## SUBDOMINANT SOLUTIONS OF THE DIFFERENTIAL EQUATION

$$
y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0
$$

BY

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## I. Introduction

1. The solutions $f_{k}$ and $f_{-k}$

In order to state the main results of this paper, we shall start with the differential equation

$$
\begin{equation*}
d^{2} y / d z^{2}-P\left(z, c_{1}, c_{2}, \ldots, c_{m}\right) y=0 \tag{1.1}
\end{equation*}
$$

where $c_{1}, c_{2}, \ldots, c_{m}$ are complex parameters and

$$
\begin{equation*}
P(z, c)=\left(z-c_{1}\right)\left(z-c_{2}\right) \ldots\left(z-c_{m}\right) . \tag{1.2}
\end{equation*}
$$

Let

$$
\begin{equation*}
\left\{\prod_{j=1}^{m}\left(1-c_{j} / z\right)\right\}^{\frac{1}{2}}=1+\sum_{n=1}^{\infty} b_{h}(c) z^{-h} \tag{1.3}
\end{equation*}
$$

where $b_{h}(c)$ are homogeneous polynomials of $c_{1}, \ldots, c_{m}$ of degree $h$ respectively, and let us put

$$
\mathcal{A}_{m}(z, c)=\left\{\begin{array}{l}
z^{\frac{1}{2} m}\left\{1+\sum_{h=1}^{\frac{1}{2}(m+1)} b_{h}(c) z^{-h}\right\} \quad(m=\text { odd })  \tag{1.4}\\
z^{\frac{1}{2} m}\left\{1+\sum_{h=1}^{\frac{1}{2} m} b_{h}(c) z^{-h}\right\}+\frac{b_{1+\frac{1}{2} m}(c)}{1+z} \quad(m=\text { even })
\end{array}\right.
$$

where

$$
\begin{equation*}
z^{r}=\exp \{r[\log |z|+i \arg z]\} \tag{1.5}
\end{equation*}
$$

for any constant r. Previously, P. F. Hsieh and Y. Sibuya [7] constructed a unique solution

$$
\begin{equation*}
y=y_{m}(z, c) \tag{1.6}
\end{equation*}
$$

of the differential equation (1.1) such that
(1) This paper was written with partial support from the National Science Foundation (GP-7041 X).
(i) $\mathcal{Y}_{m}$ is an entire function of $\left(z, c_{1}, \ldots, c_{m}\right)$;
(ii) $Y_{m}$ and $d Y_{m} / d z$ admit respectively the asymptotic representations

$$
\left.\begin{array}{rl}
\mathcal{Y}_{m} & =z^{-\frac{1}{m}}\left\{1+O\left(z^{-\frac{1}{2}}\right)\right\} \exp \left\{-\int_{0}^{z} \mathcal{A}_{m}(t, c) d t\right\}, \\
d Y_{m} / d z & =z^{\frac{1}{2} m}\left\{-1+O\left(z^{-\frac{1}{2}}\right)\right\} \exp \left\{-\int_{0}^{z} \mathcal{A}_{m}(t, c) d t\right\} \tag{1.7}
\end{array}\right\}
$$

uniformly on each compact set in the $\left(c_{1}, \ldots, c_{m}\right)$ space as $z$ tends to infinity in any sector of the form

$$
\begin{equation*}
|\arg z| \leqslant \frac{3 \pi}{(m+2)}-\delta_{0} \tag{1.8}
\end{equation*}
$$

where the path of integration should be taken in the sector (1.8) and $\delta_{0}$ is an arbitrary positive constant.

The solution $\mathcal{Y}_{m}$ tends to zero as $z$ tends to infinity in the sector $|\arg z|<\pi /(m+2)$. We call $\mathcal{Y}_{m}$ a subdominant solution of the differential equation in the sector $|\arg z|<$ $\pi /(m+2)$.

