

$$
\begin{array}{ccc}
\mathcal{O}_{x=0} & \xrightarrow{x^{-(\lambda+1)} M x^{\lambda+1}} & x^{m} \mathcal{O}_{x=0} \\
{ }^{\lambda+1} \downarrow \\
x^{\lambda+1} \downarrow \\
x^{\lambda+1} \mathcal{O}_{x=0} & \xrightarrow{M} & \downarrow^{\lambda+1}
\end{array}
$$

We continue the explanation of how to construct a $j$-null solution. Assume that we are given a microdifferential operator $E_{j}\left(\zeta, x^{\prime}, D^{\prime}\right)$ satisfying the following conditions:

$$
L\left(\zeta, x^{\prime}, \partial_{\zeta}, D_{x}^{\prime}\right) E_{j}\left(\zeta, x^{\prime}, D^{\prime}\right)=0
$$

and that $E_{j}$ is defined in $\{\operatorname{Re} \zeta>0\} \times\left(\right.$ a conic neighborhood $\subset i T^{*} N$ of $\left.p^{\prime}\right)$ where $p^{\prime}=$ $\rho(p), \rho: N \times i T^{*} M \rightarrow i T^{*} N, N=\left\{x_{1}=0\right\} \subset M$. $E_{j}$ is regular at $\zeta=i \alpha_{j}\left(0, x^{\prime}, \xi^{\prime}\right) \xi_{n}^{-1}$. At $\zeta=\infty, \stackrel{M}{E}$ j has the form

$$
E_{j}=\sum_{k=0}^{\infty} \zeta^{-1-\frac{k}{x+1}} E_{j k}\left(x^{\prime}, D^{\prime}\right)
$$

$E_{j}\left(\zeta, x^{\prime}, D^{\prime}\right) f\left(x^{\prime}\right) \in \mathcal{C O}_{+}^{\infty}$ for any $f\left(x^{\prime}\right) \in \mathcal{C}_{N, p^{\prime}}$.
Then according to [Kat 2], $E_{j} f$ defines a j-null solution. Moreover, there exists a nonzero constant $C_{k}$ such that

$$
D_{1}^{k}\left(E_{j} f\right)\left(+0, x^{\prime}\right)=C_{k} D_{n}^{\frac{k}{\lambda+1}} E_{j k}\left(x^{\prime}, D^{\prime}\right) f\left(x^{\prime}\right) .
$$

So we have a better understanding of $j$-null solutions from the viewpoint of boundary value problems.

In certain cases, this approach is really powerfull. In PARTS 1 and 2, we will present a very detailed analysis of some second and third order operators.

Next, we give a result about a particlular class of m -th order hyperbolic operator.
Let us consider

$$
P(x, D)=D_{1}^{m}+\sum_{l=1}^{m} \sum_{q} b_{q} x_{1}^{q} D_{1}^{m-l} D_{n}^{\frac{q+l}{\lambda+1}} .
$$

where the second summation is taken with respect to $q=0,1,2, \ldots, \lambda l, q \equiv \lambda l \bmod \lambda+1$. Then

$$
\sigma(P)=\xi_{1}^{m}+\sum_{l=1}^{m} b_{\lambda l, l} x_{1}^{\lambda l} \xi_{1}^{m-l} \xi_{n}^{l}
$$

Moreover, in this case $Q$ is an ordinary differential operator of Fuchs type without parameter $\left(x^{\prime}, \xi^{\prime}\right)$. That is, $Q$ has the same form as $L$ in Lemma 2. We assume that

$$
\sigma(P)(x, \xi)=\left(\xi_{1}-x_{1}^{\lambda} a_{1} \xi_{n}\right) \ldots\left(\xi_{1}-x_{1}^{\lambda} a_{m} \xi_{n}\right), a_{j} \in \mathbb{R}, a_{j} \neq a_{k}(j \neq k) .
$$

Set $\alpha_{j}\left(x, \xi^{\prime}\right)=a_{j} x_{1}^{\lambda} \xi_{n}$. According to [K-K] and [Kat 1bis], we have the isomorphism

$$
\text { b.v. : } \begin{aligned}
\left(\Gamma_{\left\{x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p} & \rightarrow \stackrel{m}{\oplus} \mathcal{C}_{N, p^{\prime}} \\
u & \mapsto\left(D_{1}^{k} u\left(+0, x^{\prime}\right)\right)_{k=0}^{m-1}
\end{aligned}
$$

Here $p^{\prime}=\rho(p)$ and $\rho$ is the projection $N \times i T^{*} M \rightarrow i T^{*} N, N=\left\{x_{1}=0\right\} \subset M$. Assume that for each $j$, there is a characteristic esponent $\notin \mathbb{Z}$ at $\zeta=i \alpha_{j} \xi_{n}^{-1}=i a_{j}$. Our result is

Theorem 6. (i) There is an isomorphism

$$
N_{j}^{+}: \stackrel{m-1}{\oplus} \mathcal{C}_{N, p^{\prime}} \rightarrow N u l l(j,+)
$$

(ii) The image of $\operatorname{Null}(j,+) \subset\left(\Gamma_{\left\{x_{1}>0\right\}} C_{M}^{P}\right)_{p}$ under b.v. is characterized by a relationship written in terms of microdifferential operators of fractional order.
PROOF. There exists a solution $v_{c}^{j}(\zeta)$ to $Q v_{c}^{j}=0$ in the right half plane which is not holomorphic if $c=m$ and holomorphic if $c=1,3, \ldots, m-1$. We may assume that $v_{c}^{j}$ 's are linearly independent. According to [Kat 2], $v_{\mathrm{c}}^{j}(\zeta) f\left(x^{\prime}\right), f \in \mathcal{C}_{N, p^{\prime}}$ defines an element of $\left(\Gamma_{\left\{x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p}$. Thus we have constructed

$$
N_{j}^{+, c}: \mathcal{C}_{N, p^{\prime}} \rightarrow\left(\Gamma_{\left\{x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p}
$$

Since $N_{j}^{+, c} f$ is j -null if $1 \leq c \leq m-1$, we can define

$$
\begin{aligned}
N_{j}^{+}: \quad{ }^{m-1} \mathcal{C}_{N, p^{\prime}} & \rightarrow N u l l(j,+) \\
\left(f_{1}, \ldots, f_{m-1}\right) & \mapsto \sum_{c=1}^{m-1}\left(N_{j}^{+, c} f\right)(x) .
\end{aligned}
$$

On the other hand, we can define

$$
\begin{aligned}
\tilde{N}_{j}^{+}: \quad{ }_{\oplus}^{m} \mathcal{C}_{N, p^{\prime}} & \rightarrow\left(\Gamma_{\left\{x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p} \\
\left(f_{1}, \ldots, f_{m-1}, f_{m}\right) & \mapsto \sum_{c=1}^{m}\left(N_{j}^{+, c} f\right)(x)
\end{aligned}
$$

Obviously $\tilde{N}_{j}^{+}=\left(N_{j}^{+}, N_{j}^{+, m}\right)$. We define $B_{j}^{+}$with the following commutative diagram


We show that $\tilde{B}_{j}^{+}$is an isomorphism. Let

$$
v_{c}^{j}(\zeta)=\sum_{k=0}^{\infty} v_{c, k}^{j} \zeta^{-1-\frac{k}{x+1}} .
$$

Then $\tilde{B}_{j}^{+}$is represented by

$$
\tilde{B}_{j}^{+}=\operatorname{diag}\left(C_{0}, C_{1} D_{n}^{\frac{1}{\lambda}+1}, C_{2} D_{n}^{\frac{\lambda}{2}+1}, \ldots, C_{m-1} D_{n}^{\frac{m-1}{\lambda+1}}\right)\left(\begin{array}{ccc}
v_{1,0}^{j} & \cdots & v_{m, 0}^{j} \\
v_{1,1}^{j} & \cdots & v_{m, 1}^{j} \\
\vdots & & \vdots \\
v_{1, m-1}^{j} & \cdots & v_{m, m-1}^{j}
\end{array}\right)
$$

The second matrix is invertible because it comes from $m$ linearly independent solutions. The first one is obviously invertible. So $\tilde{B}_{j}^{+}$is an isomorphism.

By using the commutativity of the diagram, we see that $\tilde{N}_{j}^{+}$is an isomorphism. Now let us prove (i). $N_{j}^{+}$is obviously injective. Surjectivity follows because $\left(N_{j}^{+, m} f_{m}\right)(x) \neq 0$ on $b_{j}^{+}$if $f_{m} \neq 0$ and $N_{j}^{+}$is surjective. Next, let us prove (ii). We denote by $B_{j}^{+}$the restriction of $\tilde{B}_{j}^{+}$on ${ }^{m-1} \mathcal{C}_{N, p^{\prime}}=\left\{\left(f_{1}, \ldots, f_{m-1}, 0\right) \in \stackrel{m}{\oplus} \mathcal{C}_{N, p^{\prime}}\right\}$. We have

$$
B_{j}^{+}=\operatorname{diag}\left(C_{0}, C_{1} D_{n}^{\frac{1}{\lambda+1}}, C_{2} D_{n}^{\frac{2}{\lambda+1}}, \ldots, C_{m-1} D_{n}^{\frac{m-1}{\lambda+1}}\right)\left(\begin{array}{ccc}
v_{1,0}^{j} & \ldots & v_{m-1,0}^{j} \\
v_{1,1}^{j} & \ldots & v_{m, 1}^{j} \\
\vdots & & \vdots \\
v_{1, m-1}^{j} & \ldots & v_{m-1, m-1}^{j}
\end{array}\right)
$$

are rank of the second matrix is $m-1$. Since the components of the two matrices are mmutative, we can use the same argument a in the usual linear albebra.
F: $\quad y$, we introduce a notion which will be important in PARTS 1 and 2.
Set

$$
\operatorname{Sol}(j, \pm):=\left\{u \in\left(\Gamma_{\left\{ \pm x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p} ; u=0 \text { on } b_{k}(k \neq j)\right\}
$$

An element of it is called a $j$-pure solution (in $\pm x_{1}>0$ ). In other words, a solution is j -pure if and only if it is k -null for all $k \neq j$. Obviously, a null solution is a sum of pure solutions. The study of pure solutions is more difficult than that of null solutions.

## PART 1 SECOND ORDER CASE

## §1 statement of the theorems

Let

$$
P(x, D)=D_{1}^{2}-\frac{1}{i}\left(\beta_{1}+\beta_{2}\right) x_{1} D_{1} D_{n}-\beta_{1} \beta_{2} x_{1}^{2} D_{n}^{2}-\frac{2}{i} \gamma D_{n}+\sum_{l=0}^{\text {finite }} \alpha_{-l}\left(x_{1}^{2}, x^{\prime}, D^{\prime}\right) x_{1}^{l} D_{1}^{l}
$$

be a microdifferential operator defined in a neighborhood of $p \in\left\{(x, i \xi) \in i T^{*} M ; x_{1}=\right.$ $\left.0, \xi_{n}>0\right\}$ such that ord $\alpha_{-l} \leq-l-1$ and $\alpha_{-l}$ is a polynomial in $t=\frac{1}{2} x_{1}^{2}$ and $x_{n}$. Here we write $x=\left(x_{1}, \ldots, x_{n}\right)=\left(x_{1}, x^{\prime}\right) \in \mathbb{R}^{n}=M$. We also assume that $\beta_{1}$ and $\beta_{2}$ are purely imaginary constants with $\frac{\beta_{1}}{i}>\frac{\beta_{2}}{i}$. The principal symbol of P , denoted by $\sigma(P)=\left(\xi_{1}-\frac{\beta_{1}}{i} x_{1} \xi_{n}\right)\left(\xi_{1}-\frac{\beta_{2}}{i} x_{1} \xi_{n}\right) . P$ is microhyperbolic and doubly characteristic over the initial surface $N=\left\{x ; x_{1}=0\right\}$. $\operatorname{Char}(P)$, the (purely imaginary) characteristic variety, is the union of two hypersurfaces $\xi_{1}= \pm \frac{\beta_{j}}{i} x_{1} \xi_{n}(j=1,2)$, which have an non-involutory intersection $\left\{x_{1}=\xi_{1}=0\right\} \ni p$. Let $b_{j}$ be the bicharacteristic strip of $\left\{\xi_{1}-\frac{\beta_{j}}{i} x_{1} \xi_{n}=0\right\}$ issuing from $p$, and $b_{j}^{ \pm}$be its intersection with $\left\{(x ; i \xi d x) ; \pm x_{1}>0\right\}$. Since $P$ has simple characteristics in $x_{1} \neq 0$, we can apply the propagation theorem in [SKK]. That is, if a microfunction $u$ satisfies $P u=0$ in $\pm x_{1}>0, b_{j}^{ \pm} \subset$ supp $u$ or $b_{j}^{ \pm} \cap \operatorname{supp} u=\phi$. Moreover, the general theory on microhyperbolic operators due to $[\mathrm{KK}]$ implies that we have the commutative diagram:

where $\mathcal{C}_{M}^{P}$ is the kernel sheaf of $P$, the horizontal arrow is the restriction, and the vertical arrows are the initial and the boundary value morphisms. Set

$$
\operatorname{Sol}(j, \pm)=\left\{u \in\left(\Gamma_{\left\{ \pm x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p} ; u=0 \text { on } b_{k}^{ \pm}(k \neq j)\right\}
$$

An element of $\operatorname{Sol}(j, \pm)$ is called a $j$-pure solution. Assume

$$
(*): \quad c=\frac{\frac{3}{2} \beta_{1}-2 \beta_{2}+\gamma}{\beta_{1}-\beta_{2}} \notin \frac{1}{2} \mathbb{Z}=\left\{0, \pm \frac{1}{2}, \pm 1, \pm \frac{3}{2}, \ldots\right\} .
$$

Then we have the following three theorems.
Theorem A. (boundary value problem with purity)
The map

$$
\begin{gathered}
\operatorname{Sol}(j, \pm) \rightarrow \mathcal{C}_{N, p} \\
u \mapsto u\left(+0, x^{\prime}\right)
\end{gathered}
$$

is an isomorphism. Moreover, if $\alpha_{-l}=0$ for all $l,\left(^{*}\right)$ can be replaced by a weaker
condition

$$
(*)^{\prime}: \quad c \notin \frac{3}{2}+\mathbb{N} \cup 2-\mathbb{N}, \quad \mathbb{N}=\{1,2,3, \ldots\} .
$$

Theorem B. (characterization of j-pure solutions by a relationship between their boundary values)

There exists a microdifferential operator $P_{j}^{ \pm}\left(x^{\prime}, D^{\prime}\right)$ of order $\frac{1}{2}-\mathbb{N}_{0}$ with the following property. (Here $\mathbb{N}_{0}$ is the set of non-negative integers.) : An element $u$ of $\left(\Gamma_{\left\{ \pm x_{1}>40\right\}} \mathcal{C}_{M}^{P}\right)_{p}$ is $j$-pure if and only if

$$
D_{1} u\left( \pm 0, x^{\prime}\right)=P_{j}^{ \pm}\left(x^{\prime}, D^{\prime}\right)\left[u\left( \pm 0, x^{\prime}\right)\right]
$$

Moreover if $\alpha_{-l}=0$ for all $l$, then ( ${ }^{*}$ ) can be replaced by (*)'.
Theorem C. (branching of singularities)
Let $u(x)$ be an element of $\mathcal{C}_{M, p}^{P}$. If $u$ is pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{1}^{\mp} \cup b_{2}^{\mp}$ is contained in suppu. Moreover, if $\alpha_{-l}=0$ for all $l$, we can consider the following two cases not included in ( ${ }^{*}$ ).
(i) $c \in \frac{5}{2}-\mathbb{N}$

If $u$ is 1-pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{1}^{\mp} \cup b_{2}^{\mp}$ is contained in supp $u$. If $u$ is 2-pure in $\pm x_{1}>0$, then, it is 2-pure also in $\mp x_{1}>0$.
(ii) $c \in 1+\mathbb{N}$

If $u$ is 1-pure in $\pm x_{1}>0$, then it is 1-pure also in $\mp x_{1}>0$. If $u$ is 2-pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{1}^{\mp} \cup b_{2}^{\mp}$ is contained in suppu.

We can treat another kind of perturbation. The constant $\gamma$ can be replaced by a microdifferential operator. Let the coordinate of $p^{\prime}$ be $\left(\dot{x}_{2}, \ldots, \dot{x}_{n} ; i \dot{\xi}^{\prime} d x^{\prime}\right)$ and $\tilde{\gamma}=$ $\tilde{\gamma}\left(x^{\prime}, D^{\prime}\right)$ be a microdifferential operator of order $\leq 0$ defined near $p^{\prime}$. $\tilde{\gamma}$ has an expansion of the form

$$
\begin{gathered}
\tilde{\gamma}\left(x^{\prime}, D^{\prime}\right)=\sum_{j=0}^{\infty} \gamma_{j}\left(x^{\prime \prime}, D^{\prime}\right)\left(x_{n}-\dot{x}_{n}\right)^{j} \\
x^{\prime \prime}=\left(x_{2}, \ldots, x_{n-1}\right) .
\end{gathered}
$$

Let $\hat{\gamma}=\hat{\gamma}(x, D)$ be defined by

$$
\hat{\gamma}(x, D)=\sum_{j=0}^{\infty} \gamma_{j}\left(x^{\prime \prime}, D^{\prime}\right)\left(\frac{1}{2} x_{1} D_{1} D_{n}^{-1}+x_{n}-\dot{x}_{n}\right)^{j} .
$$

It is ant perator of order $\leq 0$ defined in a neighborhood of $p$. Set $\mathbb{C} \ni \gamma=\sigma_{0}(\tilde{\gamma})\left(p^{\prime}\right)=$ $\sigma_{0}\left(\gamma_{0}\right)\left(p^{\prime}\right)=\sigma_{0}(\hat{\gamma})(p)$ and $c=\frac{\frac{3}{2} \beta_{1}-2 \beta_{2}+\gamma}{\beta_{1}-\beta_{2}}$. Let us consider the operator

$$
P(x, D)=D_{1}^{2}-\frac{1}{i}\left(\beta_{1}+\beta_{2}\right) x_{1} D_{1} D_{n}-\beta_{1} \beta_{2} x_{1}^{2} D_{n}^{2}-\frac{2}{i} D_{n} \hat{\gamma}(x, D) .
$$

Purity and the related mappings are defined in the usual way. In this situation, we have the following results.

## Theorem A'.

If $c$ satisfies ( ${ }^{*}$ ), then the map

$$
\operatorname{Sol}(j, \pm) \rightarrow \mathcal{C}_{N, p}
$$

is an isomorphism.

## Theorem B'.

If c satisfies $\left({ }^{*}\right)^{\prime}$, then there exists a microdifferential operator $P_{j}^{ \pm}\left(x^{\prime}, D^{\prime}\right)$ of order $\in \frac{1}{2}-\mathbb{N}_{0}$, which has the following property: An element of $\left(\Gamma_{\left\{x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p}$ is $j$-pure if and only if

$$
D_{1} u\left( \pm 0, x^{\prime}\right)=P_{j}^{ \pm}\left[u\left( \pm 0, x^{\prime}\right)\right] .
$$

## Theorem C'.

Assume $c \notin \frac{1}{2} \mathbb{Z}$. Then we have: If $u$ is pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{1}^{\mp} \cup b_{2}^{\mp}$ is contained in suppu.

## $\S 2$ proof of the unperturbed case

We are going to construct a $\mathbb{C}$-linear mapping

$$
\begin{gathered}
E_{j}^{ \pm}: \mathcal{C}_{N, p^{\prime}} \rightarrow \operatorname{Sol}(j, \pm) \\
f\left(x^{\prime}\right) \mapsto\left(E_{j}^{ \pm} f\right)(x) .
\end{gathered}
$$

Here $p^{\prime}=\rho(p)$ and $\rho: i T^{*} M \underset{M}{\times} N \rightarrow i T^{*} N$ is the pull-back of the inclusion map $N \hookrightarrow M$.

## 2-1 construction of $E_{j}^{ \pm}$

Let us consider

$$
P(x, D)=D_{1}^{2}-\frac{1}{i}\left(\beta_{1}+\beta_{2}\right) x_{1} D_{1} D_{n}-\beta_{1} \beta_{2} x_{1}^{2} D_{n}^{2}-\frac{2}{i} \gamma D_{n},
$$

where is a complex constant. In $x_{1}>0$, we have $\operatorname{Ker} P=\operatorname{Ker} \frac{1}{4} x_{1}^{2}$. We perform the change of variables $t=\frac{1}{2} x_{1}^{2}$ in the latter operator. By using $x_{1} D_{1}=2 t D_{t}$ and $x_{1}^{2} D_{1}^{2}=x_{1} D_{1}\left(x_{1} D_{1}-1\right)$, wcobtain

$$
\left\{\begin{array}{l}
\frac{1}{4} x_{1}^{2} D_{1}^{2}=t D_{\cdot}\left(t D_{t}-\frac{1}{2}\right) \\
\frac{1}{4} x_{1}^{3} D_{1} D_{n}=t D_{n} \cdot t D_{t} \\
\frac{1}{4} x_{1}^{4} D_{n}^{2}=\left(t D_{n}\right)^{2}
\end{array}\right.
$$

Hence

$$
\frac{1}{4} x_{1}^{2} P=t D_{t}\left(t D_{t}-\frac{1}{2}\right)+i\left(\beta_{1}+\beta_{2}\right) t D_{n} t D_{t}-\beta_{1} \beta_{2}\left(t D_{n}\right)^{2}-\frac{\gamma}{i} t D_{n}
$$

Next, we apply the quantized Legendre transform $\mathcal{L}$ with respect to $\left(t, x^{\prime}\right)$. ( $\mathcal{L}$ is denoted by $\beta_{n}^{+}$in [Kat1]). $\mathcal{L}$ is a quantized contact transformation defined by

$$
\left\{\begin{array}{l}
\mathcal{L} D_{t} \mathcal{L}^{-1}=-i \zeta D_{n}, \quad \mathcal{L} D_{k} \mathcal{L}^{-1}=D_{k}(k \neq 1) \\
\mathcal{L} t \mathcal{L}^{-1}=-i D_{\zeta} D_{n}^{-1}, \mathcal{L} x_{k} \mathcal{L}^{-1}=x_{k}(k \neq 1, n) \\
\mathcal{L} x_{n} \mathcal{L}^{-1}=x_{n}+D_{\zeta} \zeta D_{n}^{-1} .
\end{array}\right.
$$

In particular, we have

$$
\left\{\begin{array}{l}
\mathcal{L} t D_{t} \mathcal{L}^{-1}=-D_{\zeta} \zeta=-\left(\zeta D_{\zeta}+1\right) \\
\mathcal{L} t D_{n} \mathcal{L}^{-1}=-i D_{\zeta}
\end{array}\right.
$$

Here $\zeta$ is the dual variable of (the complexification of ) $t$.
Then the Legendre image, denoted by $\mathcal{Q}\left(\zeta, D_{\zeta}\right)$, is

$$
\begin{aligned}
Q & =\left(-\zeta D_{\zeta}-1\right)\left(-\zeta D_{\zeta}-\frac{3}{2}\right)+i\left(\beta_{1}+\beta_{2}\right) \cdot i D_{\zeta}\left(\zeta D_{\zeta}+1\right)-\beta_{1} \beta_{2}\left(-D_{\zeta}^{2}\right)-\frac{\gamma}{i}\left(-i D_{\zeta}\right) \\
& =\left(\zeta^{2} D_{\zeta}^{2}+\frac{7}{2} \zeta D_{\zeta}+\frac{3}{2}\right)-\left(\beta_{1}+\beta_{2}\right)\left(\zeta D_{\zeta}^{2}+2 D_{\zeta}\right)+\beta_{1} \beta_{2} D_{\zeta}^{2}+\gamma D_{\zeta} \\
& =\left(\zeta-\beta_{1}\right)\left(\zeta-\beta_{2}\right) D_{\zeta}^{2}+\left[\frac{7}{2} \zeta-2\left(\beta_{1}+\beta_{2}\right)+\gamma\right] D_{\zeta}+\frac{3}{2}
\end{aligned}
$$

- $Q$ is transformed into Gauss hypergeometric operator $G=G\left(\frac{3}{2}, 1, c ; z, D\right)$ if we introduce a new independent variable $z$ by $\zeta=\left(-\beta_{1}+\beta_{2}\right) z+\beta_{1}$, where

$$
\begin{gathered}
G=G\left(\frac{3}{2}, 1, c ; z, D\right)=z(1-z) D^{2}+\left[c-\left(\frac{3}{2}+1+1\right) z\right] D-\frac{3}{2} \cdot 1, \quad D=D_{z} \\
c=\frac{\frac{3}{2} \beta_{1}-2 \beta_{2}+\gamma}{\beta_{1}-\beta_{2}} .
\end{gathered}
$$

Its Riemann scheme is

$$
\left\{\begin{array}{cccc}
0 & 1 & \infty & \\
0 & 0 & 1 & ; z \\
1-c & c-\frac{5}{2} & \frac{3}{2} &
\end{array}\right\}
$$

## Lemma1.

Let $u(z)$ be a solution to $G u=0$. If it is holomorphic both at $z=0,1$, then it vanishes identically.

PROOF. $u$ is analytically continued to the entire complex plane, and its exponent at $z=\infty$ is 1 . Apply Liouville's theorem.

We want to find a solution $v_{j}(z)(j=1,2$ respectively) in the upper half plane, not vanishing identically, such that $v_{j}$ is holomorphic at $z=1,0$ respectively. (Hence singular at $z=0,1$ respectively). Moreover its expansion coefficients at $z=\infty$ will be necessary in the next section. This is a kind of connection problem. It is solved by using well-known formulas. We quote from [IKSY]. Assume that

$$
c \notin \frac{3}{2}+\mathbb{N} \cup 2-\mathbb{N}, \mathbb{N}=\{1,2,3, \ldots\}
$$

and, in the up:er half plane, set $0<\arg z<\pi$.
By cho ing six different paths in the Euler integral representaion, we obtain six solution, $F_{1}(z), \ldots, F_{6}(z)$ that have the following properties:

$$
\begin{equation*}
F_{1}+F_{2}+F_{3}=0 \tag{1}
\end{equation*}
$$

(2)

$$
F_{1}-F_{4}+F_{5}=0
$$

(3)

$$
\epsilon(-c) F_{2}-F_{5}+F_{6}=0
$$

(4)

$$
-\epsilon(-c) F_{3}+F_{4}+F_{6}=0
$$

(5)

$$
F_{1}(z)=\frac{\Gamma\left(\frac{5}{2}-c\right) \Gamma(c-1)}{\Gamma\left(\frac{3}{2}\right)} z^{-\frac{3}{2}} F\left(\frac{3}{2}, \frac{5}{2}-c ; \frac{3}{2} ; \frac{1}{z}\right)
$$

(6)

$$
\begin{equation*}
F_{3}(z) \text { is holomorphic at } z=1 \tag{7}
\end{equation*}
$$

$$
\begin{equation*}
F_{6}(z)=2 i \epsilon\left(-\frac{1}{2} c\right) F\left(1,2-c, ; \frac{1}{2} ; \frac{1}{z}\right), \tag{8}
\end{equation*}
$$

where $\epsilon(\cdot)=\exp (2 \pi i \cdot)$ and $F$ is the Gauss hypergeometric series. With the notation above, we define

$$
v_{1}(z)=2 \epsilon(-c) F_{3}(z), \quad v_{2}(z)=-2 \epsilon(-c) F_{2}(z)
$$

Let us calculy r expansion coefficients at $z=\infty$. From (1), ..., (4), we obtain

$$
(1-c(-c)) F_{1}-2 \epsilon(-c) F_{3}+2 F_{6}=0
$$

and

$$
(1+\epsilon(-c)) F_{1}+2 \epsilon(-c) F_{2}+2 F_{6}=0 .
$$

Hence

$$
\begin{aligned}
& v_{1}(z)=(1-\epsilon(-c)) F_{1}(z)+2 F_{6}(z) \\
& v_{2}(z)=(1+\epsilon(-c)) F_{1}(z)+2 F_{6}(z) .
\end{aligned}
$$

When we expand $v_{j}(z)$ into the form

$$
\begin{gather*}
v_{j}(z)=\sum_{n=0}^{\infty} v_{j,-1-\frac{n}{2}} z^{-1-\frac{n}{2}} \text { at } z=\infty, \operatorname{Im} z>0  \tag{9}\\
(0<\arg z<\pi)
\end{gather*}
$$

we see easily that

$$
\left(\begin{array}{cc}
v_{1,-1} & v_{2,-1}  \tag{10}\\
v_{1,-\frac{3}{2}} & v_{2,-\frac{3}{2}}
\end{array}\right)=\left(\begin{array}{cc}
4 i \epsilon\left(-\frac{1}{2} c\right) & 0 \\
0 & \frac{\Gamma\left(\frac{5}{2}-c\right) \Gamma(c-1)}{\Gamma\left(\frac{3}{2}\right)}
\end{array}\right)\left(\begin{array}{cc}
1 & 1 \\
1-\epsilon(-c) & 1+\epsilon(-c)
\end{array}\right)
$$

Now we come back to the $\zeta$-plane. Since there is a correspondence

$$
\begin{gathered}
\zeta=\beta_{1}, \beta_{2}, \infty \longleftrightarrow z=0,1, \infty \\
\operatorname{Re} \zeta>0 \longleftrightarrow \operatorname{Im} z>0 \quad(0<\arg z<\pi),
\end{gathered}
$$

$v_{j}(z)(j=1,2)$ defines a solution to $Q$ in the right half plane $\subset \mathbb{C}_{\zeta}$ which is singular at $\beta_{j}$ and holomorphic at $\beta_{3-j}$. We denote it by $V_{j}(\zeta)$. Let us calculate its expansion coefficients at $\zeta=\infty, \operatorname{Re} \zeta>0 . z=\frac{\zeta-\beta_{1}}{-\beta_{1}+\beta_{2}}$ leads to

$$
z^{-1}=\frac{-\beta_{1}+\beta_{2}}{\zeta} \sum_{k=0}^{\infty}\left(1+\frac{\beta_{1}}{\zeta}\right)^{k}
$$

and

$$
z^{-\frac{3}{2}}=\left(-\beta_{1}+\beta_{2}\right)^{\frac{3}{2}} \zeta^{-\frac{3}{2}}\left(1+\frac{3}{2} \frac{\beta_{1}}{\zeta}+O\left(\zeta^{-2}\right)\right)
$$

Here

$$
-\frac{\pi}{2}<\arg \zeta<\frac{\pi}{2}, \quad \arg \left(-\beta_{1}+\beta_{2}\right)=-\pi / 2
$$

When we expand $V_{j}(\zeta)$ into the form

$$
V_{j}(\zeta)=\sum_{n=0}^{\infty} V_{j,-1-\frac{n}{2}} \zeta^{-1-\frac{n}{2}} \text { at } \zeta=\infty, \operatorname{Re} \zeta>0 \quad\left(-\frac{\pi}{2}<\arg \zeta<\frac{\pi}{2}\right)
$$

we have

$$
\binom{V_{j,-1}}{V_{j,-\frac{3}{2}}}=\left(\begin{array}{cc}
-\beta_{1}+\beta_{2} & 0 \\
0 & \left(-\beta_{1}+\beta_{2}\right)^{\frac{3}{2}}
\end{array}\right)\binom{v_{j,-1}}{v_{j,-\frac{3}{2}}}
$$

Combining this and (10), we obtain

$$
\begin{align*}
& \left(\begin{array}{cc}
V_{1,-1} & V_{2,-1} \\
V_{1,-\frac{3}{2}} & V_{2,-\frac{3}{2}}
\end{array}\right) \\
& =\left(\begin{array}{cc}
4 i \epsilon\left(-\frac{1}{2}\right)\left(-\beta_{1}+\beta_{2}\right) & 0 \\
0 & \frac{\Gamma\left(\frac{5}{2}-c\right) \Gamma\left(c-\frac{1}{2}\right)}{\Gamma\left(\frac{3}{2}\right)}\left(-\beta_{1}+\beta_{2}\right)^{\frac{3}{2}}
\end{array}\right)\left(\begin{array}{cc}
1 & 1 \\
1-\epsilon(-c) & 1+\epsilon(-c)
\end{array}\right) \tag{11}
\end{align*}
$$

Let us define $E_{j}^{+}$. According to [Kat2], $V_{j}(\zeta) f\left(x^{\prime}\right) \in \mathcal{C O}_{+}^{\infty}$ gives an element of Sol $(\mathrm{j},+)$, which we denote by $\left(E_{j}^{+} f\right)(x)$. To define $E_{j}^{-}$, we change the sign of the time variable by introducing $\bar{x}_{1}=-x_{1}$. Since $D_{1}=-D_{\bar{x}_{1}}$, we have

$$
P(x, D)=P\left(-\bar{x}_{1}, x^{\prime},-D_{\bar{x}_{1}}, D^{\prime}\right)=P\left(\bar{x}_{1}, x^{\prime}, D_{\bar{x}_{1}}, D^{\prime}\right)
$$

That is, $P$ does not change its form. We can apply the same argument as above and the definition of $E_{j}^{-}$is obvious. Note that $\xi_{1}-\beta_{j} x_{1} \xi_{n}=-\left(\bar{\xi}_{1}-\beta_{j} \bar{x}_{1} \xi_{n}\right)$, where $\bar{\xi}_{1}$ is the dual of $\bar{x}_{1}$.

## 2-2 boundary values of $E_{j}^{ \pm}$

Let $f\left(x^{\prime}\right), g\left(x^{\prime}\right) \in \mathcal{C}_{N, p^{\prime}}$. According to [Kat2], we have

$$
\begin{equation*}
\binom{\left(E_{j}^{+} f\right)\left(+0, x^{\prime}\right)}{\left(D_{1} E_{j}^{+} f\right)\left(+0, x^{\prime}\right)}=\binom{2 \pi V_{j,-1} f\left(x^{\prime}\right)}{2 \sqrt{2 \pi}\left(\frac{D_{\mu}}{i}\right)^{\frac{1}{2}} V_{j,-\frac{3}{2}} f\left(x^{\prime}\right)}(j=1,2) \tag{12}
\end{equation*}
$$

Denote by $L^{+}$the morphism

$$
\begin{gathered}
L^{+}: \stackrel{2}{\oplus} \mathcal{C}_{N, p^{\prime}} \rightarrow \stackrel{2}{\oplus} \mathcal{C}_{N, p^{\prime}} \\
\binom{f}{g} \mapsto\binom{\left(E_{1}^{+} f+E_{2}^{+} g\right)\left(+0, x^{\prime}\right)}{D_{1}\left(E_{1}^{+} f+E_{2}^{+} g\right)\left(+0, x^{\prime}\right)} .
\end{gathered}
$$

Combination of (11) and (12) yield

$$
\begin{gathered}
L^{+}\binom{f}{g}=A\left(\begin{array}{cc}
1 & 1 \\
1-\epsilon(-c) & 1+\epsilon(-c)
\end{array}\right)\binom{f}{g} \\
A=\left(\begin{array}{cc}
2 \pi & 0 \\
0 & 2 \sqrt{2 \pi}\left(\frac{D_{n}}{i}\right)^{\frac{1}{2}}
\end{array}\right)\left(\begin{array}{cc}
4 i \epsilon\left(-\frac{1}{2}\right)\left(-\beta_{1}+\beta_{2}\right) & 0 \\
0 & \frac{\Gamma\left(\frac{5}{2}-c\right) \Gamma\left(c-\frac{1}{2}\right)}{\Gamma\left(\frac{3}{2}\right)}\left(-\beta_{1}+\beta_{2}\right)^{\frac{3}{2}}
\end{array}\right) .
\end{gathered}
$$

We may forget the explicit form of $A$. All we'll need is the fact that $A=\operatorname{diag}\left(A_{1}, A_{2} D_{n}^{\frac{1}{2}}\right)$ where $A_{1}$ and $A_{2}$ are nonzero constants. In particular, $L^{+}$is an isomorphism. Next, denote by $L^{-}$the morphism

$$
\begin{gathered}
L^{-}:{\stackrel{2}{\oplus} \mathcal{C}_{N, p^{\prime}} \rightarrow \stackrel{2}{\oplus} \mathcal{C}_{N, p^{\prime}}}^{\binom{f}{g} \mapsto\binom{\left(E_{1}^{-} f+E_{2}^{-} g\right)\left(-0, x^{\prime}\right)}{D_{1}\left(E_{1}^{-} f+E_{2}^{-} g\right)\left(-0, x^{\prime}\right)} .}
\end{gathered}
$$

Obviously,

$$
\binom{E_{1}^{-} f+E_{2}^{-} g}{\left.D_{\bar{x}_{1}}\left(E_{1}^{-} f+E_{2}^{-} g\right)\left(\bar{x}_{1}, x^{\prime}\right)\right|_{\bar{x}_{1} \rightarrow+0}}
$$

is represented by the same matrix as $L^{+}$. Since $D_{\bar{x}_{1}}=-D_{1}, L^{-}$has a slightly different representation:

$$
L^{-}=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right) L^{+}=A\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)\left(\begin{array}{cc}
1 & 1 \\
1-\epsilon(-c) & 1+\epsilon(-c)
\end{array}\right) .
$$

$L^{-}$is an isomorphism, of course.

## $2-3$ end of the proofs (of the unperturbed case)

We have the following commutative diagram:
(CD1)

$$
\begin{aligned}
&\left(\Gamma_{\left\{ \pm x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p}=\oplus_{j=1}^{2} \operatorname{Sol}(j, \pm) \stackrel{E^{ \pm}}{ } \stackrel{2}{\oplus} \mathcal{C}_{N, p^{\prime}} \\
& \text { b.v.l } \downarrow \\
& 2\left\lfloor 1 L^{ \pm}\right. \\
& \stackrel{\oplus}{\oplus} \mathcal{C}_{N, p^{\prime}}=\stackrel{2}{\oplus} \mathcal{C}_{N, p^{\prime}}
\end{aligned}
$$

where the horizontal arrow is the map

$$
E^{ \pm}=\oplus_{j=1}^{2} E_{j}^{ \pm}:^{t}\left(f\left(x^{\prime}\right), g\left(x^{\prime}\right)\right) \mapsto\left(E_{1}^{ \pm} f\right)(x)+\left(E_{2}^{ \pm} g\right)(x) .
$$

The first vertical arrow is

$$
u(x) \mapsto{ }^{t}\left(u\left( \pm 0, x^{\prime}\right), D_{1} u\left( \pm 0, x^{\prime}\right)\right)
$$

and it is known to be an isomorphism. ([K-K], [Kat1]). Therefore $E^{ \pm}$is an isomorphism. Since $E^{ \pm}=\oplus_{j=1}^{2} E_{j}^{ \pm}$, each $E_{j}^{ \pm}$is an isomorphism. In this way, we have arrived at the important identification:

$$
\begin{aligned}
& E_{1}^{ \pm}: \mathcal{C}_{N, p^{\prime}} \oplus 0 \stackrel{\sim}{\sim} \operatorname{Sol}(1, \pm) \\
& E_{2}^{ \pm}: 0 \oplus \mathcal{C}_{N, p^{\prime}} \stackrel{\sim}{\rightarrow} \operatorname{Sol}(2, \pm) .
\end{aligned}
$$

From (CD1), we obtain the following commutative diagram:

where the first horizontal arrow is the identification above, the left vertical arrow is

$$
u \mapsto u\left( \pm 0, x^{\prime}\right),
$$

and the second vertical arrow is $f\left(x^{\prime}\right) \mapsto A_{1} f\left(x^{\prime}\right)$. This implies the latter part of Theorem A.

Next, we prove Theorem B. We want to characterize the image of $\operatorname{Sol}(j, \pm)$ under b.v. Because of $(\mathrm{CD} 1)$, it is $L^{ \pm}\left(\mathcal{C}_{N, p^{\prime}} \oplus 0\right)$ if $j=1$ and $L^{ \pm}\left(0 \oplus \mathcal{C}_{N, p^{\prime}}\right)$ if $j=2$. Here

$$
L^{ \pm}\binom{f}{0}=\left(\begin{array}{cc}
1 & 0 \\
0 & \pm 1
\end{array}\right) A\binom{1}{1-\epsilon(-c)} f
$$

and

$$
L^{ \pm}\binom{0}{g}=\left(\begin{array}{cc}
1 & 0 \\
0 & \pm 1
\end{array}\right) A\binom{1}{1+\epsilon(-c)} g .
$$

Theorem B follows immediately.
Finally, let us prove Theorem C. We have the isomorphisms below:

where $E^{ \pm}$is the direct sum of the identification maps $E_{j}^{ \pm}, j=1,2$. We set $B=$ $\left(L^{-}\right)^{-1} L^{+}$and call it the branching matrix. It is easy to see that

$$
\begin{gathered}
B=\left(\begin{array}{cc}
1 & 1 \\
1-\epsilon(-c) & 1+\epsilon(-c)
\end{array}\right)^{-1}\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)\left(\begin{array}{cc}
1 & 1 \\
1-\epsilon(-c) & 1+\epsilon(-c)
\end{array}\right) \\
=\left(\begin{array}{cc}
\epsilon(c) & 1+\epsilon(c) \\
1-\epsilon(c) & \epsilon(c)
\end{array}\right) .
\end{gathered}
$$

The identification above enables us to reduce the problem of branching to the study of the branching matrix $B$. We have only to know when a certain component of $B$ is (not) zero. The proof of Theorem C is now complete.

A nonzero constant is an elliptic microdifferential operator of order 0 . Even if it is perturbed in the lower order terms, it remains elliptic. This observation will be important in the following section.

## $\S 3$ proof of the perturbed case

In this section we assume that

$$
c \notin \frac{1}{2} \mathbb{Z}=\left\{0, \pm \frac{1}{2}, \pm 1, \pm \frac{3}{2} \ldots\right\} .
$$

We only explain the construction of $E_{2}^{+}$. The remaining three maps are constructed in the same way.

## 3-1 right inverse

We make some preparation for the symbol calculus in the next subsection. $G\left(z, D_{z}\right)$ is an ordinary differential operator of Fuchs type with three regular singular points $\zeta=0,1, \infty$. Its Riemann scheme is

$$
\left(\begin{array}{cccc}
0 & 1 & \infty & \\
0 & 0 & 1 & z \\
1-c & c-\frac{5}{2} & \frac{3}{2} &
\end{array}\right)
$$

and no logarithmic term appears. The exponent of the Wronskian is $-c$ at $\zeta=0$ and $c-\frac{7}{2}$ at $\zeta=1$. Let $\Omega \subset \mathbb{C}_{z}$ be a domain as in the figure.


$$
\text { Figure } 1
$$

$G(z, D)$ induces a linear mapping:

$$
G: \mathcal{O}(\Omega) \rightarrow \mathcal{O}(\Omega)
$$

We are going to construct a right inverse $G^{-1}$ by using the method of variation of parameters. Let $F_{1}, F_{2}$ be the series solutions of exponent $0,1-c$ respectively defined near $z=0$. Let $W$ be their Wronskian.

$$
W=\operatorname{det}\left(\begin{array}{ll}
F_{1} & F_{2} \\
F_{1}^{\prime} & F_{2}^{\prime}
\end{array}\right) .
$$

It is easy to see that

$$
\begin{aligned}
f(z) \rightarrow\left(G^{-1} f\right)(z)= & -F_{1}(z) \int_{0}^{z} \frac{F_{2}(y)}{y(1-y) W(y)} f(y) d y \\
& +F_{2}(z) \int_{0}^{z} \frac{F_{1}(y)}{y(1-y) W(y)} f(y) d y
\end{aligned}
$$

gives a right inverse of $G$. Here the integrals are taken in the sense of Riemann-Liouville. We want to obtain some estimate on the integral operator $G^{-1}$. We say that a function $f$ has exponent $(p, q)$ at $z=a$ if $f$ has the form

$$
f(z)=z^{p} f_{1}(z)+z^{q} f_{2}(z), \quad(p-q \notin \mathbb{Z})
$$

where $f_{1}$ and $f_{2}$ are holomorphic at $z=a$ and $f_{1}(a), f_{2}(a) \neq 0$. Set, for $\delta, 0<\delta \ll 1$,

$$
K_{\delta}=\{z \in \Omega ; \operatorname{dist}(z, \partial \Omega) \leq \delta\}
$$

## Proposition 2.

There exist positive constants $\tilde{c}$ and $C$, not depending on $\delta$, such that for all $f \in \mathcal{O}(\Omega)$,
have we have

$$
\sup _{K_{\delta}}\left|G^{-1} f\right| \leq C \delta_{K_{\delta}}^{-\bar{c}_{\sup }}|f|
$$

In the proof, we see that $\tilde{c}=\left[\left|\operatorname{Rec}-\frac{5}{2}\right|\right]+1$, where $[a]$ is the smallest integer not exceeding $a$.