Now let $\lambda$ be a complex parameter, and $a_{1}, \ldots, a_{m}$ fixed real constants such that

$$
\begin{equation*}
-\infty<a_{m}<a_{m-1}<\ldots<a_{2}<a_{1}<+\infty \tag{1.9}
\end{equation*}
$$

Then let us put

$$
\begin{equation*}
\omega=\exp (2 \pi i /(m+2)) \tag{1.10}
\end{equation*}
$$

$$
\begin{gather*}
\varrho=\lambda^{2 /(m+2)}, \quad \arg \varrho=\frac{2}{m+2} \arg \lambda,  \tag{l.11}\\
k= \begin{cases}0,1,2, \ldots,(m+1) / 2 & (m=\mathrm{odd}), \\
0,1,2, \ldots, m / 2,1+m / 2 & (m=\text { even }),\end{cases} \tag{1.12}
\end{gather*}
$$

and

$$
\left.\begin{array}{rl}
f_{k}(x, \lambda) & =y_{m}\left(\varrho \omega^{k} x, \varrho \omega^{k} a_{1}, \varrho \omega^{k} a_{2}, \ldots, \varrho \omega^{k} a_{m}\right)  \tag{1.13}\\
f_{-k}(x, \lambda) & =y_{m}\left(\varrho \omega^{-k} x, \varrho \omega^{-k} a\right) .
\end{array}\right\}
$$

It is easily verified that the functions $y=f_{k}(x, \lambda)$ and $y=f_{-k}(x, \lambda)$ are solutions of the differential equation
(E)

$$
y^{\prime \prime}-\lambda^{2} P(x) y=0
$$

where " $=d^{2} / d x^{2}$ and

$$
\begin{equation*}
P(x)=P(x, a)=\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) \tag{1.14}
\end{equation*}
$$

Furthermore, from the properties (i) and (ii) of the function $y_{m}$ we can derive the following properties of $f_{k}$ and $f_{-k}$ :
(i) The functions $f_{k}$ and $f_{-k}$ are entire functions of ( $x, \varrho$ );
(ii) $f_{k}, f_{k}^{\prime}, f_{-k}$ and $f_{-k}^{\prime}$ admit the asymptotic representations

$$
\begin{align*}
& f_{k}=\left(\varrho \omega^{k}\right)^{-\frac{1}{2} m} x^{-\frac{1}{2} m}\left\{1+O\left(x^{-\frac{1}{2}}\right)\right\} \exp \left\{-\varrho \omega^{k} \int_{0}^{x} \mathcal{A}_{m}\left(\varrho \omega^{k} t, \varrho \omega^{k} a\right) d t\right\}  \tag{1.15-k}\\
& f_{k}^{\prime}=\left(\varrho \omega^{k}\right)^{-\frac{1}{2} m} \lambda x^{\frac{1}{m} m}\left\{(-1)^{k+1}+O\left(x^{-\frac{1}{2}}\right)\right\} \exp \left\{-\varrho \omega^{k} \int_{0}^{x} \mathcal{A}_{m}\left(\varrho \omega^{k} t, \varrho \omega^{k} a\right) d t\right\},
\end{align*}
$$

uniformly on each compact set of the @-plane as $x$ tends to infinity in the sector

$$
\begin{equation*}
\left|\arg \varrho+\arg x+\frac{2 \pi k}{(m+2)}\right| \leqslant \frac{3 \pi}{(m+2)}-\delta_{0}, \tag{1.16-k}
\end{equation*}
$$

and
$f_{-k}=\left(\varrho \omega^{-k}\right)^{-\frac{1}{4} m} x^{-\frac{1}{2} m}\left\{1+O\left(x^{-\frac{1}{2}}\right)\right\} \exp \left\{-\varrho \omega^{-k} \int_{0}^{x} \mathcal{A}_{m}\left(\varrho \omega^{-k} t, \varrho \omega^{-k} a\right) d t\right\}$,
$f_{-k}^{\prime}=\left(\varrho \omega^{-k}\right)^{-\frac{1}{2} m} \lambda x^{\frac{1}{2} m}\left\{(-1)^{k+1}+O\left(x^{-\frac{1}{2}}\right)\right\} \exp \left\{-\varrho \omega^{-k} \int_{0}^{x} \mathcal{A}_{m}\left(\varrho \omega^{-k} t, \varrho \omega^{-k} a\right) d t\right\}$,
uniformly on each compact set of the $\varrho$-plane as $x$ tends to infinity in the sector

$$
\begin{equation*}
\left|\arg \varrho+\arg x-\frac{2 \pi k}{(m+2)}\right| \leqslant \frac{3 \pi}{(m+2)}-\delta_{0} . \tag{-k}
\end{equation*}
$$

The paths of integrations in the right-hand members of (1.15) should be taken in the sectors (1.16) respectively.

Let $S_{k}$ be the sector in the $x$-plane which is defined by

$$
\begin{equation*}
\left|\arg x+\frac{2 \pi k}{(m+2)}\right|<\frac{\pi}{(m+2)} \tag{1.17.k}
\end{equation*}
$$

and let $S_{-k}$ be the sector in the $x$-plane which is defined by

$$
\begin{equation*}
\left|\arg x-\frac{2 \pi k}{(m+2)}\right|<\frac{\pi}{(m+2)} \tag{-k}
\end{equation*}
$$

If $\varrho>0, f_{k}$ tends to zero as $x$ tends to infinity in $S_{k}$, while $f_{-k}$ tends to zero as $x$ tends to infinity in $S_{-k}$. Hence we call $f_{k}$ and $f_{-k}$ subdominant solutions in $S_{k}$ and $S_{-k}$ respectively.

We shall denote by $I_{h}$ the intervals

$$
\begin{equation*}
a_{h+1}<x<a_{n} \quad(h=1,2, \ldots, m-1) \tag{1.18}
\end{equation*}
$$

on the real axis of the $x$-plane respectively. The intervals $-\infty<x<a_{m}$ and $a_{1}<x<+\infty$ will be denoted by $I_{m}$ and $I_{0}$ respectively. Furthermore, $I_{m+1}$ denotes the empty set. Then one of our main concerns in this paper will be to evaluate $f_{k}, f_{k}^{\prime}, f_{-k}$ and $f_{-k}^{\prime}$ when $x \in I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$ as $\lambda$ tends to $\infty$ in a sector $|\arg \lambda| \leqslant \theta_{0}$, where

$$
k= \begin{cases}1,2, \ldots,(m+1) / 2 & (m=\text { odd })  \tag{1.19}\\ 1,2, \ldots, m / 2 & (m=\text { even })\end{cases}
$$

By using these results we shall also compute large positive eigenvalues $\lambda$ of the bound-ary-value problem
(P)

$$
y^{\prime \prime}-\lambda^{2} P(x) y=0, \quad \int_{-\infty}^{+\infty}|y(x)|^{2} d x<+\infty
$$

when $m$ is an even integer.

## 2. The associated Riccati equation

We denote by $\Omega_{0}$ the domain obtained from the complex $x$-plane by deleting the line-segments $a_{2 j}<x<a_{2 j-1}$ on the real axis, where

$$
j= \begin{cases}1,2, \ldots,(m+1) / 2 & (m=\text { odd })  \tag{2.1}\\ 1,2, \ldots, m / 2 & (m=\text { even })\end{cases}
$$

and $a_{m+1}=-\infty$. Let $\mathcal{A}(x)$ be a branch of $P(x)^{\frac{1}{2}}$ such that

$$
\begin{equation*}
\mathcal{A}(x)>0 \tag{2.2}
\end{equation*}
$$

for large positive values of $x$. It is easily verified that $\mathcal{A}(x)$ is single-valued and holomorphic in the interior of $\Omega_{0}$.

By the transformation $w=y^{-1} d y / d x$ we derive the Riccati equation $d w / d x+w^{2}=$ $\lambda^{2} P(x)$ from the equation (E) of Section I. In the domain $\Omega_{0}$ we can write this Riccati equation in the form

$$
\begin{equation*}
w^{\prime}+w^{2}=(\lambda A(x))^{2} \tag{2.3}
\end{equation*}
$$

In order to evaluate the quantities $f_{k}, f_{-k}, f_{k}^{\prime}$ and $f_{-k}^{\prime}$ when $x \in I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$ as $\lambda$ tends to $\infty$, we shall study the Riccati equation (2.3).