PROOF. We consider the second term in the definition of $G^{-1}$. (The first term is easier to deal with.) Let us introduce the following notation:

$$
(J f)(z)=F_{2}(z) \int_{0}^{z} \cdot \frac{F_{1}(y)}{y(1-y) W(y)} f(y) d y .
$$

We will deduce an estimate on $J$ in several steps.

Lemma 3. Fix a sufficiently small constant $R>0$. Then, there exists a constant $C_{1}>0$, independent of $\delta$, such that

$$
\sup _{|z| \leq \frac{R}{2}}\left|z^{1-c} \int_{0}^{z} y^{1-c} g(y) d y\right| \leq C_{K_{\delta}} \sup _{K_{\delta}}|g(y)|
$$

for all $g(z) \in \mathcal{O}(\Omega)$.
PROOF. We may assume that $\{|z| \leq 2 R\} \cdot \subset$ int $K_{\delta}$ for any $\delta, 0<\delta \ll 1$. So $g(z)$ has a Taylor expansion

$$
g(z)=\sum_{n=0}^{\infty} g_{n} z^{n} \quad \text { in } \quad\{|z|<2 R\}
$$

Since Riemann-Liouville integral can be carried out term by term ([IKSY]), we have

$$
\begin{gather*}
\int_{0}^{z} y^{c-1} g(y) d y=z^{c} \sum_{n=0}^{\infty} \frac{1}{c+n} g_{n} z^{n}, \\
z^{1-c} \int_{0}^{z} y^{c-1} g(y) d y=\sum_{n=0}^{\infty} \frac{1}{c+n} g_{n} z^{n+1}, \tag{13}
\end{gather*}
$$

By the way, the assumption $c \notin \frac{1}{2} \mathbb{Z}$ implies that there exists a constant $C_{c}$ such that

$$
\left|\frac{1}{c+n}\right| \leq C_{c} \text { for all } n=0,1,2, \ldots
$$

Moreover, Cauchy's estimate shows that

$$
\left|g_{n}\right| \leq \frac{1}{R^{n}} \sup _{|y|=R}|g(y)| .
$$

Therefore, in view of (13), we obtain, in $|z|<R$,

$$
\begin{aligned}
\left|z^{1-c} \int_{0}^{z} z^{c-1} g(y) d y\right| & \leq \sum_{n=o}^{\infty} C_{c} \cdot \frac{1}{R^{n}} \sup _{|y|=R}|g(y)| \cdot|z|^{n+1} \\
& =C_{c} \frac{|z|}{1-\frac{|z|}{R}|y|=R} \sup |g(y)| .
\end{aligned}
$$

This leads to the lemma because we have

$$
\sup _{|z| \leq \frac{R}{2}} \frac{|z|}{1-\frac{|z|}{R}}=R .
$$

## Lemma 4.

There exists a constant $C_{2}>0$, independent of $\delta$, such that

$$
\sup _{K_{6}}\left|z^{1-c} \int_{0}^{z} y^{c-1} g(y) d y\right| \leq C_{2} \sup _{K_{6}}|g(y)|
$$

for all $g(z) \in \mathcal{O}(\Omega)$.
PROOF. What remains is the estimate for $z \in\left\{|z|>\frac{R}{2}\right\} \cap K_{\delta}$. We write the function in question as the sum of two terms.

$$
z^{1-c} \int_{0}^{z}=z^{1-c} \cdot \int_{0}^{\frac{R x_{x}}{\| \eta}}+z^{1-c} \int_{\left.\frac{R}{2 \lambda} \right\rvert\,}^{z}
$$

We can apply Lemma 3 to the first term. In fact,

$$
z^{1-c} \int_{0}^{\frac{R z}{2 \lambda \pi}}=\left(\frac{2|z|}{R}\right)^{1-c} \cdot\left(\frac{R z}{2|z|}\right)^{1-c} \int_{0}^{\frac{R z}{2 \pi \mid}},
$$

where the first factor is bounded in $\Omega \cap\{|z|>R / 2\}$ and the second factor is estimated by using Lemma 1. Let us consider the second term. We may assume that the length of the path of integration $\subset K_{\delta}$ from $\frac{R z}{2|z|}$ to $z$ is estimated by a constant $C_{R, \Omega}^{1}>0$ independent of $\delta$. Additionally, in $\left\{|z|>\frac{R}{2}\right\} \cap K_{\delta}, z^{1-c}$ is estimated by a constant $C_{R, \Omega}^{2}>0$, independent of $\delta$. Therefore we have

$$
\left|z^{1-c} \int_{\frac{R_{2}}{2 \nabla \nabla}}^{z}\right| \leq C_{R, \Omega}^{1} C_{R, \Omega}^{2} \sup _{K_{s}}|g(y)|
$$

in $\left\{|z|>\frac{R}{2}\right\} \cap K_{\delta}$.
Lemma 5. Put $c^{\prime}=\left|\operatorname{Rec}-\frac{5}{2}\right| \geq 0$. There exists a constant $C_{J}>0$, independent of $\delta$, such that for all $f \in \mathcal{O}(\Omega)$, we have

$$
\sup _{K_{\delta}}|J f| \leq C_{J} \delta^{-c^{\prime}} \sup _{K_{\delta}}|f| .
$$

## PROOF.

We have

$$
\frac{F_{1}(y)}{y(1-y) W(y)}=y^{c-1} G(y)
$$

where $G(y)$ is holomorphic in $\Omega$, or more precisely, in the universal covering space of $\mathbb{C} \backslash\{1\}$, and has exponent $\left(-c+\frac{5}{2}, 0\right)$ at $y=1$. Obviously,

$$
\frac{F_{1}(y)}{y(1-y) W(y)} f(y)=y^{c-1} \times G(y) f(y)
$$

and we consider $G f$ as $g$ in Lemma 2. Since

$$
\sup _{K_{\delta}}|G f| \leq \sup _{K_{\delta}}|G| \sup _{K_{\delta}}|f| \leq C_{G} \delta^{\min \left(-\operatorname{Rec}+\frac{5}{2}, 0\right)} \sup _{K_{\delta}}|f|,
$$

Lemma 4 implies that

$$
\begin{equation*}
\sup _{K_{6}}\left|z^{1-c} \int_{0}^{z} y^{c-1} G(y) f(y) d y\right| \leq C_{2} C_{G} \delta^{\min \left(-\operatorname{Rec}+\frac{5}{2}, 0\right)} \sup _{K_{\delta}}|f| \tag{14}
\end{equation*}
$$

On the other hand, $F_{2}(z) / z^{1-c}$ is holomorphic in $\Omega$ and has exponent $\left(0, c-\frac{5}{2}\right)$ at $z=1$. So, there exists a constant $C_{3}>0$ independent of $\delta$ such that

$$
\begin{equation*}
\sup _{K_{b}}\left|\frac{F_{2}(z)}{z^{1-c}}\right| \leq C_{3} \delta^{\min \left(\operatorname{Rec}-\frac{s}{2}, 0\right)} . \tag{15}
\end{equation*}
$$

Combination of (14) and (15) yields the lemma, because

$$
\min \left(-\operatorname{Rec}+\frac{5}{2}, 0\right)+\min \left(\operatorname{Rec}-\frac{5}{2}, 0\right)=\min \left( \pm\left(\operatorname{Rec}-\frac{5}{2}\right)\right)=-c^{\prime} .
$$

PROOF OF PROPOSITION 2 CONTINED. The first term in the definition of $G^{-1}$ satisfies the same estimate as Lemma 5, with a larger $C$, if necessary. Then the proposition follows immediately, because $0 \geq-c^{\prime}>-\tilde{c}$.

## 3-2 successive approximation

Let us consider

$$
\begin{aligned}
P(x, D)= & D_{1}^{2}-\frac{1}{i}\left(\beta_{1}+\beta_{2}\right) x_{1} D_{1} D_{n}-\beta_{1} \beta_{2} x_{1}^{2} D_{n}^{2}-\frac{2}{i} \gamma D_{n} \\
& +\sum_{l=0}^{\text {finite }} \alpha_{-l}\left(x_{1}^{2}, x^{\prime}, D^{\prime}\right) x_{1}^{l} D_{1}^{l} .
\end{aligned}
$$

As in $\S 2$, we put $t=\frac{1}{2} x_{1}^{2}$ in $\frac{1}{4} x_{1}^{2} P(x, D)$ and use the quantized Legendre transform $\mathcal{L}$. Let us calculate the contribution of the perturbation term

$$
P^{\prime}(x, D)=\sum_{l=0}^{\text {finite }} \alpha_{-l}\left(x_{1}^{2}, x^{\prime}, D^{\prime}\right) x_{1}^{l} D_{1}^{l} .
$$

First, we consider $x_{1}^{2} \cdot x_{1}^{l} D_{1}^{l}$. It is easy to see that

$$
x_{1}^{2} \cdot x_{1}^{l} D_{1}^{l}=2 t \cdot 2 t D_{t}\left(2 t D_{t}-1\right) \ldots\left(2 t D_{t}-l+1\right) .
$$

Lemma 6. Let $W(\mathbb{C})$ be the Weyl algebra of variable $t$ and $V$ be the subalgebra generated by $t$ and $\vartheta=t D_{t}$. Then we have $t^{j} V \subset V t^{j} \subset W(\mathbb{C}) t^{j} .(j=0,1,2, \ldots)$
PROOF. Obviously $[t, \vartheta]=-t$, so $t \vartheta \in V t$. Hence the case $j=1$ is proved. The remaining cases are proved by induction.

This lemma $(j=1)$ implies that

$$
\frac{1}{4} x_{1}^{2} \cdot x_{1}^{l} D_{1}^{l} \in \mathcal{D}_{t, w^{\prime}} t
$$

Therefore $\frac{1}{4} x_{1}^{2} P^{\prime}(x, D)$ belongs to $\mathcal{E}_{t, x^{\prime}} t \cap \mathcal{E}_{t, x^{\prime}}(-1)$ and is a polynomial in $t$ and $x_{n}$ by the assumption on $\alpha_{-l}$. Its image under $\mathcal{L}$, denoted by $Q^{\prime}\left(\zeta, x^{\prime}, D_{\zeta}, D^{\prime}\right)$ belongs to $\mathcal{E}_{\zeta, x^{\prime}} D_{\zeta} \cap \mathcal{E}_{\zeta, x^{\prime}}(-1)$ and is a polynomial in $D_{\zeta}$ and $\zeta$. More precisely, it has the form

$$
Q^{\prime}\left(\zeta, x^{\prime}, D_{\zeta}, D^{\prime}\right)=\sum_{m=1}^{\text {finite }} \sum_{j=0}^{m-1} \tilde{\alpha}_{m, j}\left(x^{\prime}, D^{\prime}\right) \zeta^{j} D_{\zeta}^{m} \in \mathcal{E}(-1) .
$$

where ord $\tilde{\alpha}_{m, j} \leq-m-1$. If we write it in terms of the other complex variable $z, Q^{\prime}$ is transformed into

$$
G^{\prime}\left(z, x^{\prime}, D_{z}, D^{\prime}\right)=\sum_{m=1}^{\dot{m}} \sum_{j=0}^{m-1} \alpha_{m, j}\left(x^{\prime}, D^{\prime}\right) z^{j} D_{z}^{m} \in \mathcal{E}(-1) .
$$

Here $\bar{m}$ is a positive integer and $\alpha_{m, j}$ is a microdifferential operator defined in a neighborhood of $p^{\prime}=\left(x^{\prime} ; i \xi^{\prime} d x^{\prime}\right)=\rho(p) \in i T^{*} N, N=\mathbb{R}^{n-1}$. Thus $-\frac{1}{4} x_{1}^{2} P$ is transformed into $G-G^{\prime}$, where $G$ is the Gauss hypergeometric operator in $\S 2$. We will construct a microdifferential operator $E\left(z, x^{\prime}, D^{\prime}\right)$ of order 0 that satisfies

$$
\left(G-G^{\prime}\right)\left(z, x^{\prime}, \partial_{z}, D^{\prime}\right) E=0, \quad \partial_{z}=\left[D_{z}, \cdot\right]
$$

In addition, we require that $E$ should be defined in (a neighborhood $\subset \mathbb{C}_{z}$ of $\{z ; \operatorname{Im} z \geq$ $0, z \neq 1\}) \times\left(\right.$ a conic neighborhood $\subset i T^{*} N$ of $\left.p^{\prime}=\rho(p)\right)$ and that

$$
E \in z^{-1} \mathcal{E}(0)+z^{-3 / 2} \mathcal{E}(0) \quad \text { at } \quad z=\infty
$$

where $\mathcal{E}(0)$ is regarded as a sheaf on $T^{*}\left(\mathbb{P}^{1} \times \mathbb{C}^{n-1}\right)$. There is another requirement to be explained in 3-3. Put

$$
\left\{\begin{array}{l}
E_{0}(z)=v_{2}(z) \quad \text { where } v_{2} \text { is defined in } 2-1, \\
E_{k+1}\left(z, x^{\prime}, D^{\prime}\right)=G^{-1}\left[G^{\prime}\left(z, x^{\prime}, \partial_{z}, D^{\prime}\right) E_{k}\left(z, x^{\prime}, D^{\prime}\right)\right]
\end{array}\right.
$$

Here $G^{-1}$ and $G^{\prime}$ are mappings on $\mathcal{O}_{z} \otimes_{\mathrm{C}} \mathcal{E}_{x^{\prime}}$, to which $E_{k}$ belongs. We want to show that $E=\sum_{k \geq 0} E_{k}$ converges in $\mathcal{E}_{x^{\prime}} \mathcal{O}_{z}$. We have only to prove it when $z$ belongs to a fixed $\Omega$, where $\Omega$ is as in 3-1. Obviously

$$
\begin{aligned}
E_{k}\left(z, x^{\prime}, D^{\prime}\right) & =\left(G^{-1} G^{\prime}\right)^{k} E_{0}(z) \\
& =\sum_{\left(m_{k}, \ldots, m_{1}\right)} \sum_{\left(j_{k}, \ldots, j_{1}\right)}\left(G^{-1} \alpha_{m_{k}, j_{k}} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} \alpha_{m_{1}, j_{1}} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z) .
\end{aligned}
$$

where $\left(m_{k}, \ldots, m_{1}\right)$ runs through the set $\{1, \ldots, \bar{m}\} \times \ldots\{1, \ldots, \bar{m}\} \quad(k$ times $)$ and $\left(j_{k}, \ldots, j_{1}\right)$ through $\left\{0, \ldots, m_{k}-1\right\} \times \cdots \times\left\{0, \ldots, m_{1}-1\right\}$. Therefore

$$
\begin{aligned}
E_{k}\left(z, x^{\prime}, D^{\prime}\right) & =\sum_{\left(m_{k}, \ldots, m_{1}\right)\left(j_{k}, \ldots, j_{1}\right)} \alpha_{m_{k}, j_{k} \ldots \alpha_{m_{1}, j_{1}} \otimes\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}} \\
& \in \mathcal{E}_{x^{\prime}}\left(-\left(m_{k}+\cdots+m_{1}\right)-k\right) \otimes \mathrm{C} \mathcal{O}_{z} \subset \mathcal{E}_{z, x^{\prime}}(0)
\end{aligned}
$$

We will show the convergence of $\sum E_{k}$ for $z \in \Omega$ in three steps. They are:
STEP 1 estimate of $\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z)$
STEP 2 estimate of $\alpha_{m_{k}, j_{k}} \ldots \alpha_{m_{1}, j_{1}}$
STEP 3 convergence of $\sum E_{k}$ STEP1

Proposition 7. With the notation of 3-1, there exists a constant $C^{\prime \prime}$ independent of $\delta$, such that

$$
\begin{gathered}
\sup _{K_{k}}\left|\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z)\right| \\
\leq C^{\prime k+1}\left\{\tilde{c}+\left(m_{k}+\cdots+m_{1}\right)\right\}!\delta^{-(k+1) \bar{c}-\left(m_{k}+\cdots+m_{1}\right)}
\end{gathered}
$$

for all $k \geq 0$ and all $\delta, 0<\delta \ll 1$, where $\tilde{c}$ is the one in Proposition 2. We refer to this inequality as $(*)_{k, \delta}$.

PROOF. It is true for $k=0$. We proceed by induction on $k$. Assume that $(*)_{k, \delta}$ is true for all sufficiently small $\delta$. Take $\delta^{\prime}=\left(1+\frac{1}{\bar{c}+\left(m_{k}+\ldots m_{1}\right)}\right)^{-1}<\delta .(*)_{k, \delta^{\prime}}$, which is true by assumption, states that

$$
\begin{gathered}
\sup _{K_{夕^{\prime}}}\left|\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z)\right| \\
\leq C^{\prime k+1}\left\{\tilde{c}+\left(m_{k}+\cdots+m_{1}\right)\right\}!\delta^{-(k+1) \tilde{c}-\left(m_{k}+\cdots+m_{1}\right)} \\
\quad \times\left(1+\frac{1}{\tilde{c}+\left(m_{k}+\ldots m_{1}\right)}\right)^{(k+1) \tilde{c}+\left(m_{k}+\cdots+m_{1}\right)} .
\end{gathered}
$$

Here

$$
\begin{aligned}
\text { the last factor } & \leq\left(1+\frac{1}{\tilde{c}+\left(m_{k}+\ldots m_{1}\right)}\right)^{\tilde{\tilde{c}}+\left(m_{k}+\ldots m_{1}\right)}\left(1+\frac{1}{\tilde{c}+\left(m_{k}+\ldots m_{1}\right)}\right)^{k \bar{c}} \\
& \leq e\left\{\left(1+\frac{1}{k}\right)^{k}\right\}^{\bar{c}} \\
& \leq e^{\overline{\tilde{c}+1}}
\end{aligned}
$$

because $m_{k}, \ldots, m_{1} \geq 1$. Next, we emply Cauchy's estimate. A circle with center in $K_{\delta}$ and radius $\delta /\left\{\tilde{c}+\left(m_{k}+\cdots+m_{1}\right)+1\right\}$ is contained in $K_{\delta^{\prime}}$. Therefore

$$
\sup _{K_{\delta}}\left|\partial_{z}^{m_{k+1}}\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z)\right|
$$

$$
\begin{aligned}
\leq & m_{k+1}!\left\{\tilde{c}+\left(m_{k}+\cdots+m_{1}\right)+1\right\}^{m_{k+1}} \delta^{-m_{k+1}} \\
& \times C^{\prime k+1}\left\{\tilde{c}+\left(m_{k}+\cdots+m_{1}\right)\right\}!\delta^{-(k+1) \tilde{c}-\left(m_{k}+\cdots+m_{1}\right)} \times e^{\bar{c}+1} \\
\leq & m_{k+1}!e^{\bar{c}+1} C^{\prime k+1}\left\{\tilde{c}+\left(m_{k+1}+m_{k}+\cdots+m_{1}\right)\right\}! \\
& \times \delta^{-(k+1) \tilde{c}-\left(m_{k+1}+m_{k}+\cdots+m_{1}\right)}
\end{aligned}
$$

Here remark that $m_{k+1}!\leq \bar{m}$ (independent of $k$ ). By the way $\left|z^{j_{k+1}}\right|$ is bounded by a positive constant $C^{\prime \prime}$ independent of $k$. Then we finish the proof by choosing $C^{\prime}>C \cdot C^{\prime \prime} \cdot \bar{m}!e^{\bar{c}+1}$, where $C$ is the constant in Proposition 2.

## Propositon 8.

There is a constant $C_{\delta}$ such that

$$
\sup _{K_{\delta}}\left|\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z)\right| \leq\left(m_{k}+\cdots+m_{1}\right)!C_{\delta}^{k+1} .
$$

PROOF. First, we have

$$
\delta^{-\left(m_{k}+\cdots+m_{1}\right)} \leq\left(\delta^{-m_{n}}\right)^{k} .
$$

Secondly, since there is a constant $C_{\tilde{c}}>1$ such that

$$
\begin{aligned}
\frac{(\tilde{c}+l)!}{l!} & =\text { a polynomial in } l \text { of degree } \tilde{c} \\
& \leq C_{\tilde{c}}^{l+1}
\end{aligned}
$$

for any positive integer $l$, we have

$$
\frac{\left\{\tilde{c}+\left(m_{k}+\cdots+m_{1}\right)\right\}!}{\left(m_{k}+\ldots m_{1}\right)!} \leq C_{\tilde{c}}^{\left(m_{k}+\cdots+m_{1}\right)+1} \leq\left(C_{\tilde{c}}^{\tilde{m}_{1}}\right)^{k} C_{\tilde{c}} .
$$

Thus the present proposition follows from the preceding one.
A holomorphic function $f(z)$ in $\Omega$ can be regarded as a microdifferential operator in $\left(z, x^{\prime}\right)$, and its formal norm $N_{0}^{K_{\delta}}(f ; T)$ is defined. Here $T$ is an indeterminate.
Proposition 9.

$$
\begin{gathered}
N_{0}^{K_{\delta}}\left(\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z) ; T\right) \\
\quad \leq 2\left(m_{k}+\cdots+m_{1}\right)!C_{\delta / 2}^{k+1} \frac{1}{1-\frac{2 T}{\delta}}
\end{gathered}
$$

PROOF. Use Cauchy's estimate. The path of integration should by centered in $K_{\delta}$ and with radius $\delta / 2$.

## STEP2

First, we prepare some generalities.