First of all, we shall prove the following lemma.
Lemma 1. The Riccati equation (2.3) has two formal solutions

$$
\text { SUBDOMINANT SOLUTIONS OF } y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0
$$

and

$$
\begin{gather*}
w=\lambda A(x)-\frac{1}{4} R(x)+\boldsymbol{p}(x, \lambda)  \tag{2.4}\\
w=-\lambda, A(x)-\frac{1}{4} R(x)+\boldsymbol{p}(x,-\lambda),  \tag{2.5}\\
R(x)=P^{\prime}(x) / P(x),  \tag{2.6}\\
p(x, \lambda)=\sum_{n=1}^{\infty}(\lambda A(x))^{-n} p_{n}(x),  \tag{2.7}\\
\boldsymbol{p}_{n}(x)=Q_{n}(x) P(x)^{-n-1}, \tag{2.8}
\end{gather*}
$$

where
and the quantity $Q_{n}$ is a polynomial in $x$ of degree not greater than $(m-1)(n+1)$.
In fact, if we define $p_{n}$ by

$$
\begin{equation*}
p_{1}(x)=\frac{1}{2}\left\{\frac{1}{4} R^{\prime}(x)-\frac{1}{16} R(x)^{2}\right\} \tag{2.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\boldsymbol{p}_{n}(x)=\frac{1}{2}\left\{\frac{1}{2} n R(x) p_{n-1}(x)-\boldsymbol{p}_{n-1}^{\prime}(x)-\sum_{j+n=n-1} p_{j}(x) \boldsymbol{p}_{h}(x)\right\} \quad(n>1) \tag{2.10}
\end{equation*}
$$

the two formal series (2.4) and (2.5) satisfy the equation (2.3), since $\mathcal{A}^{\prime}(x)=\frac{1}{2} \mathcal{A}(x) R(x)$.
Secondly, we shall prove the following lemma.
Lemma 2. For each fixed point $\xi$ in $I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$, where

$$
k= \begin{cases}1,2, \ldots,(m+1) / 2 & (m=\text { odd })  \tag{2.11}\\ 1,2, \ldots, m / 2 & (m=\text { even })\end{cases}
$$

there is a curve $\quad \mathcal{C}_{k, 1}(\xi): \quad x=z_{k}(s ; \xi) \quad(0 \leqslant s<+\infty)$
and a positive constant $\theta_{0}(\xi)$ such that
(i) $z_{k}(s ; \xi)$ is continuous in $s$ for $0 \leqslant s<+\infty$;
(ii) $\dot{z}_{k}(s ; \xi)$ is bounded and piecewise continuous in $s$ for $0 \leqslant s<+\infty$, where denotes $d / d s$;
(iii) $z_{k}(0 ; \xi)=\xi$;
(iv) $\lim _{s \rightarrow+\infty} z_{k}(s ; \xi)=\infty, \lim _{s \rightarrow+\infty} \arg z_{k}(s ; \xi)=2 \pi k /(m+2)$;
(v) $\operatorname{Im} z_{k}(s ; \xi)>0$ for $0<s<+\infty$;
(vi) the quantity $(-1)^{k} \operatorname{Re}\left[e^{i \theta} \int_{0}^{t} \mathcal{A}\left(z_{k}(s ; \xi)\right) \dot{z}_{k}(s ; \xi) d s\right]$ is nondecreasing for $0 \leqslant t<$ $+\infty i f|\theta| \leqslant \theta_{0}(\xi)$.