Lemma 10. Let $P(x, D)$ be a microdifferential operator of order $\leq-m<0$ defined in a neighborhood of a compact set $\omega \subset T^{*} \mathbb{C}_{x}^{n}$, where $m$ is a positive integer. Then we have

$$
N_{0}^{\omega}(P ; T) \ll \frac{(2 n)^{-m}}{m!} T^{2 m} N_{-m}^{\omega}(P ; T)
$$

PROOF. By definition,

$$
N_{0}^{\omega}(P ; T)=\sum_{k, \alpha, \beta} \frac{2(2 n)^{-k} k!}{(|\alpha|+k)!(|\beta|+k)!} \sup _{\omega}\left|D_{x}^{\alpha} D_{\xi}^{\beta} p_{-k}(x, \xi)\right| T^{2 k+|\alpha+\beta|}
$$

where $P=\sum_{k \geq 0} p_{-k}$ and $p_{-k}$ is the homogeneous part of degree $-k$. There's no contribution by the terms corresponding to $k=0,1,2, \ldots, m-1$. Hence, if we put $l=k-m$,

$$
\begin{aligned}
N_{0}^{\omega}(P ; T)= & \sum_{l \geq 0, \alpha \beta} \frac{2(2 n)^{-(l+m)}(l+m)!}{(|\alpha|+l+m)!(|\beta|+l+m)!} \\
& \times \sup _{\omega}\left|D_{x}^{\alpha} D_{\xi}^{\beta} p_{-(l+m)}(x, \xi)\right| T^{2(l+m)+|\alpha+\beta|}
\end{aligned}
$$

We have only to prove that

$$
\frac{2(2 n)^{-(l+m)}(l+m)!}{(|\alpha|+l+m)!(|\beta|+l+m)!} \leq \frac{(2 n)^{-m}}{m!} \frac{2(2 n)^{-l} l!}{(|\alpha|+l)!(|\beta|+l)!} .
$$

This inequality is obtained by the calculation below.

$$
\begin{aligned}
& \frac{2(2 n)^{-(l+m)}(l+m)!}{(|\alpha|+l+m)!(|\beta|+l+m)!} \times \frac{(|\alpha|+l)!(|\beta|+l)!}{2(2 n)^{-l} l!} \\
& \leq(2 n)^{-m} \times \frac{1}{(|\alpha|+l+m) \ldots(|\alpha|+l+1)} \times \frac{(l+m) \ldots(l+1)}{(|\beta|+l+m) \ldots(|\beta|+l+1)} \\
& \leq(2 n)^{-m} \times \frac{1}{m!} \times 1 .
\end{aligned}
$$

Lemma 11. Let $P_{1}(x, D), \ldots, P_{k}(x, D)$ be microdifferential operators of order $\leq-m_{1}, \ldots,-m_{k}$ respectively, where $m_{1}, \ldots, m_{k}$ are positive integers. Then we have

$$
\begin{aligned}
& N_{0}^{\omega}\left(P_{k} \ldots P_{1} ; T\right) \\
& \ll \frac{(2 n)^{-\left(m_{k}+\cdots+m_{1}\right)}}{\left(m_{k}+\cdots+m_{1}\right)!} T^{2\left(m_{k}+\cdots+m_{1}\right)} N_{-m_{k}}\left(P_{k} ; T\right) \ldots N_{-m_{1}}\left(P_{1} ; T\right)
\end{aligned}
$$

PROOF. Since $\operatorname{ord}\left(P_{k} \ldots P_{1}\right) \leq-\left(m_{k}+\cdots+m_{1}\right)$, the preceding lemma implies that

$$
\begin{aligned}
& N_{0}\left(P_{k} \ldots P_{1} ; T\right) \\
& \ll \frac{(2 n)^{-\left(m_{k}+\cdots+m_{1}\right)}}{\left(m_{k}+\cdots+m_{1}\right)!} T^{2\left(m_{k}+\cdots+m_{1}\right)} N_{-\left(m_{k}+\cdots+m_{1}\right)}\left(P_{k} \ldots P_{1} ; T\right) .
\end{aligned}
$$

Moreover, according to [Bou-Kr], we have

$$
N_{-\left(m_{k}+\cdots+m_{1}\right)}\left(P_{k} \ldots P_{1} ; T\right) \ll N_{-m_{k}}\left(P_{k} ; T\right) \ldots N_{-m_{1}}\left(P_{1} ; T\right)
$$

In the lemma above, let $P_{1}, \ldots P_{k}$ be our $\alpha_{m_{1}, j_{1}}, \ldots, \alpha_{m_{k}, j_{k}}$ respectively. Regard them as operators of $n$ variables $\left(z, x^{\prime}\right)$. Then we have

$$
\begin{aligned}
& N_{0}\left(\alpha_{m_{k}, j_{k}} \ldots \alpha_{m_{1}, j_{1}} ; T\right) \\
& \ll \frac{1}{\left(m_{k}+\cdots+m_{1}+k\right)!}\left(\frac{T^{2}}{2 n}\right)^{m_{k}+\cdots+m_{1}+k} \\
& \times N_{-m_{k}-1}\left(\alpha_{m_{k}, j_{k}} ; T\right) \ldots N_{-m_{1}-1}\left(\alpha_{m_{1}, j_{1}} ; T\right) .
\end{aligned}
$$

## STEP3

Combining Proposition 9 and the estimate immediately above, we obtain

$$
\begin{aligned}
& N_{0}^{K_{\delta} \times \omega}\left(\alpha_{m_{k}, j_{k}}\left(x^{\prime}, D^{\prime}\right) \ldots \alpha_{m_{1}, j_{1}}\left(x^{\prime}, D^{\prime}\right)\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z) ; T\right) \\
& \ll \frac{2}{1-\frac{2 T}{\delta}} \frac{1}{k!} C_{\delta / 2}^{k+1} \\
& \quad \times\left\{\left(\frac{T^{2}}{2 n}\right)^{m_{k}+1} N_{-m_{k}-1}\left(\alpha_{m_{k}, j_{k}} ; T\right)\right\} \ldots\left\{\left(\frac{T^{2}}{2 n}\right)^{m_{1}+1} N_{-m_{1}-1}\left(\alpha_{m_{1}, j_{1}} ; T\right)\right\} .
\end{aligned}
$$

Here $\omega \ni p^{\prime}$ is a compact set of $T^{*} \mathbb{C}^{n-1}$ in a neighborhood of which $\alpha_{m_{k}, j_{k}}, \ldots, \alpha_{m_{1}, j_{1}}$ are defined. Since

$$
\begin{aligned}
E_{k}\left(z, x^{\prime}, D^{\prime}\right) & =\left(\sum_{m=1}^{\dot{m}} \sum_{j=0}^{m-1} G^{-1} \alpha_{m, j}\left(x^{\prime}, D^{\prime}\right) z^{j} \partial_{z}^{m}\right)^{k} E_{0}(z) \\
& =\sum_{\left(m_{k}, \ldots, m_{1}\right)} \sum_{\left(j_{k}, \ldots, j_{1}\right)} \alpha_{m_{k}, j_{k}} \ldots \alpha_{m_{1}, j_{1}}\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z) \\
& \in \mathcal{E}(0),
\end{aligned}
$$

we have

$$
\begin{aligned}
& N_{0}\left(E_{k} ; T\right) \\
& <\sum_{\left(m_{k}, \ldots, m_{1}\right)} \sum_{\left(j_{k}, \ldots, j_{1}\right)} \frac{2}{k!} \frac{1}{1-\frac{2 T}{\delta}} C_{\delta / 2}^{k+1} \\
& \quad \times\left\{\left(\frac{T^{2}}{2 n}\right)^{m_{k}+1} N_{-m_{k}-1}\left(\alpha_{m_{k}, j_{k}} ; T\right)\right\} \ldots\left\{\left(\frac{T^{2}}{2 n}\right)^{m_{1}+1} N_{-m_{1}-1}\left(\alpha_{m_{1}, j_{1} ;} ; T\right)\right\} \\
& \ll \frac{2}{k!} \frac{1}{1-\frac{2 T}{\delta}} C_{\delta / 2}^{k+1}\left[\sum_{m=1}^{\bar{m}} \sum_{j=0}^{m-1}\left(\frac{T^{2}}{2 n}\right)^{m+1} N_{-m-1}\left(\alpha_{m, j} ; T\right)\right]^{k} .
\end{aligned}
$$

Now the convergence of $E=\sum_{k \geq 0} E_{k}$ is clear. Here remark that its principal part is $E_{0}(z)=v_{2}(z)$.

Next, we have to study the behaviour of $E$ near $z=\infty$.
Lemma 12. Let $f(z), g(z)$ be holomorphic functions in the upper half plane such that $G\left(z, \partial_{z}\right) g(z)=f(z)$. Assume that in a neighborhood of $z=\infty, f$ is a finite sum of functions of exponent $2, \frac{5}{2}, 3, \frac{7}{2}, \ldots$. Then $g$ is a finite sum of functions of exponent $1, \frac{3}{2}, 2, \frac{5}{2}, 3, \ldots$.
PROOF. This is a consequence of PART 0 Lemma 2. An alternative proof is the use of the variation of parameters method. Let $F_{1}, F_{2}$ be two linearly independent homogeneous solutions and $W$ be their Wronskian. $\frac{F_{j}(y)}{y(1-y) W(y)} f(y)$ is a sum of terms of exponent $\frac{3}{2}, 2, \frac{5}{2}, 3, \ldots$ So $\int^{z} \frac{F_{j}(y)}{y(1-y) W(y)} f(y) d y$ is a sum of terms of exponent $0,1 / 2,1,3 / 2,2, \ldots$ at $z=\infty$. Note that no logarighmic term appears. Since $F_{3-j}$ is of exponent $(1,3 / 2)$, the lemma follows immediately.
Lemma 13. Let $f(z)$ be of exponent $\alpha$ at $z=\infty$. Then $z^{j} \partial_{z}^{m} f(z)$ is of exponent $\alpha+m-j$, or larger by a positive integer.
PROOF Easy.
We will use this lemma in the case $m \geq 1,0 \leq j \leq m-1$. Then the exponent increases because $m-j \geq 1$. Combining Lemmas 12 and 13 , we conclude that $\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z)$ is a finite sum of functions of exponent $1,3 / 2,2,5 / 2,3, \ldots$ at $z=\infty$. Therefore $E\left(z, x^{\prime}, D^{\prime}\right)$ can be written

$$
E\left(z, x^{\prime}, D^{\prime}\right)=z^{-1} E^{\prime}\left(z, x^{\prime}, D^{\prime}\right)+z^{-3 / 2} E^{\prime \prime}\left(z, x^{\prime}, D^{\prime}\right)
$$

where $E^{\prime}$ and $E^{\prime \prime}$ are formal microdifferential operators in a neighborhood of $z=\infty$. In fact, we have
Lemma 14. $E^{\prime}\left(z, x^{\prime}, D^{\prime}\right)$ and $E^{\prime \prime}\left(z, x^{\prime}, D^{\prime}\right)$ are microdifferential operators. (That is, they satisty a suitable growth condition.)
PROOF. E satisfies the growth condition of microdifferential operators in the universal covering space of $\{1 \ll|z|<\infty\}$. We can derive the lemma by using the sublemma below (with $\lambda=1 / 2$ ) and the lemma of Schwarz.
Sublemma. Let $D=\{z \in \mathbb{C} ; 0<|z|<r\}$ be a punctured disk, $\tilde{D}$ its universal covering and $\lambda$ a non-integer. Then the sum $\mathcal{O}(D)+z^{\lambda} \mathcal{O}(D) \subset \mathcal{O}(\tilde{D})$ is a direct sum. Moreover, if $K$ is a compact set in $D$ and $\tilde{K}$ is the closure of $\underset{0 \leq \theta \leq 2 \pi}{\cup} e^{i \theta} K \subset \tilde{D}$, then, there exists a constant $C=C_{\lambda, K}$ such that : for $f(z)=g(z)+z^{\lambda} h(z), \quad g, h \in \mathcal{O}(D)$, we have

$$
\sup _{K}|g| \leq C \sup _{\hat{K}}|f|, \quad \sup _{K}|h| \leq C \sup _{\hat{K}}|f|
$$

PROOF. Consider the variation of $f$,

$$
\begin{aligned}
\operatorname{Var} f(z) & =f\left(e^{2 \pi i} z\right)-f(z) \\
& =(1-2 \pi i \lambda) z^{\lambda} h(z)
\end{aligned}
$$

Note that $1-e^{2 \pi i \lambda} \neq 0$ by the assumption on $\lambda$. Obviously we have

$$
\sup _{K}|\operatorname{Var} f(z)| \leq \underset{\tilde{K}}{2 \sup ^{2}}|f|
$$

## 3-3 construction of $E_{j}^{ \pm}$

We are interested in the behaviour of $E\left(z, x^{\prime}, D^{\prime}\right)$ at $z=1$.
Proposition 15. $E\left(z, x^{\prime}, D^{\prime}\right) f\left(x^{\prime}\right)$ belongs to (the inverse image under $\zeta=\left(-\beta_{1}+\right.$ $\left.\beta_{2}\right) z+\beta_{1}$ of $) \mathcal{C O}_{+}^{\infty} \subset \mathcal{C}_{x^{\prime}} \mathcal{O}_{\zeta}$.
PROOF. We construct a defining function which is holomorphic in $\{\operatorname{Im} z>0\} \times$ (an infinitesimal wedge $\omega$ in $\mathbb{C}_{z^{\prime}}^{n-1}$ ). We employ the action of Bony-Schapira. We may assume that $p^{\prime}=\left(0^{\prime} ; i d x_{n}\right)$ and choose $z_{n}=i \sigma$ as the initial surface of the action. Let $F\left(z^{\prime}\right)$ be a defining function of $f\left(x^{\prime}\right)$, which is holomorphic in a flat domain $\omega \subset \mathbb{C}^{n-1}$ as in [B-S] p.107. By virtue of the flabbiness of $\mathcal{C}$ and the remark in [B-S]p.99,1.7-9, we may work in a domain where $F$ is bounded, thus satisfying the assumption of [B-S] Proposition 2.4.3. By Lemma 11, we have

$$
\begin{aligned}
& N_{0}\left(\alpha_{m_{k}, j_{k}} \ldots \alpha_{m_{1}, j_{1}} ; T\right) \\
& \lll \frac{1}{\left(m_{k}+\cdots+m_{1}+k\right)!} \\
& \quad \times\left(\frac{T^{2}}{2(n-1)}\right)^{m_{k}+\cdots+m_{1}+k} \quad N_{-m_{k}-1}\left(\alpha_{m_{k}, j_{k}} ; T\right) \ldots N_{-m_{1}-1}\left(\alpha_{m_{1}, j_{1}} ; T\right) .
\end{aligned}
$$

Because there are only a finite number of $\alpha_{m, j}$ 's, there is a constant $B>0$ such that

$$
N_{-m-1}\left(\alpha_{m, j} ; T\right)<B \text { for all } m \text { and } j \text {. }
$$

Therefore, there is a constant $A>0$ such that

$$
N_{0}\left(\alpha_{m_{k}, j_{k}} \ldots \alpha_{m_{1}, j_{1}} ; T\right) \ll \frac{A^{k+1}}{\left(k+m_{k}+\cdots+m_{1}\right)!}
$$

for any choice of $\left(m_{k}, j_{k}\right), \ldots,\left(m_{1}, j_{1}\right)$. Let us derive an estimate like the one in [B-S] p.94. We see easily that there is a constant $M_{0}$ such that the homogeneous part of degree $(-l)$ of $\alpha_{m_{k}, j_{k}} \ldots \alpha_{m_{1}, j_{1}}$ is estimated by $\frac{A^{k+1}}{\left(k+m_{k}+\cdots+m_{1}\right)!} M_{0}^{l+1} l!$. Then [BS] Proposition 2.4.3 implies that there is constant $C$ such that

$$
\left|\left(\alpha_{m_{k}, j_{k}} \ldots \alpha_{m_{1}, j_{1}}\right)_{\Sigma} F\left(z^{\prime}\right)\right| \leq C \frac{A^{k+1}}{\left(k+m_{k}+\cdots+m_{1}\right)!} d_{I}\left(z^{\prime}\right)^{-\alpha} d_{J}\left(z^{\prime}\right)^{-\beta} .
$$

in $\omega$. Combining this with Proposition 8, we obtain

$$
\sup _{z \in K_{\delta}}\left|\left(G^{-1} z^{j_{k}} \partial_{z}^{m_{k}}\right) \ldots\left(G^{-1} z^{j_{1}} \partial_{z}^{m_{1}}\right) E_{0}(z) \times\left(\alpha_{m_{k}, j_{k}} \ldots \alpha_{m_{1}, j_{1}}\right)_{\Sigma} F\left(z^{\prime}\right)\right|
$$

$$
\leq C \frac{\left(A C_{\delta}\right)^{k+1}}{k!} d_{I}\left(z^{\prime}\right)^{-\alpha} d_{J}\left(z^{\prime}\right)^{-\beta} .
$$

Set

$$
S_{k}\left(z, z^{\prime}\right) \underset{\operatorname{def}}{=}\left(\sum_{m=1}^{\dot{m}} \sum_{j=0}^{m-1} G^{-1} \alpha_{m, j}\left(z^{\prime}, D^{\prime}\right)_{\Sigma} z^{j} \partial_{z}^{m}\right)^{k} E_{0}(z) F\left(z^{\prime}\right) .
$$

Since the summation above consists of $\bar{m}(\bar{m}-1) / 2 \leq \bar{m}^{2}$ terms,

$$
\begin{aligned}
\sup _{z \in K_{\delta}}\left|S_{k}\left(z, z^{\prime}\right)\right| & \leq \sum_{\left(m_{k}, \ldots, m_{1}\right)} \sum_{\left(j_{k}, \ldots, j_{1}\right)} C \frac{\left(A C_{\delta}\right)^{k+1}}{k!} d_{I}\left(z^{\prime}\right)^{-\alpha} d_{J}\left(z^{\prime}\right)^{-\beta} \\
& \leq C A C_{\delta} d_{I}^{-\alpha} d_{J}^{-\beta} \frac{1}{k!}\left(\bar{m}^{2} A C_{\delta}\right)^{k} .
\end{aligned}
$$

This proves that $\sum_{k} S_{k}$ converges in $\{\operatorname{Im} z>0\} \times \omega$ locally uniformly. This completes the proof.

Again according to [Kat 2], $E f$, or rather its counterpart in $\zeta$-variable, defines a 2-pure solution. We denote it by $\left(E_{2}^{+} f\right)(x), x \in \mathbb{R}^{n}$. All the other $E_{j}^{ \pm}$'s are defined similarly. A special emphasis is laid on the fact that the principal part of $E$ is $E_{0}$. There's no contribution of the perturbation terms in this respect.

## 3-4 end of the proofs

In this subsection, we prove the remaining parts of Theorems $\mathrm{A}, \mathrm{B}$ and C . The mappings $L^{ \pm}, E^{ \pm}$and $B$ are defined and calculated in the same way as before. Because of the remark at the end of the preceding subsection, the principal part ( $=$ the 0 -th order part) remains the same as the unperturbed case. This preserves the ellipticity of the components of the above mappings.

## $\S 4$ proof of the case $\gamma$ is an operator <br> 4-1 substitution of operators into a convergent power series

Proposition 16. Let $S\left(w_{1}, w_{2}\right)=\sum_{j, k \geq 0} a_{j k}\left(w_{1}-\dot{w}_{1}\right)^{j}\left(w_{2}-\dot{w}_{2}\right)^{k}$ be a convergent power series, and $P=P(z, D) \in \mathcal{E}_{\mathbb{C}^{n}}(0)$ be a microdifferential operator of order $\leq 0$ defined in a neighborhood of $p \in T^{*} \mathbb{C}^{n}$. If $\sigma_{0}(P)(p)=\dot{w}_{1}$, then

$$
S\left(P, w_{2}\right)=\sum_{j, k \geq 0} a_{j k}\left(P-\dot{w}_{1}\right)^{j}\left(w_{2}-\dot{w}_{2}\right)^{k} \in \mathcal{E}_{\mathrm{C}^{n+1}}(0)
$$

is a well-defined microdifferential operator. Moreover we have

$$
S\left(P, w_{2}\right)=\sum_{j \geq 0}\left(\sum_{k \geq 0} a_{j k}\left(w_{2}-\dot{w}_{2}\right)^{k}\right)\left(P-\dot{w}_{1}\right)^{j}
$$

PROOF. We use the formal norm $N_{0}(\cdot, t)$, which we denote by $\|\cdot\|$ for brevity. We have

$$
\|S\| \ll \sum_{j, k}\left|a_{j k}\right|\left\|P-\dot{w}_{1}\right\|^{j}\left\|w_{2}-\dot{w}_{2}\right\|^{k}<\infty
$$

Remark 17. The last expression in the above proposition justifies analytic continua-
tion in the $w_{2}$-direction.

## Example 18.

The hypergeometric function $F(a, b ; c ; w)$ is a holomorphic function in $\{(c, w) ; c \neq$ $0,-1,-2, \ldots,|w|<1\}$. We can define $F(a, b ; P(z, D) ; w)$ for $P \in \mathcal{E}(0)$ if $\sigma_{0}(P)$ avoids
$0,-1,-2, \ldots$.