This lemma will be proved in Sections 5, 6 of Chapter II. In the proof we shall use the idea that was suggested in the paper of M. A. Evgrafov and M. V. Fedorjuk [3; p. 9]. We shall use some properties of the trajectories of the autonomous system

$$
\begin{equation*}
d x / d t=i \mathcal{A}(\bar{x}) \tag{2.13}
\end{equation*}
$$

where $t$ is a real independent variable and-denotes the complex conjugate. If $x=$ $x(t)$ is a trajectory of (2.13), then we have

$$
\begin{aligned}
\frac{d}{d t} \int_{0}^{t} \mathcal{A}(x(s)) \dot{x}(s) d s & =\mathcal{A}(x(t)) \dot{x}(t)=i \mathcal{A}(x(t)) \mathcal{A}(\overline{x(t)}) \\
& =i \mathcal{A}(x(t)) \overline{\mathcal{A}(x(t)}=i|\mathcal{A}(x(t))|^{2}
\end{aligned}
$$

Hence the real part of the holomorphic function

$$
\int^{x} \mathcal{A}(\tau) d \tau
$$

is constant along the trajectory $x=x(t)$ of (2.13). This is the reason why we shall study the system (2.13).

As a corollary of Lemma 2, we obtain the following lemma.
Lemma 3. Put

$$
\begin{equation*}
\zeta_{k}(s ; \xi)=\overline{z_{k}(s ; \xi)} \quad(0 \leqslant s<+\infty) . \tag{2.14}
\end{equation*}
$$

Then the curve

$$
\mathcal{C}_{k, 2}(\xi): \quad x=\zeta_{k}(s ; \xi)
$$

satisfies the conditions
(i) $\zeta_{k}(s ; \xi)$ is continuous in $s$ for $0 \leqslant s<+\infty$;
(ii) $\dot{\zeta}_{k}(s ; \xi)$ is bounded and piecewise continuous in s for $0 \leqslant s<+\infty$;
(iii) $\zeta_{k}(0 ; \xi)=\xi$;
(iv) $\lim _{s \rightarrow+\infty} \zeta_{k}(s ; \xi)=\infty, \lim _{s \rightarrow+\infty} \arg \zeta_{k}(s ; \xi)=-2 \pi k /(m+2)$;
(v) $\operatorname{Im} \zeta_{k}(s ; \xi)<0$ for $0<s<+\infty$;
(vi) the quantity $(-1)^{k} \operatorname{Re}\left[e^{i \theta} \int_{0}^{t} \mathcal{A}\left(\zeta_{k}(s ; \xi)\right) \xi_{k}(s ; \xi) d s\right]$ is nondecreasing for $0 \leqslant t<$ $+\infty$ if $|\theta| \leqslant \theta_{0}(\xi)$.

Since $\xi$ is real, the conditions (i), (ii), (iii), (iv) and (v) of Lemma 3 can be easily derived from the corresponding conditions of Lemma 2. In order to derive the condition (vi) of Lemma 3 from that of Lemma 2, we use the identity $\mathcal{A}(\bar{x})=\overline{\mathcal{A}(x)}$.

Finally we shall prove the following lemma.

Lemma 4. For each $\xi$ in $I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$, where

$$
k= \begin{cases}1,2, \ldots,(m+1) / 2 & (m=\text { odd })  \tag{2.15}\\ 1,2, \ldots, m / 2 & (m=\text { even })\end{cases}
$$

there exists a function $w_{k, 1}(s, \lambda ; \xi)$ such that
(i) $w_{k, 1}$ is defined and continuous in $(s, \lambda)$ and holomorphic in $\lambda$ for

$$
\begin{equation*}
|\lambda| \geqslant M_{0}(\xi), \quad|\arg \lambda| \leqslant \theta_{0}(\xi), \quad 0 \leqslant s<+\infty, \tag{2.16}
\end{equation*}
$$

where $M_{0}$ is a sufficiently large positive constant depending on $\xi$;
(ii) $\dot{w}_{k, 1}$ is piecewise continuous in $s$ for $0 \leqslant s<+\infty$, where denotes $d / d s$;
(iii) $w_{k, 1}$ satisfies the differential equation