## Example 19. (microdifferential connection formula)

The classical connection formula for hypergeometric functions asserts that

$$
\begin{aligned}
& F\left(\frac{3}{2}, 1 ; c ; w\right) \\
= & \frac{c-1}{c-\frac{5}{2}} F\left(\frac{3}{2}, 1, \frac{7}{2}-c ; 1-w\right) \\
+ & \frac{\Gamma(c) \Gamma\left(\frac{5}{2}-c\right)}{\Gamma\left(\frac{3}{2}\right)}(1-w)^{c-\frac{5}{2}} F\left(c-\frac{3}{2}, c-1 ; c-\frac{3}{2} ; 1-w\right)
\end{aligned}
$$

If $\sigma_{0}(P) \notin \frac{1}{2} \mathbb{Z}$, we can replace $c$ by $P(z, D)$. We obtain

$$
\begin{aligned}
& F\left(\frac{3}{2}, 1 ; P ; w\right) \\
= & \frac{P-1}{P-\frac{5}{2}} F\left(\frac{3}{2}, 1, \frac{7}{2}-P ; 1-w\right) \\
& +\frac{\Gamma(P) \Gamma\left(\frac{5}{2}-P\right)}{\Gamma\left(\frac{3}{2}\right)}(1-w)^{P-\frac{5}{2}} F\left(P-\frac{3}{2}, P-1 ; P-\frac{3}{2} ; 1-w\right)
\end{aligned}
$$

In the example above, we encountered an operator of the form $w^{P(z, D)}$, which is defined by using Proposition 16. On the other hand, in [Tah] and [ O ], this kind of
operator is defined by

$$
\begin{aligned}
w^{P(z, D)} & =\exp (P(z, D) \log w) \\
& =\sum_{l \geq 0} \frac{1}{l!}\{P(z, D) \log w\}^{l} .
\end{aligned}
$$

Proposition 20. Our definition coincides with that of [Tah] and [O].
PROOF. Let $w^{P}$ be defined by

$$
w^{P(z, D)}=\sum_{j, k} a_{j k}\left(P(z, D)-\dot{w}_{1}\right)^{j}\left(w-\dot{w}_{2}\right)^{k}
$$

where $w_{2}^{w_{1}}=\sum_{j, k} a_{j k}\left(w_{1}-\dot{w}_{1}\right)^{j}\left(w_{2}-\dot{w}_{2}\right)^{k}$ is a convergent power series (in the classi-
cal sense). Set

$$
w_{1}=\left(w_{1}-\dot{w}_{1}\right)+\dot{w}_{1}, \quad \log w_{2}=\sum_{m \geq 0} b_{m}\left(w-\dot{w}_{2}\right)^{m},
$$

then

$$
\begin{aligned}
\sum_{l \geq 0} \frac{1}{l!}\left[\left\{\left(w_{1}-\dot{w}_{1}\right)+\dot{w}_{1}\right\} \sum_{m} b_{m}\left(w-\dot{w}_{2}\right)^{m}\right]^{l} & =w_{2}^{w_{1}} \\
& =\sum a_{j k}\left(w_{1}-\dot{w}_{1}\right)^{j}\left(w_{2}-\dot{w}_{2}\right)^{k}
\end{aligned}
$$

Moreover we have

$$
\begin{aligned}
& N_{0}\left(\sum_{l \geq 0} \frac{1}{l!}(P \log w)^{l}\right) \\
\ll & \sum_{l \geq 0} \frac{1}{l!}\left[\left\{N_{0}\left(P-\dot{w}_{1}\right)+N_{0}\left(\dot{w}_{1}\right)\right\} \sum_{m \geq 0}\left|b_{m}\right| N_{0}\left(w-\dot{w}_{2}\right)^{m}\right]^{l} \\
< & \infty
\end{aligned}
$$

Therefore we may rearrange the order of the sum in the same way as in the classical case and obtain

$$
\sum_{l \geq 0} \frac{1}{l!}(P \log w)^{l}=\sum_{j, k} a_{j k}\left(P-\dot{w}_{1}\right)^{j}\left(w-\dot{w}_{2}\right)^{k}
$$

Lemma 21. Let $U$ be a conic open set of $i T^{*} \mathbb{R}_{x}^{n}$ and $P(x, D)$ be a 0 -th order microdifferential operator defined there. $w^{P(x, D)}$ is defined in $\{$ Rew $>0\} \times U$. Then, for any microfunction $f(x)$ in $U, w^{P(x, D)} f(x)$ is an element of $\mathcal{C O}_{+}^{\infty}(\{$ Rew $=0\} \times U)$.
PROOF. Although this fact is well-known to specialists, there seems to be no published proof. Here we give a sketch of a proof based on the action of Bony-Schapira. We borrow some notation from them. We construct a defining function which is holomorphic in $\{\operatorname{Re} w>0\} \times\left(\right.$ an infinitesimal wedge in $\left.\mathbb{C}_{z}^{n}\right)$. We may assume that $U$ is a neighborhood of $\left(0, i d x_{n}\right)$ and we choose $z_{n}=i \sigma$ as the initial surface of the action. Let $F(z)$ be a defining function of $f(x)$. We have only to prove the convergence of

$$
\left(w^{P}\right)_{\Sigma} F=\sum_{l \geq 0} \frac{1}{l!}(\log w)^{l}\left(P(z, D)_{\Sigma}\right)^{l} F(z) .
$$

We may assume that $F$ satisfies the assumption of [B-S] Proposition 2.4.3 without loss of generality by virtue of the flabbiness of the sheaf of microfunctions and the remark in [B-S] p. 99 1.7-9. We have

$$
\left|\left(P(z, D)_{\Sigma}\right)^{t}\right| \leq C^{l} d_{I}(z)^{-\alpha} d_{J}(z)^{-\beta} .
$$

Then the convergence follows. Here the factor $1 / l!$ is essential.

## Example 22.

Put $\zeta=\left(-\beta_{1}+\beta_{2}\right) z+\beta_{1}$, for complex variables $\zeta$ and $x$. Here $\beta_{1}$ and $\beta_{2}$ are purely imaginary and $\beta_{1} / i>\beta_{2} / i$. Let $x^{\prime} \in \mathbb{R}^{n-1}$ be a real coordinate and $f\left(x^{\prime}\right)$ be a microfunction. We see easily that $F\left(3 / 2,1 ; c\left(x^{\prime}, D^{\prime}\right) ; z\right) f\left(x^{\prime}\right)$ belongs to $\mathcal{C} \mathcal{O}_{+}^{\infty} \subset \mathcal{C}_{x^{\prime}} \mathcal{O}_{\zeta}$. In fact, $\mathcal{C O}_{+}^{\infty}$ is an $\mathcal{E}$-Module and we know that $(1-z)^{c-\frac{5}{2}} f\left(x^{\prime}\right)$ belongs to $\mathcal{C O}_{+}^{\infty}$.

## 4-2 end of the proofs

We calculate in the same way as in the beggining of 2-1. Then $G$ should be replaced by

$$
\begin{gathered}
G\left(3 / 2,1 ; \tilde{c}\left(x^{\prime}, D^{\prime}\right) ; z, D\right)=z(1-z) D^{2}+\left\{\tilde{c}\left(x^{\prime}, D^{\prime}\right)-\left(\frac{3}{2}+1+1\right) z\right\} D-\frac{3}{2} \cdot 1 \\
\tilde{c}=\frac{\frac{3}{2} \beta_{1}-2 \beta_{2}+\tilde{\gamma}}{\beta_{1}-\beta_{2}}
\end{gathered}
$$

Here we have used the fact that

$$
\begin{aligned}
\tilde{\gamma}(x, D) & =\sum_{j} \gamma_{j}\left(x^{\prime \prime}, D^{\prime}\right)\left(\frac{1}{2} x_{1} D_{1} D_{n}^{-1}+x_{n}-\dot{x}_{n}\right)^{j} \\
& =\sum_{j} \gamma_{j}\left(x^{\prime \prime}, D^{\prime}\right)\left(t D_{t} D_{n}^{-1}+x_{n}-\dot{x}_{n}\right)^{j}
\end{aligned}
$$

is transformed under $\mathcal{L}$ into

$$
\begin{aligned}
& \sum_{j} \gamma_{j}\left\{\left(-\zeta D_{\zeta}-1\right) D_{n}^{-1}+x_{n}+D_{\zeta} \zeta D_{n}^{-1}-\dot{x}_{n}\right\}^{j} \\
= & \sum_{j} \gamma_{j}\left(x_{n}-\dot{x}_{n}\right)^{j}=\tilde{\gamma}\left(x^{\prime}, D^{\prime}\right)
\end{aligned}
$$

Obviously we have

$$
G\left(3 / 2,1 ; \tilde{\gamma} ; z, \partial_{z}\right) F(3 / 2,1 ; \tilde{\gamma} ; z)=0, \quad \text { etc. }
$$

Therefore we may replace $\gamma, c$ in $\S 2$ by $\tilde{\gamma}, \tilde{c} . L^{ \pm}$and $B$ (in the present context) are calculated easily. For example, we have

$$
B=\left(\begin{array}{cc}
\epsilon\left(\tilde{c}\left(x^{\prime}, D^{\prime}\right)\right) & 1+\epsilon\left(\tilde{c}\left(x^{\prime}, D^{\prime}\right)\right) \\
1-\epsilon\left(\tilde{c}\left(x^{\prime}, D^{\prime}\right)\right) & \epsilon\left(\tilde{c}\left(x^{\prime}, D^{\prime}\right)\right)
\end{array}\right)
$$

To prove Theorem C', we have to prove the ellipticity of all the components. We have

$$
\begin{gathered}
\sigma_{0}\left(\epsilon\left(\tilde{c}\left(x^{\prime}, D^{\prime}\right)\right)\right)=\epsilon\left(\sigma_{0}\left(\tilde{c}\left(x^{\prime}, D^{\prime}\right)\right)\right) \\
\sigma_{0}(\tilde{c})\left(p^{\prime}\right)=c
\end{gathered}
$$

Hence $\sigma_{0}\left(\epsilon\left(\tilde{c}\left(x^{\prime}, D^{\prime}\right)\right)\right)\left(p^{\prime}\right)=\epsilon(c) \neq 0$. The other components are dealt with in the same way. Theorems A' and B' are proved similarly.

## PART 2 THIRD ORDER CASE

## §1 the statement of the theorems

Let

$$
\begin{aligned}
P(x, D)= & D_{1}^{3}-x_{1}^{2} D_{n}^{2} D_{1}+2(a-b) D_{n} D_{1}+\{2(a+b)-3\} x_{1} D_{n}^{2} \\
& +\sum_{l=0}^{\text {finite }} \alpha_{-l}\left(x_{1}^{2}, x^{\prime}, D^{\prime}\right) x_{1}^{l+1} D_{1}^{t}
\end{aligned}
$$

be a microdifferential operator defined in a neighborhood of $p \in\left\{(x, i \xi d x) \in i T^{*} M ; x_{1}=\right.$ $\left.0, \xi_{n}>0\right\}$, such that ord $\alpha_{-l} \leq-l-1$ and that $\alpha_{-l}$ is a polynomial in $t=\frac{1}{2} x_{1}^{2}$ and $x_{n}$. Here we write $x=\left(x_{1}, \ldots, x_{n}\right)=\left(x_{1}, x^{\prime}\right) \in \mathbb{R}^{n}=M$. The principal symbol of $P$, denoted by $\sigma(P)(x, \xi)$ is factorized in the form $\sigma(P)=\left(\xi_{1}-x_{1} \xi_{n}\right) \xi_{1}\left(\xi_{1}-x_{1} \xi_{n}\right)$. $P$ is microhyperbolic and triply characteristic over the initial surface $N=\left\{x_{1}=0\right\}$. $\operatorname{Char}(P)$, the (purely imaginary) characteristic variety, is the union of three hypersurfaces $\xi_{1}=0, \pm x_{1} \xi_{n}$, which have a non-involutory intersection $\left\{x_{1}=\xi_{1}=0\right\} \ni p$. Let $b_{j}$ be the bicharacteristic strip of $\left\{\xi_{1}=x_{1} \xi_{n}\right\},\left\{\xi_{1}=0\right\},\left\{\xi_{1}=-x_{1} \xi_{n}\right\}$ for $j=1,2,3$ respectively, issuing from $p$, and $b_{j}^{ \pm}$be its intersection with $\left\{(x ; i \xi d x) ; \pm x_{1}>0\right\}$. We set, as in the second order case,

$$
\operatorname{Sol}(j, \pm)=\left\{u \in\left(\Gamma_{\left\{ \pm x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p} ; u=0 \text { on } b_{k}^{ \pm}(k \neq j)\right\}
$$

An element of $\operatorname{Sol}(j, \pm)$ is called a $j$-pure solution in $\pm x_{1}>0$. First we give the following three theorems, assuming

$$
\text { (\#): } \quad \alpha_{-l}=0 \text { for all } l
$$

Set

$$
\begin{aligned}
Z= & \left\{(a, b) \in \mathbb{C}^{2} ; a=0,-1,-2, \ldots, \quad\right. \text { or } \\
& \quad b=0,-1,-2, \ldots, \quad \text { or } a+b=3 / 2,5 / 2,7 / 2, \ldots\} .
\end{aligned}
$$

Let $p^{\prime}=\rho(p)$, where $\rho$ is the projection $N \times \underset{M}{ } i T^{*} M \rightarrow i T^{*} N, N=\left\{x_{1}=0\right\} \subset M$.
Theorem D. (boundary value problem with purity)
If $(a, b) \notin Z$, the the map

$$
\begin{aligned}
& \operatorname{Sol}(j, \pm) \rightarrow \mathcal{C}_{N, p^{\prime}} \\
& \quad u \mapsto D_{1} u\left(+0, x^{\prime}\right)
\end{aligned}
$$

is an isomorphism.

## Remark.

There is an open dense subset of $\mathbb{C}^{2} \backslash Z$ such that if $(a, b)$ belongs to it, then the mappings

$$
\begin{aligned}
& S o l(j, \pm) \rightarrow \mathcal{C}_{N, p} \\
& u \mapsto u\left(+0, x^{\prime}\right) \\
& u \mapsto D_{1}^{2} u\left(+0, x^{\prime}\right)
\end{aligned}
$$

are isomorphisms.

Theorem E. (characterization of $j$-pure solutions by a relationship between their boundary values)

If $(a, b) \notin Z$, then there exist microdifferential operators $P_{j}^{ \pm(0)}\left(x^{\prime}, D^{\prime}\right)$ and $P_{j}^{ \pm(2)}\left(x^{\prime}, D^{\prime}\right)$ of half integer order that have the following property: an element of $\left(\Gamma_{\left\{ \pm x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p}$ is $j$-pure if and only if

$$
\left\{\begin{array}{l}
u\left( \pm 0, x^{\prime}\right)=P_{j}^{ \pm(0)}\left(\dot{x}^{\prime}, D^{\prime}\right)\left\{D_{1} u\left( \pm 0, x^{\prime}\right)\right\} \\
D_{1}^{2} u\left( \pm 0, x^{\prime}\right)=P_{j}^{ \pm(2)}\left(x^{\prime}, D^{\prime}\right)\left\{D_{1} u\left( \pm 0, x^{\prime}\right)\right\}
\end{array}\right.
$$

Theorem F. (branching of singularities)
(1) There is an open dense subset of $\mathbb{C}^{2} \backslash Z$ such that if $(a, b)$ belongs to it, we have: Let $u(x)$ be an element of $\mathcal{C}_{M, p}^{P}$. If $u$ is pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{1}^{\mp} \cup b_{2}^{\mp} \cup b_{3}^{\mp}$ is contained in suppu.
(2) Assume that $b \in \mathbb{N}=\{1,2,3, \ldots\}$ and $a+b=1 / 2,-1 / 2,-3 / 2,-5 / 2, \ldots$, then we have;
(2-1) If $u$ is 1-pure and $u \neq 0$ in $\pm x_{1}>0$, then $u$ is 1-pure in $\mp x_{1}>0$.
(2-2) If $u$ is 2-pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{1}^{\mp} \cup b_{2}^{\mp} \subset$ suppu and $u=0$ on $b_{3}^{\mp}$.
(2-3) If $u$ is 3-pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{1}^{\mp} \cup b_{3}^{\mp} \subset$ suppu and $u=0$ on $b_{2}^{\mp}$.
(3) Assume that $(a, b) \in \mathbb{N} \times \mathbb{N}$, then we have;
(3-1) If $u$ is 1-pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{1}^{\mp} \cup b_{2}^{\mp} \subset$ supp $u$ and $u=0$ on $b_{3}^{\mp}$.
(3-2) If $u$ is 2-pure and $u \neq 0$ in $\pm x_{1}>0$, then $u$ is 2-pure in $\pm x_{1}>0$.
(3-3) If $u$ is 3-pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{2}^{\mp} \cup b_{3}^{\mp} \subset$ supp $u$ and $u=0$ on $b_{1}^{\mp}$.
(4) Assume that $a \in \mathbb{N}$ and $a+b=1 / 2,-1 / 2,-3 / 2,-5 / 2, \ldots$, then we have;
(4-1) If $u$ is 1-pure and $u \neq 0$ in $\pm x_{1}>.0$, then $b_{1}^{\mp} \cup b_{3}^{\mp} \subset$ supp $u$ and $u=0$ on $b_{2}^{\mp}$.
(4-2) If $u$ is 2-pure and $u \neq 0$ in $\pm x_{1}>0$, then $b_{2}^{\mp} \cup b_{3}^{\mp} \subset$ supp $u$ and $u=0$ on $b_{1}^{\mp}$.
(4-3) If $u$ is 3-pure and $u \neq 0$ in $\pm x_{1}>0$, then $u$ is 3-pure in $\pm x_{1}>0$.
Next, we remove the condition ( $\sharp$ ) and consider the case $\alpha_{-l}$ is not necessarily 0 . We have the following three results. Set

$$
\tilde{Z}=\left\{(a, b) \in \mathbb{C}^{2} ; a \in \mathbb{Z} \quad \text { or } \quad b \in \mathbb{Z} \quad \text { or } \quad a+b+\frac{1}{2} \in \mathbb{Z}\right\} .
$$

## Theorem D'.

There is an open dense subset of $\mathbb{C}^{2} \backslash \tilde{Z}$ such that if $(a, b)$ belongs to it, then, the mappings

$$
\begin{aligned}
\operatorname{Sol}(j, \pm) & \rightarrow \mathcal{C}_{N, p^{\prime}} \\
u & \mapsto u\left(+0, x^{\prime}\right) \\
u & \mapsto D_{1} u\left(+0, x^{\prime}\right) \\
u & \mapsto D_{1}^{2} u\left(+0, x^{\prime}\right)
\end{aligned}
$$

are isomorphisms.

Theorem E'. There is an open dense subset of $\mathbb{C}^{2} \backslash \tilde{Z}$ such that if $(a, b)$ belongs to it, then the same conclusion as Theorem E holds.
Theorem $\mathbf{F}^{\prime}$. There is an open dense subset of $\mathbb{C}^{2} \backslash Z$ such that if $(a, b)$ belongs to it, then the same conclusion as Theorem $F(1)$ holds.
Remark. It is a generic condition that $(a, b)$ belongs to an open dense subset. So in the following proofs, we sometimes say "for a generic $(a, b)$ ", or "generically" instead of mentioning an open dense subset. Those generic conditions will be the avoidance by $(a, b)$ of the zeroes of holomorphic functions $\not \equiv 0$.

Finally we state some results about the case $a$ and $b$ are replaced by 0 -th order microdifferential operators. Let the coordinate of $p^{\prime}$ be ( $\left.\dot{x}_{2}, \ldots, \dot{x}_{n} ; i \dot{\xi}^{\prime} d x^{\prime}\right)$ and $\tilde{a}=$ $\tilde{a}\left(x^{\prime}, D^{\prime}\right), \tilde{b}=\tilde{b}\left(x^{\prime}, D^{\prime}\right)$ be microdifferential operators of order $\leq 0$ defined near $p^{\prime}$ which are commutative: $[\tilde{a}, \tilde{b}]=0$. They have an expansion of the form

$$
\left\{\begin{array}{l}
\tilde{a}\left(x^{\prime}, D^{\prime}\right)=\sum_{j=0}^{\infty} a_{j}\left(x^{\prime \prime}, D^{\prime}\right)\left(x_{n}-\dot{x}_{n}\right)^{j} \\
\tilde{b}\left(x^{\prime}, D^{\prime}\right)=\sum_{j=0}^{\infty} b_{j}\left(x^{\prime \prime}, D^{\prime}\right)\left(x_{n}-\dot{x}_{n}\right)^{j} \\
x^{\prime \prime}=\left(x_{2}, \ldots, x_{n-1}\right) .
\end{array}\right.
$$

Let $\hat{a}=\hat{a}(x, D)$ and $\hat{b}=\hat{b}(x, D)$ be defined by

$$
\begin{aligned}
& \hat{a}(x, D)=\sum_{j=0}^{\infty} a_{j}\left(x^{\prime \prime}, D^{\prime}\right)\left(\frac{1}{2} x_{1} D_{1} D_{n}^{-1}+x_{n}-\dot{x}_{n}\right)^{j} \\
& \hat{b}(x, D)=\sum_{j=0}^{\infty} b_{j}\left(x^{\prime \prime}, D^{\prime}\right)\left(\frac{1}{2} x_{1} D_{1} D_{n}^{-1}+x_{n}-\dot{x}_{n}\right)^{j} .
\end{aligned}
$$

They are operators of order $\leq 0$ defined in a neighborhood of $p$. Set

$$
\begin{aligned}
& a=\sigma_{0}(\tilde{a})\left(p^{\prime}\right)=\sigma_{0}\left(a_{0}\right)\left(p^{\prime}\right)=\sigma_{0}(\hat{a})(p) \\
& b=\sigma_{0}(\tilde{b})\left(p^{\prime}\right)=\sigma_{0}\left(b_{0}\right)\left(p^{\prime}\right)=\sigma_{0}(\hat{b})(p) .
\end{aligned}
$$

Let us consider the operator

$$
\begin{aligned}
P(x, D)= & D_{1}^{3}-x_{1}^{2} D_{n}^{2} D_{1} \\
& +2 D_{n} D_{1}\{\hat{a}(x, D)-\hat{b}(x, D)\} \\
& +x_{1} D_{n}^{2}\{2 \hat{a}(x, D)+2 \hat{b}(x, D)-3\} .
\end{aligned}
$$

Purity and the related mappings are defined in the usual way. In this situation, we have the following theorems $\mathrm{D}^{\prime \prime}, \mathrm{E}^{\prime \prime}$ and F ".
Theorem D". The same statement as Theorem D is true.
Remark. The same statement as the Remark following Theorem D is true. We can take the same open dense subset.

Theorem E". The same statement as Theorem E is true.
Theorem $\mathbf{F}$ ". The same statement as Theorem $F(1)$ is true. (We may take the same open dense subset. )
§2 Jordan-Pochhammer operator and Euler integral representation
Let us consider the following ordinary differential equation of Fuchs type.

$$
\begin{aligned}
J[y]= & \left(x-p_{1}\right)\left(x-p_{2}\right)\left(x-p_{3}\right) y^{\prime \prime \prime} \\
& -\left\{\left(\lambda_{1}-3\right)\left(x-p_{2}\right)\left(x-p_{3}\right)+\left(\lambda_{2}-3\right)\left(x-p_{3}\right)\left(x-p_{1}\right)\right. \\
& \left.\quad+\left(\lambda_{3}-3\right)\left(x-p_{1}\right)\left(x-p_{2}\right)\right\} y^{\prime \prime} \\
& -2\left\{\left(\lambda_{2}+\lambda_{3}-3\right)\left(x-p_{1}\right)+\left(\lambda_{3}+\lambda_{1}-3\right)\left(x-p_{2}\right)\right. \\
& \left.\quad+\left(\lambda_{1}+\lambda_{2}-3\right)\left(x-p_{3}\right)\right\} y^{\prime} \\
& -2\left(\lambda_{1}+\lambda_{2}+\lambda_{3}-3\right) y=0 .
\end{aligned}
$$

We refer to this equation as (JP).
Later, we will set

$$
\begin{gathered}
\lambda_{1}=a, \lambda_{2}=\frac{3}{2}-(a+b), \lambda_{3}=b \\
p_{1}=i, p_{2}=0, p_{3}=-i
\end{gathered}
$$

Lemma 1. If $\lambda_{j} \neq 0,-1,-2,-3, \ldots,($ for $j=1,2,3), \lambda_{1}+\lambda_{2}+\lambda_{3}-4 \neq 1,2,3, \ldots$, then

$$
y_{j}(x)=\int_{p_{j}}^{\infty}\left(u-p_{1}\right)^{\lambda_{1}-1}\left(u-p_{2}\right)^{\lambda_{2}-1}\left(u-p_{3}\right)^{\lambda_{3}-1}(u-x)^{-1} d u
$$

is a solution to (JP). Here the integral is taken in the sense of finite part, if necessary.
PROOF. Although several textbooks (e.g. [Huk], [IKSY]) treat Jordan-Pochhammer equations, ours does not belong to the class solved in them. Therefore, we give an independent proof. See [M] and [I]. Since the finite part is holomorphic in $\lambda_{1}, \lambda_{2}, \lambda_{3}$, we may assume that $\operatorname{Re} \lambda_{j}>0(j=1,2,3)$ and that $\operatorname{Re}\left(\lambda_{1}+\lambda_{2}+\lambda_{3}\right)-4<0$ without loss of generality. We set $y=y_{j}(x)$ in the left hand side of (JP) and write it in terms of powers of $u-x$ by using $x=u-(u-x)$. We have

$$
\begin{aligned}
&\left(x-p_{1}\right)\left(x-p_{2}\right)\left(x-p_{3}\right) \\
&=\left\{\left(u-p_{1}\right)-(u-x)\right\}\left\{\left(u-p_{2}\right)-(u-x)\right\}\left\{\left(u-p_{3}\right)-(u-x)\right\} \\
&=\left(u-p_{1}\right)\left(u-p_{2}\right)\left(u-p_{3}\right) \\
&-\left\{\left(u-p_{1}\right)\left(u-p_{2}\right)+\left(u-p_{2}\right)\left(u-p_{3}\right)+\left(u-p_{3}\right)\left(u-p_{2}\right)\right\}(u-x) \\
&+\left(3 u-p_{1}-p_{2}-p_{3}\right)(u-x)^{2}-(u-x)^{3} \\
&\left(\lambda_{1}-3\right)\left(x-p_{2}\right)\left(x-p_{3}\right)+\ldots \\
&=\left(\lambda_{1}-3\right)\left\{\left(u-p_{2}\right)-(u-x)\right\}\left\{\left(u-p_{3}\right)-(u-x)\right\}+\ldots \\
&=\left(\lambda_{1}-3\right)\left(u-p_{2}\right)\left(u-p_{3}\right)+\left(\lambda_{2}-3\right)\left(u-p_{3}\right)\left(u-p_{1}\right)+\left(\lambda_{3}-3\right)\left(u-p_{1}\right)\left(u-p_{2}\right) \\
&-\left\{\left(\lambda_{1}-3\right)\left(2 u-p_{2}-p_{3}\right)+\left(\lambda_{2}-3\right)\left(2 u-p_{3}-p_{1}\right)+\left(\lambda_{3}-3\right)\left(2 u-p_{1}-p_{2}\right)\right\}(u-x) \\
&+\left(\lambda_{1}+\lambda_{2}+\lambda_{3}-9\right)(u-x)^{2}
\end{aligned}
$$

$$
\begin{aligned}
& \left(\lambda_{2}+\lambda_{3}-3\right)\left(x-p_{1}\right)+\ldots \\
= & \left(\lambda_{2}+\lambda_{3}-3\right)\left\{\left(u-p_{1}\right)-(u-x)\right\}+\ldots \\
= & \left\{\left(\lambda_{2}+\lambda_{3}-3\right)\left(u-p_{1}\right)+\left(\lambda_{3}+\lambda_{1}-3\right)\left(u-p_{2}\right)+\left(\lambda_{1}+\lambda_{2}-3\right)\left(u-p_{3}\right)\right\} \\
& -\left(2 \lambda_{1}+2 \lambda_{2}+2 \lambda_{3}-9\right)(u-x)
\end{aligned}
$$

Moreover

$$
D_{x}^{n} y_{j}(x)=n!\int_{p_{j}}^{\infty}\left(u-p_{1}\right)^{\lambda_{1}-1}\left(u-p_{2}\right)^{\lambda_{2}-1}\left(u-p_{3}\right)^{\lambda_{3}-1}(u-x)^{-1-n} d u
$$

From the equalities above, we have

$$
\begin{aligned}
& J\left[y_{j}\right]=\int_{p_{j}}^{\infty}\left(u-p_{1}\right)^{\lambda_{1}-1}\left(u-p_{2}\right)^{\lambda_{2}-1}\left(u-p_{3}\right)^{\lambda_{3}-1}(u-x)^{-4} \\
& \times\left\{c_{0}+c_{1}(u-x)+c_{2}(u-x)^{2}+c_{3}(u-x)^{3}\right\} d u
\end{aligned}
$$

where

$$
\begin{aligned}
c_{0}= & 6\left(u-p_{1}\right)\left(u-p_{2}\right)\left(u-p_{3}\right) \\
c_{1}= & -6\left\{\left(u-p_{1}\right)\left(u-p_{2}\right)+\ldots\right\}-2\left\{\left(\lambda_{1}-3\right)\left(u-p_{2}\right)\left(u-p_{3}\right)+\ldots\right\} \\
= & -2\left\{\lambda_{1}\left(u-p_{2}\right)\left(u-p_{3}\right)+\lambda_{2}\left(u-p_{3}\right)\left(u-p_{1}\right)+\lambda_{3}\left(u-p_{1}\right)\left(u-p_{2}\right)\right\} \\
c_{2}= & 6\left(3 u-p_{1}-p_{2}-p_{3}\right) \\
& +2\left\{\left(\lambda_{1}-3\right)\left(2 u-p_{2}-p_{3}\right)+\left(\lambda_{2}-3\right)\left(2 u-p_{3}-p_{1}\right)+\left(\lambda_{3}-3\right)\left(2 u-p_{1}-p_{2}\right)\right\} \\
& -2\left\{\left(\lambda_{2}+\lambda_{3}-3\right)\left(u-p_{1}\right)+\left(\lambda_{3}+\lambda_{1}-3\right)\left(u-p_{2}\right)+\left(\lambda_{1}+\lambda_{2}-3\right)\left(u-p_{3}\right)\right\} \\
= & 0 \\
c_{3}= & -6-2\left(\lambda_{1}+\lambda_{2}+\lambda_{3}-9\right) \\
& +2\left(2 \lambda_{1}+2 \lambda_{2}+2 \lambda_{3}-9\right) \\
& -2\left(\lambda_{1}+\lambda_{2}+\lambda_{3}-3\right) \\
= & 0
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
& J\left[y_{j}\right] \\
& =\int_{p_{j}}^{\infty}\left(u-p_{1}\right)^{\lambda_{1}-1}\left(u-p_{2}\right)^{\lambda_{2}-1}\left(u-p_{3}\right)^{\lambda_{3}-1}(u-x)^{-4} \\
& \times\left[6\left(u-p_{1}\right)\left(u-p_{2}\right)\left(u-p_{3}\right)\right. \\
& \left.\quad-2\left\{\lambda_{1}\left(u-p_{2}\right)\left(u-p_{3}\right)+\lambda_{2}\left(u-p_{3}\right)\left(u-p_{1}\right)+\lambda_{3}\left(u-p_{1}\right)\left(u-p_{2}\right)\right\}(u-x)\right] d u \\
& =\int_{p_{j}}^{\infty}\left[-2\left(u-p_{1}\right)^{\lambda_{1}}\left(u-p_{2}\right)^{\lambda_{2}}\left(u-p_{3}\right)^{\lambda_{3}} \frac{\partial}{\partial u}(u-x)^{-3}\right. \\
& \left.\quad-2\left\{\lambda_{1}\left(u-p_{2}\right)\left(u-p_{3}\right)+\ldots\right\}\left(u-p_{1}\right)^{\lambda_{1}-1}\left(u-p_{2}\right)^{\lambda_{2}-1}\left(u-p_{3}\right)^{\lambda_{3}-1}(u-x)^{-3}\right] d u \\
& =\left[-2\left(u-p_{1}\right)^{\lambda_{1}}\left(u-p_{2}\right)^{\lambda_{2}}\left(u-p_{3}\right)^{\lambda_{3}}(u-x)^{-3}\right]_{p_{j}}^{\infty} \\
& +2 \int_{p_{j}}^{\infty} \frac{\partial}{\partial u}\left\{\left(u-p_{1}\right)^{\lambda_{1}}\left(u-p_{2}\right)^{\lambda_{2}}\left(u-p_{3}\right)^{\lambda_{3}}\right\}(u-x)^{-3} d u \\
& \\
& -2 \int_{p_{j}}^{\infty}\left\{\lambda_{1}\left(u-p_{2}\right)\left(u-p_{3}\right)+\ldots\right\}\left(u-p_{1}\right)^{\lambda_{1}-1}\left(u-p_{2}\right)^{\lambda_{2}-1}\left(u-p_{3}\right)^{\lambda_{3}-1}(u-x)^{-3} d u \\
& =
\end{aligned}
$$

Here we have used integration by parts.
If $\lambda_{1}+\lambda_{2}+\lambda_{3}=\frac{3}{2}$, it is easy to see that the Riemann scheme of (JP) is

$$
\left\{\begin{array}{ccccc}
p_{1} & p_{2} & p_{3} & \infty & \\
0 & 0 & 0 & 1 & x \\
1 & 1 & 1 & 3 / 2 & \\
\lambda_{1}-1 & \lambda_{2}-1 & \lambda_{3}-1 & 2 &
\end{array}\right\}
$$

and that $\infty$ is a non-logarithmic singularity. Moreover, if $\lambda_{j} \notin \mathbb{Z}$, then $p_{j}$ is nonlogarithmic.

Lemma 2. An entire solution to (JP) vanishes identically.
PROOF. The characteristic exponents at $\infty$ are larger than 1. Use Liouville's theorem.

Hereafter, we consider the case

$$
\begin{gathered}
\lambda_{1}=a, \lambda_{2}=\frac{3}{2}-(a+b), \lambda_{3}=b \\
p_{1}=i, p_{2}=0, p_{3}=-i
\end{gathered}
$$

$(a, b) \notin Z \underset{\text { def }}{=}\left\{(a, b) \in \mathbb{C}^{2} ; a=0,-1,-2, \ldots\right.$ or $b=0,-1,-2, \ldots$ or $\left.a+b=\frac{3}{2}, \frac{5}{2}, \frac{7}{2}, \ldots\right\}$
Then, our operator is

$$
\begin{aligned}
& Q(x, D) \\
= & \left(x^{3}+x\right) D^{3}+\left\{\frac{15}{2} x^{2}-i(a-b) x+a+b+\frac{3}{2}\right\} D^{2}+\{12 x-2 i(a-b)\} D+3
\end{aligned}
$$

Proposition 3. Take the path of integration from $p_{j}$ to $\infty$ in Reu $\leq 0$. Then $y_{j}$ is holomorphic in Rex $>0$. Moreover, it is holomorphically extended to $x=a_{k}(k \neq j)$, but not to $a_{j}$.
PROOF. $y_{j}$ is obviously holomorphic in the right half plane and at $x=a_{k}$. If it is holomorphic at $a_{j}$, then it is entire. The preceding lemma implies that it vanishes identically. But this is not the case as will be seen when we calculate the expansion coefficients of $y_{j}$ at $x=\infty$.

Let us calculate the expansion coefficients of $y_{j}(j=1,2,3)$ at $x=\infty$, Rex $>0$. We will need the coefficients of $x^{-1}, x^{-\frac{3}{2}}, x^{-2}$ in the next section. For convenience, set $\tau=i / x$. Obviously,

$$
\begin{gathered}
\operatorname{Re} x>0, x=\infty \quad\left(-\frac{\pi}{2}<\arg x<\frac{\pi}{2}\right) \Leftrightarrow \operatorname{Im} \tau>0, \tau=0 \quad(0<\arg \tau<\pi) . \\
\tau^{3 / 2}=e^{3 \pi i / 4} x^{-3 / 2}, \tau^{2}=-1 / x^{2} .
\end{gathered}
$$

We consider the expansion coefficients at

$$
\tau=0, \operatorname{Im} \tau>0
$$

In the following three propositions, we give the coefficients of $\tau, \tau^{\frac{3}{2}}, \tau^{2}$.
Let C be the path in the figure below. Here the $w$-plane has a cut in $\{w ; w \geq 1$ or $w \leq$ $-1\}$.


In the following three propositions, the integrands are continuous on C and

$$
0 \leq \arg w \leq 2 \pi,-\pi \leq \arg (1+w) \leq \pi,-\pi \leq \arg (1-w) \leq \pi .
$$

## Proposition 4.

There exists a nonzero constant $C_{1}$ such that $V_{1}(x) \underset{\text { def }}{=} C_{1} y_{1}(x)$ has expansion coefficients

$$
\begin{aligned}
& p=\frac{1}{\text { def }} \frac{2 \pi i}{2} \int_{C}(1-w)^{a-1} w^{-1 / 2}(1+w)^{b-1} d w \\
& 1 \\
& q=\frac{1}{=} \frac{1}{2 \pi i} \int_{C}(1-w)^{a-1} \cdot w^{-3 / 2}(1+w)^{b-1} d w
\end{aligned}
$$

Remark that $p$ and $q$ are holomorphic in $\{(a, b) ; a \neq 0,-1,-2, \ldots\}$

## Proposition 5.

There exists a nonzero constant $C_{2}$ such that $V_{2}(x) \underset{\text { def }}{=} C_{2} y_{2}(x)$ has expansion coefficients

$$
\begin{aligned}
& r=\frac{1}{\text { def }} \frac{2 \pi}{2 \pi} w^{-\frac{1}{2}}(1-w)^{-(a+b)+\frac{1}{2}}(1-2 w)^{b-1} d w \\
& 1 \\
& s=\frac{-1}{\text { def }} \frac{2 \pi}{2 \pi} w^{-\frac{3}{2}}(1-w)^{-(a+b)+\frac{3}{2}}(1-2 w)^{b-1} d w
\end{aligned}
$$

Remark that $r$ and $s$ are holomorphic in $\{(a, b) ; a+b \neq 3 / 2,5 / 2,7 / 2, \ldots\}$

## Proposition 6.

There exists a nonzero constant $C_{3}$ such that $V_{3}(x) \underset{\text { def }}{=} C_{3} y_{3}(x)$ has expansion coefficients

$$
\begin{aligned}
& t=\frac{1}{\operatorname{def}} \frac{2 \pi}{2 \pi} \int_{C}(1+w)^{a-1} w^{-1 / 2}(1-w)^{b-1} d w \\
& 1 \\
& u \underset{\text { def }}{=}-\frac{1}{2 \pi} \int_{C}(1+w)^{a-1} w^{-3 / 2}(1-w)^{b-1} d w
\end{aligned}
$$

Remark that $t$ and $u$ are holomorphic in $\{(a, b) ; b \neq 0,-1,-2, \ldots\}$
PROOF OF PROPOSITION 4.

$$
y_{1}(x)=\int_{i}^{\infty}(u-i)^{a-1} u^{\frac{1}{2}-(a+b)}(u+i)^{b-1}(u-x)^{-1} d u .
$$

Put $u=i / w, x=i / \tau$. Then

$$
\int_{u=i}^{\infty}=\int_{w=1}^{0}, \quad d u=-\frac{i}{w^{2}} d w .
$$

The path of integration was taken in $\operatorname{Re} u \leq 0$, which corresponds to $\operatorname{Im} w \leq 0$. We have

$$
\begin{aligned}
& y_{1}(x) \\
= & \text { const } \int_{0}^{1}\left(\frac{1}{w}-1\right)^{a-1} w^{a+b-\frac{1}{2}}\left(\frac{1}{w}+1\right)^{b-1}\left(\frac{1}{w}-\frac{1}{\tau}\right)^{-1} \frac{d w}{w^{2}} \\
= & \text { const } \times \tau \int_{0}^{1} w^{\frac{1}{2}}(1-w)^{a-1}(1+w)^{b-1}(\tau-w)^{-1} d w
\end{aligned}
$$

- $\tau(\operatorname{Im} \tau>0)$


$$
=\text { const } \times \tau \int w^{\frac{1}{2}}(1-w)^{a-1}(1+w)^{b-1}(\tau-w)^{-1} d w
$$

In the last expression, the path is the left one in the following figure. Here remark that $\tau$ is outside the path.


Figure 4

The left path is homologous to the right one. The integration around $\tau$ is calculated by means of Cauchy's formula. (Take care of the orientation.)

$$
\begin{aligned}
\frac{y_{1}(x)}{\text { const }}=\tau[ & -2 \pi i \tau^{1 / 2}(1-\tau)^{a-1}(1+\tau)^{b-1} \\
& \left.+\int_{C} w^{1 / 2}(1-w)^{a-1}(1+w)^{b-1}(\tau-w)^{-1} d w\right] .
\end{aligned}
$$

Let us calculate

$$
I=\int_{C} w^{1 / 2}(1-w)^{a-1}(1+w)^{b-1}(\tau-w)^{-1} d w
$$

by deforming C. If $|\tau| \ll 1$, then on the path of integration,

$$
(\tau-w)^{-1}=-w^{-1}\left(1-\frac{\tau}{w}\right)^{-1}=-w^{-1} \sum_{n=0}^{\infty}\left(\frac{\tau}{w}\right)^{n} .
$$

Hence

$$
I=I(\tau, a)=-\int_{C} w^{-1 / 2}(1-w)^{a-1}(1+w)^{b-1} \sum_{n=0}^{\infty}\left(\frac{\tau}{w}\right)^{n} d w .
$$

We can change the order of the integration and the infinite sum at least if $\operatorname{Re} a>0$. The proposition follows in this case.

On the other hand, $I=I(\tau, a)$ is holomorphic in $\tau$ and $a \neq 0,-1,-2, \ldots$. Here we take finite part at $w=1$. Taylor coefficients with respect to $\tau$ is calculated by

$$
\oint \frac{1}{\tau^{m}} I(\tau, a) d \tau
$$

This is holomorphic in $a \neq 0,-1,-2, \ldots$. Therefore the proposition is proved for $a \neq$ $0,-1,-2, \ldots$ by analytic continuation with respect to $a$.

## PROOF OF PROPOSITION 5.

$$
y_{2}(x)=\int_{0}^{\infty}(u-i)^{a-1} u^{\frac{1}{2}-(a+b)}(u+i)^{b-1}(u-x)^{-1} d u .
$$

Set $x=i / \tau$ as before. Moreover, set $u=i-\tilde{w}$ first, and then $\tilde{w}=i / w$. The paths of integration is in $\operatorname{Re} \tilde{w} \geq 0$ and $\operatorname{Im} w \geq 0$ respectively. We have

$$
\begin{aligned}
y_{2}(x) & =\int_{i}^{\infty}(-\tilde{w})^{a-1}(i-\tilde{w})^{-(a+b)+\frac{1}{2}}(2 i-\tilde{w})^{b-1}(i-x-\tilde{w})^{-1}(-d \tilde{w}) \\
& =\text { const } \int_{0}^{1} w^{-a+1}\left(1-\frac{1}{w}\right)^{-(a+b)+\frac{1}{2}}\left(2-\frac{1}{w}\right)^{b-1}\left(\frac{1-\tau}{\tau}+\frac{1}{w}\right)^{-1} \frac{d w}{w^{2}} \\
& =\text { const } \frac{\tau}{1-\tau} \int_{0}^{1} w^{\frac{1}{2}}(1-w)^{-(a+b)+\frac{1}{2}}(1-2 w)^{b-1}\left(w+\frac{\tau}{1-\tau}\right)^{-1} .
\end{aligned}
$$

Set $\theta=\frac{\tau}{1-\tau},|\tau| \ll 1, \operatorname{Im} \tau>0$ and deform the path of integration as below. Here -C is the path obtained by reversing the orientation of C .


Figure 5
We have


$$
\begin{aligned}
\frac{y_{2}(x)}{\text { const }}= & \theta\left[2 \pi i(-\theta)^{\frac{1}{2}}(1+\theta)^{-(a+b)+\frac{1}{2}}(1+2 \theta)^{b-1}\right. \\
& \left.-\int_{C} w^{\frac{1}{2}}(1-w)^{-(a+b)+\frac{1}{2}}(1-2 w)^{b-1}(w+\theta)^{-1} d w\right] \\
& (\pi<\arg (-\theta)<2 \pi) .
\end{aligned}
$$

Here $(-\theta)^{\frac{1}{2}}=i \theta^{\frac{1}{2}}, 0<\arg \theta<\pi$. We expand the right hand side in powers of $\theta$. Then the coefficients of $\theta, \theta^{\frac{3}{2}}, \theta^{2}$ are

$$
\begin{aligned}
& -\int_{C} w^{-\frac{1}{2}}(1-w)^{-(a+b)+\frac{1}{2}}(1-2 w)^{b-1} d w=: I_{1} \\
& \int_{C} w^{-\frac{3}{2}}(1-w)^{-(a+b)+\frac{1}{2}}(1-2 w)^{b-1} d w=: I_{2}
\end{aligned}
$$

Since $\theta=\frac{\tau}{1-\tau}=\tau+\tau^{2}+\tau^{3}+\ldots, \theta^{\frac{1}{2}}=\tau^{\frac{1}{2}}\left(1+\frac{1}{2} \tau+\ldots\right)$, the coefficients of $\tau, \tau^{\frac{3}{2}}$ are $I_{1},-2 \pi$. On the other hand, because

$$
I_{1} \theta+I_{2} \theta^{2}=I_{1}\left(\tau+\tau^{2}+\tau^{3}+\ldots\right)+I_{2}\left(\tau^{2}+2 \tau^{3}+\ldots\right),
$$

the coefficient of $\tau^{2}$ is

$$
\begin{aligned}
I_{1}+I_{2} & =\int_{C}(-w+1) w^{-\frac{3}{2}}(1-w)^{-(a+b)+\frac{1}{2}}(1-2 w)^{b-1} d w \\
& =\int_{C} w^{-\frac{3}{2}}(1-w)^{-(a+b)+\frac{3}{2}}(1-2 w)^{b-1} d w .
\end{aligned}
$$

## PROOF OF PROPOSITION 6.

$$
y_{3}(x)=\int_{-i}^{\infty}(u-i)^{a-1} u^{\frac{1}{2}-(a+b)}(u+i)^{b-1}(u-x)^{-1} d u .
$$

Set $x=i / \tau$ as usual. In addition, we perform a change of variables $u=-i / w$. Then the path of integration would be in $\operatorname{Im} w \geq 0$. We have

$$
\begin{aligned}
y_{3}(x) & =\text { const } \int_{1}^{0}\left(\frac{1}{w}+1\right)^{a-1} w^{(a+b)-\frac{1}{2}}\left(\frac{1}{w}-1\right)^{b-1}\left(\frac{1}{w}+\frac{1}{\tau}\right)^{-1} \frac{d w}{w^{2}} \\
& =\text { const } \times \tau \int_{0}^{1} w^{\frac{1}{2}}(1-w)^{b-1}(1+w)^{a-1}(w+\tau)^{-1} d w .
\end{aligned}
$$

We deform the path in the following way.


Then we have

$$
\begin{aligned}
\frac{y_{3}(x)}{\text { const }}= & \tau\left[2 \pi i(-\tau)^{\frac{1}{2}}(1+\tau)^{b-1}(1-\tau)^{a-1}\right. \\
& \left.-\int_{C} w^{\frac{1}{2}}(1-w)^{b-1}(1+w)^{a-1}(w+\tau)^{-1} d w\right]
\end{aligned}
$$

We expand it into a power series in $\tau$. Since $(-\tau)^{\frac{1}{2}}=i \tau^{\frac{1}{2}}$, the coefficient of $\tau^{\frac{3}{2}}$ is $-2 \pi$. Next, by using

$$
(w+\tau)^{-1}=w^{-1} \sum_{n=0}^{\infty}\left(-\frac{\tau}{w}\right)^{n},
$$

we see that the coefficient of $\tau$ is

$$
-\int_{C} w^{-\frac{1}{2}}(1-w)^{b-1}(1+w)^{a-1} d w
$$

and that the coefficient of $\tau^{2}$ is

$$
\int_{C} w^{-\frac{3}{2}}(1-w)^{b-1}(1+w)^{a-1} d w
$$

In the following section, we use the complex variable $\zeta=x$. $\zeta$ is to be the dual variable of $t$.

## $\S 3$ proof of the unperturbed case

Let us consider

$$
-\frac{1}{8} x_{1}^{3} P(x, D)=-\frac{1}{8} x_{1}^{3}\left[D_{1}^{3}-x_{1}^{2} D_{n}^{2} D_{1}+2(a-b) D_{n} D_{1}+\{2(a+b)-3\} x_{1} D_{n}\right]
$$

Set $t=\frac{1}{2} x_{1}^{2}$ and apply the quantized Legendre transform $\mathcal{L}$. Since

$$
\frac{1}{2} x_{1} D_{1}=t D_{t}, t D_{t} \mapsto-\left(\zeta D_{\zeta}+1\right), t D_{n} \mapsto-i D_{\zeta}
$$

we have

$$
\begin{aligned}
-\frac{1}{8} x_{1}^{3} D_{1}^{3} & =-\frac{1}{8} x_{1} D_{1}\left(x_{1} D_{1}-1\right)\left(x_{1} D_{1}-2\right) \\
& =-\frac{1}{8} \cdot 2 t D_{t}\left(2 t D_{t}-1\right)\left(2 t D_{t}-2\right) \\
& =-t D_{t}\left(t D_{t}-\frac{1}{2}\right)\left(t D_{t}-1\right) \\
\mapsto & \mapsto\left(\zeta D_{\zeta}+1\right)\left(\zeta D_{\zeta}+\frac{3}{2}\right)\left(\zeta D_{\zeta}+2\right) \\
& =\zeta^{3} D_{\zeta}^{3}+\frac{15}{2} \zeta^{2} D_{\zeta}^{2}+12 \zeta D_{\zeta}+3 \\
\frac{1}{8} x_{1}^{5} D_{n}^{2} D_{1}= & \left(\frac{1}{2} x_{1}^{2} D_{n}\right)^{2} \frac{1}{2} x_{1} D_{1}=\left(t D_{n}\right)^{2} t D_{t} \\
\mapsto & -D_{\zeta}^{2}(-1)\left(\zeta D_{\zeta}+1\right)=D_{\zeta}^{2}\left(\zeta D_{\zeta}+1\right) \\
= & \zeta D_{\zeta}^{3}+3 D_{\zeta}^{2} \\
-\frac{1}{4}(a-b) x_{1}^{3} D_{n} D_{1} & =-(a-b) \cdot \frac{1}{2} x_{1}^{2} D_{n} \cdot \frac{1}{2} x_{1} D_{1} \\
& =-(a-b) t D_{n} \cdot t D_{t} \\
& \mapsto-(a-b)\left(-i D_{\zeta}\right)(-1)\left(\zeta D_{\zeta}+1\right) \\
& =-i(a-b) D_{\zeta}\left(\zeta D_{\zeta}+1\right) \\
& =-i(a-b)\left(\zeta D_{\zeta}^{2}+2 D_{\zeta}\right)
\end{aligned}
$$

$$
\begin{aligned}
-\frac{1}{8} x_{1}^{3}\{2(a+b)-3\} x_{1} D_{n}^{2} & =-\frac{1}{2}\{2(a+b)-3\}\left(\frac{1}{2} x_{1}^{2} D_{n}\right)^{2} \\
& =-\frac{1}{2}\{2(a+b)-3\}\left(t D_{n}\right)^{2} \\
& \mapsto-\frac{1}{2}\{2(a+b)-3\}\left(-i D_{\zeta}\right)^{2} \\
& =\frac{1}{2}\{2(a+b)-3\} D_{\zeta}^{2}
\end{aligned}
$$

Summing up, we obtain from $-\frac{1}{8} x_{1}^{3} P$

$$
\begin{gathered}
Q\left(a, b, \zeta, D_{\zeta}\right)=\left(\zeta_{\text {def }}^{=}+\zeta\right) D_{\zeta}^{3}+\left\{\frac{15}{2} \zeta^{2}-i(a-b) \zeta+a+b+\frac{3}{2}\right\} D_{\zeta}^{2} \\
+\{12 \zeta-2 i(a-b)\} D_{\zeta}+3 .
\end{gathered}
$$

We encountered this operator in the previous section. $V_{j}(\zeta)$ is a solution to it.
In the same way as in the second order case, we can construct $E_{j}^{ \pm}$from $V_{j}(\zeta) \cdot L^{ \pm}$ and $B$ are defined accordingly. Let the expansion of $V_{j}(\zeta)(\mathrm{j}=1,2,3)$ at $\zeta=\infty, \operatorname{Re} \zeta>0$ be

$$
V_{j}(\zeta)=\sum_{n=0}^{\infty} V_{j,-1-\frac{n}{2}} \zeta^{-1-\frac{n}{2}}
$$

Then, the matrix $V$, defined by

$$
\left(\begin{array}{ccc}
V_{1,-1} & V_{2,-1} & V_{3,-1} \\
V_{1,-\frac{3}{2}} & V_{2,-\frac{3}{2}} & V_{3,-\frac{3}{2}} \\
V_{1,-2} & V_{2,-2} & V_{3,-2}
\end{array}\right)
$$

is

$$
V=\left(\begin{array}{ccc}
i & 0 & 0 \\
0 & \exp \left(\frac{3}{4} \pi i\right) & 0 \\
0 & 0 & -1
\end{array}\right)\left(\begin{array}{ccc}
p & r & t \\
1 & 1 & 1 \\
q & s & u
\end{array}\right) .
$$

$L^{ \pm}$is expressed by

$$
L^{ \pm}=\left(\begin{array}{ccc}
2 \pi & 0 & 0 \\
0 & \pm 2 \sqrt{2 \pi}\left(\frac{D_{n}}{i}\right)^{\frac{1}{2}} & 0 \\
0 & 0 & 2 \pi \frac{D_{n}}{i}
\end{array}\right) V .
$$

Moreover, we have

$$
\begin{aligned}
& B_{\text {def }}^{=}\left(L^{-}\right)^{-1} L^{+}=\left(\begin{array}{lll}
p & r & t \\
1 & 1 & 1 \\
q & s & u
\end{array}\right)^{-1}\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & -1 & 0 \\
0 & 0 & 1
\end{array}\right)\left(\begin{array}{ccc}
p & r & t \\
1 & 1 & 1 \\
q & s & u
\end{array}\right) \\
&=\frac{-1}{d_{1}+d_{2}+d_{3}}\left(\begin{array}{ccc}
d_{1}-d_{2}-d_{3} & 2 d_{1} & 2 d_{1} \\
2 d_{3} & -d_{1}-d_{2}+d_{3} & 2 d_{3} \\
2 d_{2} & 2 d_{2} & -d_{1}+d_{2}-d_{3}
\end{array}\right)
\end{aligned}
$$

where

$$
d_{1}=\left|\begin{array}{ll}
r & t \\
s & u
\end{array}\right|, d_{2}=\left|\begin{array}{ll}
p & r \\
q & s
\end{array}\right|, d_{3}=\left|\begin{array}{ll}
t & p \\
u & q
\end{array}\right| .
$$

This is checked easily by using the observation that

$$
p d_{1}+r d_{3}+t d_{2}=q d_{1}+s d_{3}+u d_{2}=0 .
$$

## Proposition 7.

$d_{1}, d_{2}$ and $d_{3}$ are holomorphic functions in

$$
\mathbb{C}^{2} \backslash Z=\{a \neq 0,-1,-2, \ldots\} \cap\{b \neq 0,-1,-2, \ldots\} \cap\left\{a+b \neq \frac{3}{2}, \frac{5}{2}, \frac{7}{2}, \ldots\right\} .
$$

Moreover, they don't vanish identically. (Hence generically $d_{1}, d_{2}, d_{3} \neq 0$.)
$P R O O F$. They are obviously holomorphic in $\mathbb{C}^{2} \backslash Z$. The latter part of the proposition follows from Propositions 10, 14 and 15 below.

## Proposition 8.

$\pm d_{1} \pm d_{2} \pm d_{3}$ doesn't vanish identically.
PROOF. This proposition follows immediately from Proposition 10 below.

## Remark 9.

$$
d_{1}+d_{2}+d_{3}=-\left|\begin{array}{lll}
p & r & t \\
1 & 1 & 1 \\
q & s & u
\end{array}\right|
$$

never vanishes, because the components of the matrix are expansion coefficients of three linearly independent solutions.

## Proposition 10.

If $(a, b) \in \mathbb{N} \times \mathbb{N}, \mathbb{N}=\{1,2,3, \ldots\}$, then $r=s=0$. (Hence $d_{1}=d_{2}=0$ and $d_{3}=d_{1}+d_{2}+d_{3} \neq 0$.)
PROOF. Let $h$ be $1 / 2$ or $3 / 2$. Then

$$
-\frac{1}{2} \int_{C} w^{-h}(1-w)^{-(a+b)+h}(1-2 w)^{b-1} d w=\int_{0}^{1} w^{-h}(1-w)^{-(a+b)+h}(1-2 w)^{b-1} d w
$$

Since $(1-2 w)^{b-1}$ is a polynomial of degree $b-1$, it suffices to prove

$$
\int_{0}^{1} w^{c-h}(1-w)^{-(a+b)+h} d w=0 \quad \text { for } \quad c=0,1,2, \ldots, b-1
$$

The left hand side is equal to

$$
B(c-h+1,-a-b+h+1)=\frac{\Gamma(c-h+1) \Gamma(-a-b+h+1)}{\Gamma(-a-b+c+2)}
$$

Here the numerator is finite. The denominator is infinite because $-a-b+c+2$ is a nonpositive integer.

## Proposition 11.

Under the condition of the proposition above, we have

$$
B=\frac{-1}{d_{3}}\left(\begin{array}{ccc}
-d_{3} & 0 & 0 \\
2 d_{3} & d_{3} & 2 d_{3} \\
0 & 0 & -d_{3}
\end{array}\right)=\left(\begin{array}{ccc}
1 & 0 & 0 \\
-2 & -1 & -2 \\
0 & 0 & 1
\end{array}\right) .
$$

## Proposition 12.

If $b \in \mathbb{N}$ and $a+b=\frac{1}{2},-\frac{1}{2},-\frac{3}{2},-\frac{5}{2}, \ldots$, then $p=q=0$.
PROOF. Let $h$ be $1 / 2$ or $3 / 2$. Since $(1+w)^{b-1}$ is a polynomial of degree $b-1$, it suffices to prove

$$
\int_{0}^{1}(1-w)^{a-1} w^{-h} w^{c} d w=0 \text { for } c=0,1,2, \ldots, b-1
$$

The left hand side is equal to

$$
B(c-h+1, a)=\frac{\Gamma(c-h+1) \Gamma(a)}{\Gamma(a+c-h+1)}
$$

Here $a+c-h+1$ is an integer such that

$$
a+c-h+1 \leq\left(\frac{1}{2}-b\right)+(b-1)-h+1=\frac{1}{2}-h \leq 0
$$

## Proposition 13.

If $a \in \mathbb{N}$ and $a+b=\frac{1}{2},-\frac{1}{2},-\frac{3}{2},-\frac{5}{2}, \ldots$, then $t=u=0$.
PROOF. As functions of $(a, b), p, q, t, u$ satisfies

$$
p(b, a)=t(a, b), \quad q(b, a)=-u(a, b) .
$$

## Proposition 14.

Under the condition of Proposition 12, we have $d_{2}=d_{3}=0$. Hence $d_{1}=d_{1}+d_{2}+d_{3} \neq$ 0 and

$$
B=\frac{-1}{d_{1}}\left(\begin{array}{ccc}
d_{1} & 2 d_{1} & 2 d_{1} \\
0 & -d_{1} & 0 \\
0 & 0 & -d_{1}
\end{array}\right)=\left(\begin{array}{ccc}
-1 & -2 & -2 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right) .
$$

## Proposition 15.

Under the condition of Proposition 13, we have $d_{1}=d_{3}=0$. Hence $d_{2}=d_{1}+d_{2}+d_{3} \neq$ 0 and

$$
B=\frac{-1}{d_{2}}\left(\begin{array}{ccc}
-d_{2} & 0 & 0 \\
0 & -d_{2} & 0 \\
2 d_{2} & 2 d_{2} & d_{2}
\end{array}\right)=\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
-2 & -2 & -1
\end{array}\right) .
$$

## PROOF OF THEOREM F.

(1) follows from Proposition 7. (2) (3) and (4) follow from Propositions 14, 11 and 15 respectively.

## PROOF OF THEOREM D.

Recall that

$$
L^{ \pm}=\left(\begin{array}{ccc}
2 \pi & 0 & 0 \\
0 & \pm 2 \sqrt{2 \pi}\left(\frac{D_{n}}{i}\right)^{\frac{1}{2}} & 0 \\
0 & 0 & 2 \pi \frac{D_{n}}{i}
\end{array}\right)\left(\begin{array}{ccc}
i & 0 & 0 \\
0 & \exp \left(\frac{3}{4} \pi i\right) & 0 \\
0 & 0 & -1
\end{array}\right)\left(\begin{array}{ccc}
p & r & t \\
1 & 1 & 1 \\
q & s & u
\end{array}\right) .
$$

We have the following commutative diagram:

$$
\begin{aligned}
\left(\Gamma_{\left\{ \pm x_{1}>0\right\}} \mathcal{C}_{M}^{P}\right)_{p}=\oplus_{j=1}^{3} \operatorname{Sol}(j, \pm) & \stackrel{E^{ \pm} \sim}{\sim}{\stackrel{3}{\oplus} \mathcal{C}_{N, p^{\prime}}}^{\text {b.v.l } \downarrow} \\
& \left\lfloor!L^{ \pm}\right. \\
{ }^{3} \mathcal{C}_{N, p^{\prime}} & =\stackrel{3}{\oplus} \mathcal{C}_{N, p^{\prime}}
\end{aligned}
$$

where $E^{ \pm}=\oplus_{j=1}^{3} E_{j}^{ \pm}$and the first vertical arrow is

$$
u(x) \mapsto\left(u\left( \pm 0, x^{\prime}\right), D_{1} u\left( \pm 0, x^{\prime}\right), D_{1}^{2} u\left( \pm 0, x^{\prime}\right)\right)
$$

By $E^{ \pm}$we identify $\operatorname{Sol}(j, \pm)$ with $\mathcal{C}_{N, p^{\prime}} \oplus 0 \oplus 0,0 \oplus \mathcal{C}_{N, p^{\prime}} \oplus 0,0 \oplus 0 \oplus \mathcal{C}_{N, p^{\prime}}(\mathrm{j}=1,2$, 3). So, in order to prove Theorem D, we have only to prove that the following maps $\in \operatorname{End}\left(\mathcal{C}_{N, p^{\prime}}\right)$ induced by $L^{ \pm}$are automorphisms.

$$
\begin{aligned}
& \underline{j=1} \quad f \mapsto \text { the second component of } L^{ \pm}\left(\begin{array}{l}
f \\
0 \\
0
\end{array}\right) \\
& \underline{j=2} \quad f \mapsto \text { the second component of } L^{ \pm}\left(\begin{array}{l}
0 \\
f \\
0
\end{array}\right) \\
& \underline{j=3} \quad f \mapsto \text { the second component of } L^{ \pm}\left(\begin{array}{l}
0 \\
0 \\
f
\end{array}\right) .
\end{aligned}
$$

For all j , these maps coincide with

$$
f \mapsto \pm 2 \sqrt{2 \pi}\left(\frac{D_{n}}{i}\right)^{\frac{1}{2}} \cdot \exp \left(\frac{3}{4} \pi i\right) \cdot 1 .
$$

They are obviously isomorphisms. The key is that 1 is an elliptic operator.

## PROOF OF THE REMARK AFTER THEOREM D.

We follow nearly the same argument as above. The main difference is that we need the non-vanishing of $p, r$ and $t$ in the Dirichlet case, and that of $q, s$ and $u$ in the case of $D_{1}^{2} u\left( \pm 0, x^{\prime}\right)$ for a generic $(a, b)$. They follow from the lemma below.

## Lemma 16.

$p, q, r, s, t$ and $u$ are holomorphic functions in $(a, b) \in \mathbb{C}^{2} \backslash Z$ which don't vanish identically.

## PROOF .

$p, q, t, u \neq 0$ at $(a, b)=(1,1), r, s \neq 0$ if $b=1$ and $a$ is a half-integer. (Use formulas about the Beta and the Gamma functions ).

## PROOF OF THEOREM E.

Use the same identification as in the proof of Theorem D.
$\S 4$ proof of the perturbed case

## 4-1 the method of the variation of parameters

Let us consider

$$
\begin{gathered}
Q(x, D)=\left(x^{3}+x\right) D^{3}+\left\{\frac{15}{2} x^{2}-i(a-b) x+a+b+\frac{3}{2}\right\} D^{2}+\{12 x-2 i(a-b)\} D+3, \\
a, b \notin \mathbb{Z} \text { and } a+b \notin \frac{1}{2}+\mathbb{Z} .
\end{gathered}
$$

Its Riemann scheme is

$$
\left\{\begin{array}{ccccc}
i & 0 & -i & \infty & \\
0 & 0 & 0 & 1 & x \\
1 & 1 & 1 & \frac{3}{2} & \\
a-1 & \frac{1}{2}-(a+. b) & b-1 & 2 &
\end{array}\right\}
$$

and all the singularities are non-logarithmic. Let $p=i, 0,-i$ and $\varphi_{1}, \varphi_{2}$ and $\varphi_{3}$ be solutions in a neighborhood of $p$. We assume that $\varphi_{1}$ and $\varphi_{2}$ are of exponent 0,1 and that $\varphi_{3}$ is of exponent $a-1, \frac{1}{2}-(a+b), b-1$ if $p=i, 0,-i$ respectively. Set

$$
W(x)=\left|\begin{array}{ccc}
\varphi_{1} & \varphi_{2} & \varphi_{3} \\
\varphi_{1}^{\prime} & \varphi_{2}^{\prime} & \varphi_{3}^{\prime} \\
\varphi_{1}^{\prime \prime} & \varphi_{2}^{\prime \prime} & \varphi_{3}^{\prime \prime}
\end{array}\right|, \quad W_{j k}(x)=\left|\begin{array}{cc}
\varphi_{j} & \varphi_{k} \\
\varphi_{j}^{\prime} & \varphi_{k}^{\prime}
\end{array}\right| .
$$

By using the classical method of the variation of parameters, we see easily that for any $v(x)$, holomorphic near $p$,

$$
\begin{aligned}
I_{p}(v)= & \varphi_{1}(x) \int_{p}^{x} \frac{W_{23}(y)}{\left(y^{3}+y\right) W(y)} v(y) d y \\
& +\varphi_{2}(x) \int_{p}^{x} \frac{W_{31}(y)}{\left(y^{3}+y\right) W(y)} v(y) d y \\
& +\varphi_{3}(x) \int_{p}^{x} \frac{W_{12}(y)}{\left(y^{3}+y\right) W(y)} v(y) d y
\end{aligned}
$$

is a holomorphic function near $p$ such that

$$
Q\left[I_{p}(v)(x)\right]=v(x)
$$

Moreover, we see that

$$
I_{p}(v)(p)=\left\{I_{p}(v)\right\}^{\prime}(p)=0,
$$

that is, $I_{p}(v)$ is of exponent $\in 2+\mathbb{N}_{0}, \mathbb{N}_{0}=\{0,1,2, \ldots\}$.
4-2 a right inverse of $Q$ in a domain containing two regular singular points Let $\Omega$ be a domain $\subset \mathbb{C}_{x}$ as in the figure below.


Figure 7

Obviously $Q$ defines a linear mapping

$$
Q: \mathcal{O}(\Omega) \rightarrow \mathcal{O}(\Omega)
$$

We want to construct a right inverse of this. Let $v(x)$ be an element of $\mathcal{O}(\Omega)$. Although $I_{i} v$ is holomorphic near $i$, there's no guarantee that it should be holomorphic near 0 . We have a similar trouble with $I_{0} v$. In order to overcome this difficulty, we use the following trick.

Let

$$
\Omega_{i} \underset{\operatorname{def}}{=}\left\{x \in \Omega ; \operatorname{Im} x>\frac{1}{3}\right\} \ni i, \quad \Omega_{0} \underset{\text { def }}{=}\left\{x \in \Omega ; \operatorname{Im} x<\frac{2}{3}\right\} \ni 0 .
$$

Obviously, these two domains constitute a-covering of $\Omega$ and

$$
I_{i} v-I_{0} v \in \mathcal{O}^{Q}\left(\Omega_{i} \cap \Omega_{0}\right),
$$

where $\mathcal{O}^{Q}$ is the kernel sheaf of $Q \in \mathcal{E} n d_{\mathrm{C}}(\mathcal{O})$.
If $\left\{F_{1}, F_{2}, F_{3}\right\}$ is a fundamental system of solutions to $Q$ in $\Omega_{i} \cap \Omega_{0}$, then there exists a unique triple of constants $(\alpha, \beta, \gamma) \in \mathbb{C}^{3}$ such that

$$
I_{i} v-I_{0} v=\alpha F_{1}+\beta F_{2}+\gamma F_{3} .
$$

Then obviously

$$
\begin{aligned}
\left(I_{i} v-I_{0} v\right)^{\prime} & =\alpha F_{1}^{\prime}+\beta F_{2}^{\prime}+\gamma F_{3}^{\prime} \\
\left(I_{i} v-I_{0} v\right)^{\prime \prime} & =\alpha F_{1}^{\prime \prime}+\beta F_{2}^{\prime \prime}+\gamma F_{3}^{\prime \prime} .
\end{aligned}
$$

Therefore
(*) $\quad\left(\begin{array}{l}\alpha \\ \beta \\ \gamma\end{array}\right)=\frac{1}{W_{F}}\left(\begin{array}{c}I_{i} v-I_{0} v \\ \left(I_{i} v-I_{0} v\right)^{\prime} \\ \left(I_{i} v-I_{0} v\right)^{\prime \prime}\end{array}\right)$, where $\quad W_{F}=\left|\begin{array}{ccc}F_{1} & F_{2} & F_{3} \\ F_{1}^{\prime} & F_{2}^{\prime} & F_{3}^{\prime \prime} \\ F_{1}^{\prime \prime} & F_{2}^{\prime \prime} & F_{3}^{\prime \prime}\end{array}\right|$.

Since the right hand side of $\left(^{*}\right)$ become constants, it can be evaluated and estimated at any $x$. This observation will be useful later. We define functionals

$$
\begin{aligned}
& \alpha, \beta, \gamma: \mathcal{O}(\Omega) \rightarrow \mathbb{C} \\
& v \mapsto \alpha, \beta, \gamma
\end{aligned}
$$

by means of (*). They mean obstructions. We have to kill them.

## Lemma 17.

For a generic $(a, b)$,

$$
M=\left(\begin{array}{lll}
\alpha(Q[1]) & \alpha(Q[x]) & \alpha\left(Q\left[x^{2}\right]\right) \\
\beta(Q[1]) & \beta(Q[x]) & \beta\left(Q\left[x^{2}\right]\right) \\
\gamma(Q[1]) & \gamma(Q[x]) & \gamma\left(Q\left[x^{2}\right]\right)
\end{array}\right)
$$

is invertible.
PROOF. We prove that $M^{t}(\lambda, \mu, \nu)=\overrightarrow{0}$ implies ${ }^{t}(\lambda, \mu, \nu)=\overrightarrow{0}$. the assumption is that

$$
\alpha\left(Q\left[\lambda+\mu x+\nu x^{2}\right]\right)=\beta\left(Q\left[\lambda+\mu x+\nu x^{2}\right]\right)=\gamma\left(Q\left[\lambda+\mu x+\nu x^{2}\right]\right)=0
$$

This means that $I_{i}\left(Q\left[\lambda+\mu x+\nu x^{2}\right]\right)$ and $I_{0}\left(Q\left[\lambda+\mu x+\nu x^{2}\right]\right)$ are patched together and define $\psi(x) \in \mathcal{O}(\Omega)$. We have $Q \psi=Q\left[\lambda+\mu x+\nu x^{2}\right]$. By the way, we proved before that $\mathcal{O}^{Q}(\Omega)$ is a one-dimensional space generated by $y_{3}(x)$. So there is a constant $c$ such that

$$
\psi-\left(\lambda+\mu x+\nu x^{2}\right)=c y_{3}(x)
$$

Later we'll prove that $c$ is generically 0 . Once we have obtained this, it is clear that

$$
\lambda+\mu x+\nu x^{2} \in \text { Image } I_{i} \cap \text { Image } I_{0}
$$

Therefore it has a zero of order $\geq 2$ at $x=i, 0$. Such a polynomial of degree $\leq 2$ must vanish identically. Hence $\lambda=\mu=\nu=0$.

Now what remains to prove that $c$ is 0 for a generic $(a, b)$. First we prove that if

$$
2 i y_{3}(0)-y_{3}^{\prime}(0)-2 i y_{3}(i)-y_{3}^{\prime}(i) \neq 0
$$

then $c=0$. In fact, if $c \neq 0$,

$$
\frac{\lambda+\mu x+\nu x^{2}}{c}+y_{3}(x)\left(=\frac{1}{c} \psi(x)\right)
$$

has a zero of order $\geq 2$ at $x=i, 0$. Set $\lambda^{\prime}=-\lambda / c, \mu^{\prime}=-\mu / c, \nu^{\prime}=-\nu / c$. Then

$$
\begin{aligned}
& \lambda^{\prime}=y_{3}(0) \\
& \mu^{\prime}=y_{3}^{\prime}(0) \\
& \lambda^{\prime}+i \mu^{\prime}-\nu^{\prime}=y_{3}(i) \\
& \mu^{\prime}+2 i \nu^{\prime}=y_{3}^{\prime}(i)
\end{aligned}
$$

From these we obtain

$$
2 i y_{3}(0)-y_{3}^{\prime}(0)-2 i y_{3}(i)-y_{3}^{\prime}(i)=0,
$$

which constradicts (\#).
Finally, we prove that (\#) holds for a generic ( $a, b$ ). Since the left hand side of (\#) is holomorphic in $(a, b)$, we have only to prove that it is different from 0 for some $(a, b)$.

Recall that

$$
y_{3}(x)=\int_{-i}^{\infty}(u-i)^{a-1} u^{-(a+b)+\frac{1}{2}}(u+i)^{b-1}(u-x)^{-1} d u .
$$

Hence

$$
y_{3}^{\prime}(x)=\int_{-i}^{\infty}(u-i)^{a-1} u^{-(a+b)+\frac{1}{2}}(u+i)^{b-1}(u-x)^{-2} d u .
$$

Let us prove that if $(a, b)=\left(3,-\frac{7}{2}\right)$, we have

$$
y_{3}(0)=y_{3}(i)=y_{3}^{\prime}(i)=0, y_{3}^{\prime}(0) \neq 0 .
$$

Set

$$
B_{3}(p, q)=\int_{-i}^{\infty}(u-i)^{p} u^{q}(u+i)^{b-1} d u .
$$

Then at $(a, b)=\left(3,-\frac{7}{2}\right)$

$$
\begin{aligned}
& y_{3}(0)=B_{3}\left(a-1,-(a+b)-\frac{1}{2}\right)=B_{3}(2,0) \\
& y_{3}^{\prime}(0)=B_{3}\left(a-1,-(a+b)-\frac{3}{2}\right)=B_{3}(2,-1) \\
& y_{3}(i)=B_{3}\left(a-2,-(a+b)+\frac{1}{2}\right)=B_{3}(1,1) \\
& y_{3}^{\prime}(i)=B_{3}\left(a-3,-(a+b)+\frac{1}{2}\right)=B_{3}(0,1)
\end{aligned}
$$

By using a change of variables $u=-i / w$, we have

$$
\begin{aligned}
B_{3}(p, q) & =\text { const } \int_{0}^{1}\left(\frac{w+1}{w}\right)^{p} w^{-q}\left(\frac{1-w}{w}\right)^{b-1} \frac{d w}{w^{2}} \\
& =\text { const } \int_{0}^{1}(1+w)^{p} w^{-p-q-b-1}(1-w)^{b-1} d w
\end{aligned}
$$

Hence at $(a, b)=\left(3,-\frac{7}{2}\right)$,

$$
\begin{aligned}
& y_{3}(0)=\text { const } \int_{0}^{1}(1+w)^{2} w^{\frac{1}{2}}(1-w)^{-\frac{9}{2}} d w=0 . \\
& y_{3}^{\prime}(0)=\text { const } \int_{0}^{1}(1+w)^{2} w^{\frac{3}{2}}(1-w)^{-\frac{2}{2}} d w \neq 0 . \\
& y_{3}(i)=\text { const } \int_{0}^{1}(1+w) w^{\frac{1}{2}}(1-w)^{-\frac{9}{2}} d w=0 . \\
& y_{3}^{\prime}(i)=\text { const } \int_{0}^{1} w^{\frac{3}{2}}(1-w)^{-\frac{9}{2}} d w=0 .
\end{aligned}
$$

This concludes the proof of the lemma.
We can define functionals $\lambda, \mu, \nu: \mathcal{O}(\Omega) \rightarrow \mathbb{C}$ by

$$
\left(\begin{array}{l}
\lambda \\
\mu \\
\nu
\end{array}\right)=M^{-1}\left(\begin{array}{l}
\alpha \\
\beta \\
\gamma
\end{array}\right) .
$$

It is easy to check that

$$
\text { (h) }\left(\begin{array}{l}
\alpha \\
\beta \\
\gamma
\end{array}\right) \circ Q \circ\left(1, x, x^{2}\right)\left(\begin{array}{l}
\lambda \\
\mu \\
\nu
\end{array}\right)=\left(\begin{array}{l}
\alpha \\
\beta \\
\gamma
\end{array}\right): \mathcal{O}(\Omega) \rightarrow \mathbb{C}^{3}
$$

Here the left hand side means

$$
v \mapsto\left(\begin{array}{c}
\alpha\left(Q\left[\lambda(v)+\mu(v) x+\nu(v) x^{2}\right]\right) \\
\beta\left(Q\left[\lambda(v)+\mu(v) x+\nu(v) x^{2}\right]\right) \\
\gamma\left(Q\left[\lambda(v)+\mu(v) x+\nu(v) x^{2}\right]\right)
\end{array}\right) .
$$

Now set

$$
\begin{aligned}
\pi(v) & =v-Q\left[\lambda(v)+\mu(v) x+\nu(v) x^{2}\right] \\
& =v-\lambda(v) Q[1]-\mu(v) Q[x]-\nu(v) Q\left[x^{2}\right] .
\end{aligned}
$$

Then (h) implies that

$$
I_{i} \pi(v)-I_{0} \pi(v)=0
$$

We can define

$$
\tilde{I}_{i, 0}= \begin{cases}I_{i} \pi(v) & \text { on } \Omega_{i} \\ I_{0} \pi(v) & \text { on } \Omega_{0} .\end{cases}
$$

$\tilde{I}_{i, 0}: \mathcal{O}(\Omega) \rightarrow \mathcal{O}(\Omega)$ is a well-defined linear mapping. Next we define $I_{i, 0}: \mathcal{O}(\Omega) \rightarrow \mathcal{O}(\Omega)$ by

$$
I_{i, 0}=\tilde{I}_{i, 0}(v)+\lambda(v)+\mu(v) x+\nu(v) x^{2} .
$$

## Lemma 18.

$Q I_{i, 0}(v)=v$. That is, $I_{i, 0}$ is a right inverse of $Q$.
PROOF. We have

$$
\begin{aligned}
Q I_{i, 0}(v) & =\pi(v)+Q\left[\lambda(v)+\mu(v) x+\nu(v) x^{2}\right] \\
& =v
\end{aligned}
$$

## 4-3 an estimate on the right inverse $I_{i, 0}$

First, let us obtain an estimate on $\alpha, \beta$ and $\gamma$. Fix compact subsets $K_{i} \ni i, K_{0} \ni 0$ of $\Omega$ such that $\operatorname{int}\left(K_{i} \cap K_{0}\right) \neq \phi$. Choose an arbitrary point $\dot{x}$ in int $\left(K_{i} \cap K_{0}\right)$. Then

$$
\left(\begin{array}{l}
\alpha(v) \\
\beta(v) \\
\gamma(v)
\end{array}\right)=\frac{1}{W_{F}(\dot{x})}\left(\begin{array}{c}
\left(I_{i} v-I_{0} v\right)(\dot{x}) \\
\left(I_{i} v-I_{0} v\right)^{\prime}(\dot{x}) \\
\left(I_{i} v-I_{0} v\right)^{\prime \prime}(\dot{x})
\end{array}\right) .
$$

By the way, as in the second order case, we can prove that there exists a constant $C>0$ such that

$$
\sup _{K_{i}}\left|I_{i} v\right| \leq C \sup _{K_{i}}|v|, \sup _{K_{0}}\left|I_{0} v\right| \leq C \sup _{K_{0}}|v| .
$$

Hence, for a larger $C$, we have

$$
|\alpha(v)|,|\beta(v)|,|\gamma(v)| \leq C \sup _{K_{i} \cup K_{0}}|v| .
$$

Therefore, again for a larger $C$, we have

$$
|\lambda(v)|,|\mu(v)|,|\nu(v)| \leq C \sup _{K_{i} \cup K_{0}}|v| .
$$

Set $K_{\delta}=\{x \in \Omega ; \operatorname{dist}(x, \partial \Omega) \geq \delta\}(0<\delta \ll 1)$. Then there exists a constant $C^{\prime}$ independent of $\delta$ such that

$$
\sup _{K_{\delta}^{\prime}}|\pi(v)| \leq C_{K_{\delta}^{\prime}}^{\prime} \sup |v| .
$$

Finally let us obtain an estimate on $I_{i, 0}$. Set

$$
K_{\delta}^{(i)}=K_{\delta} \cap\left\{\operatorname{Im} x \geq \frac{1}{2}\right\}, K_{\delta}^{(0)}=K_{\delta} \cap\left\{\operatorname{Im} x \leq \frac{1}{2}\right\} .
$$

We derive an estimate on $K_{\delta}^{(p)}(p=0, i)$ from the expression

$$
\begin{gathered}
I_{i, 0}=I_{p} \pi(v)+\lambda(v)+\mu(v) x+\nu(v) x^{2} . \\
(p=0, i)
\end{gathered}
$$

There exist constants $\lambda, C_{\Omega}>0$, independent of $\delta$, such that

$$
\sup _{K_{\delta}}\left|I_{i, 0}(v)\right| \leq C_{\Omega} \delta^{-\lambda} \sup _{K_{\delta}}|v| .
$$

( $\lambda$ is determined by the characteristic exponents at $x=-i$, hence by $b$ ). This has the same form as the estimate on $G^{-1}$ in the second order case.

8-4 the end of the proof
Let us construct $E_{1}^{+}$. Let us consider

$$
P(x, D)=D_{1}^{3}-x_{1}^{2} D_{n}^{2} D 1+2(a-b) D_{n} D_{1}+\{2(a+b)-3\} x_{1} D_{n}^{2}+P^{\prime}(x, D)
$$

$$
P^{\prime}(x, D)=\sum_{l=0}^{\text {finite }} \alpha_{-l}\left(x_{1}^{2}, x^{\prime}, D^{\prime}\right) x_{1}^{l+1} D_{1}^{l}
$$

Set $t=\frac{1}{2} x_{1}^{2}$ and apply the quantized Legendre transform $\mathcal{L}$. From $x_{1}^{3} P$, we obtain $Q+Q^{\prime}$, where

$$
\begin{aligned}
& Q=Q\left(\zeta, D_{\zeta}\right) \\
& =\left(\zeta^{3}+\zeta\right) D_{\zeta}^{3}+\left\{\frac{15}{2} \zeta^{2}-i(a-b) \zeta+a+b+\frac{3}{2}\right\} D_{\zeta}^{2}+\{12 \zeta-2 i(a-b)\} D_{\zeta}+3 \\
& \qquad Q^{\prime}=Q^{\prime}\left(\zeta, x^{\prime}, D_{\zeta}, D^{\prime}\right)=\sum_{m=2}^{\text {finite }} \sum_{j=0}^{m-2} \alpha_{m, j}\left(x^{\prime}, D^{\prime}\right) \zeta^{j} D_{\zeta}^{m} \\
& \text { ord } \alpha_{m, j} \leq-m-1
\end{aligned}
$$

Here we have used Part 1 Lemma 4.
With $I_{i, 0}$ instead of $G^{-1}$, we can calculate in the same way as in the second order case. (We don't change $\zeta$ by another complex variable).

The other $E_{j}^{ \pm}$'s are constructed similarly.

## $\S 5$ proof of the case $a$ and $b$ are replaced by operators

In this case, when we perform the process as in the begining of $\S 3$, we obtain the operator

$$
\begin{aligned}
\left(\zeta^{3}+\zeta\right) D_{\zeta}^{3} & +\left\{\frac{15}{2} \zeta^{2}-i \tilde{a}\left(x^{\prime}, D^{\prime}\right) \zeta+i \tilde{b}\left(x^{\prime}, D^{\prime}\right) \zeta+\tilde{a}\left(x^{\prime}, D^{\prime}\right)+\tilde{b}\left(x^{\prime}, D^{\prime}\right)+\frac{3}{2}\right\} D_{\zeta}^{2} \\
& +\left\{12 \zeta-2 i \tilde{a}\left(x^{\prime}, D^{\prime}\right)+2 i \tilde{b}\left(x^{\prime}, D^{\prime}\right)\right\} D_{\zeta}+3
\end{aligned}
$$

which we denote by $Q\left(\tilde{a}\left(x^{\prime}, D^{\prime}\right), \tilde{b}\left(x^{\prime}, D^{\prime}\right), \zeta, D_{\zeta}\right)$. Recall that we have

$$
Q\left(a, b, \zeta, \partial_{\zeta}\right) V_{j}(a, b, \zeta)=0
$$

Here we write $V_{j}(a, b, \zeta)$ instead of $V_{j}(\zeta)$ to specify $a, b$. Remark that $V_{j}$ is holomorphic not only in $\zeta$ but also in $(a, b)$. So we can substitute the commutative pair of operators $\left(\tilde{a}\left(x^{\prime}, D^{\prime}\right), \tilde{b}\left(x^{\prime}, D^{\prime}\right)\right)$ into $(a, b)$ and obtain $V_{j}\left(\tilde{a}\left(x^{\prime}, D^{\prime}\right), \tilde{b}\left(x^{\prime}, D^{\prime}\right), \zeta\right) \in \mathcal{E}(O)$. Obviously, for all $f\left(x^{\prime}\right) \in \mathcal{C}_{N, p^{\prime}}$, we have

$$
Q\left(\tilde{a}, \tilde{b}, \zeta, \partial_{\zeta}\right)\left[V_{j}\left(\tilde{a}\left(x^{\prime}, D^{\prime}\right), \tilde{b}\left(x^{\prime}, D^{\prime}\right), \zeta\right) f\left(x^{\prime}\right)\right]=0
$$

We can easily construct $E_{j}^{ \pm}, L^{ \pm}$and $B$ in this context. For example, we have

$$
L^{ \pm}=\left(\begin{array}{ccc}
2 \pi & 0 & 0 \\
0 & \pm 2 \sqrt{2 \pi}\left(\frac{D_{n}}{i}\right)^{\frac{1}{2}} & 0 \\
0 & 0 & 2 \pi \frac{D_{n}}{i}
\end{array}\right)\left(\begin{array}{ccc}
i & 0 & 0 \\
0 & \exp \left(\frac{3}{4} \pi i\right) & 0 \\
0 & 0 & -1
\end{array}\right)
$$

$$
\times\left(\begin{array}{ccc}
\tilde{p}\left(x^{\prime}, D^{\prime}\right) & \tilde{r}\left(x^{\prime}, D^{\prime}\right) & \tilde{t}\left(x^{\prime}, D^{\prime}\right) \\
1 & 1 & 1 \\
\tilde{q}\left(x^{\prime}, D^{\prime}\right) & \tilde{s}\left(x^{\prime}, D^{\prime}\right) & \tilde{u}\left(x^{\prime}, D^{\prime}\right)
\end{array}\right)
$$

where

$$
\tilde{p}\left(x^{\prime}, D^{\prime}\right)=p\left(\tilde{a}\left(x^{\prime}, D^{\prime}\right), \tilde{b}\left(x^{\prime}, D^{\prime}\right)\right)
$$

etc. It is obvious that

$$
\sigma_{0}\left(\tilde{p}\left(x^{\prime}, D^{\prime}\right)\right)=p\left(\sigma_{0}\left(\tilde{a}\left(x^{\prime}, D^{\prime}\right)\right), \sigma_{0}\left(\tilde{b}\left(x^{\prime}, D^{\prime}\right)\right)\right), \quad \text { etc. }
$$

This observation is used to prove the ellipticity of the components of $L^{ \pm}$and $B$. The remaining part of the proof of Theorems $\mathrm{D}^{\prime \prime}, \mathrm{E}^{\prime \prime}$ and $\mathrm{F}^{\prime \prime}$ is easy.

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7-3-1 Hongo Bunkyo-ku Tokyo, 113, Japan