$$
\begin{equation*}
\dot{w}_{k, 1}=\dot{z}_{k}(s ; \xi)\left\{\left[\lambda A\left(z_{k}(s ; \xi)\right)\right]^{2}-w_{k, 1}^{2}\right\} \tag{2.17}
\end{equation*}
$$

for (2.16) except at the points of discontinuity of $\dot{w}_{k, 1}$ in $s$;
(iv) $w_{k, 1}$ satisfies the inequalities

$$
\begin{align*}
\left\lvert\, w_{k, 1}(s, \lambda ; \xi)-(-1)^{k+1} \lambda \mathcal{A}\left(z_{k}\right)+\frac{1}{4}\right. & R\left(z_{k}\right)-\sum_{n=1}^{N}\left[(-1)^{k+1} \lambda \mathcal{A}\left(z_{k}\right)\right]^{-n} p_{n}\left(z_{k}\right) \mid \\
\leqslant & K_{N}(\xi)\left|\lambda \mathcal{A}\left(z_{k}\right)\right|^{-N-1} \quad(N=1,2, \ldots) \tag{2.18}
\end{align*}
$$

for (2.16), where $z_{k}=z_{l c}(\xi ; \xi)$ and $K_{N}$ are positive constants depending on $\xi$.
This lemma will be proved in Section 7 of Chapter II. From Lemma 4 we can derive the following lemma.

Lemma 5. Put

$$
\begin{equation*}
w_{k, 2}(s, \lambda ; \xi)=\overline{w_{k, 1}(s, \lambda ; \xi)} \tag{2.19}
\end{equation*}
$$

Then the function $w_{k, 2}$ satisfies the conditions
(i) $w_{k, 2}$ is defined and continuous in $(s, \lambda)$ and holomorphic in $\lambda$ for (2.16);
(ii) $\dot{w}_{k, 2}$ is piecewise continuous in $s$ for $0 \leqslant s<+\infty$;
(iii) $w_{k, 2}$ satisfies the differential equation

$$
\begin{equation*}
\dot{w}_{k, 2}=\dot{\zeta}_{k}(s ; \xi)\left\{\left[\lambda A\left(\zeta_{k}(s ; \xi)\right)\right]^{2}-w_{k, 2}^{2}\right\} \tag{2.20}
\end{equation*}
$$

for (2.16) except at the points of discontinuity of $\dot{w}_{k, 2}$ in $s$;
(iv) $w_{k, 2}$ satisfies the inequalities

$$
\begin{align*}
\left\lvert\, w_{k, 2}(s, \lambda ; \xi)-(-1)^{k+1} \lambda A\left(\zeta_{k}\right)+\frac{1}{4}\right. & R\left(\zeta_{k}\right)-\sum_{n=1}^{N}\left[(-1)^{k+1} \lambda \mathcal{A}\left(\zeta_{k}\right)\right]^{-n} p_{n}\left(\zeta_{k}\right) \mid \\
\leqslant & K_{N}(\xi)\left|\lambda \mathcal{A}\left(\zeta_{k}\right)\right|^{-N-1} \quad(N=1,2, \ldots) \tag{2.21}
\end{align*}
$$

for (2.16), where $\zeta_{k}=\zeta_{k}(s ; \xi)$.
It is easily seen that the quantities $\dot{w}_{k, 1}$ and $\dot{w}_{k, 2}$ have a discontinuity in $s$ only at the points of discontinuity of $\dot{z}_{k}$ and $\dot{\zeta}_{k}$ respectively. The differential equations (2.17) and (2.20) are derived from the Riccati equation (2.3) by the changes of the variable $x=z_{k}(s ; \xi)$ and $x=\zeta_{k}(s ; \xi)$ respectively. The inequalities (2.18) and (2.21) give asymptotic expansions of $w_{k, 1}$ and $w_{k, 2}$ as $\lambda \mathcal{A}\left(z_{k}\right)$ and $\lambda A\left(\zeta_{k}\right)$ tend to infinity for each fixed $\xi$. We have

$$
\begin{equation*}
w_{k, 1}(s, \lambda ; \xi) \cong(-1)^{k+1} \lambda A\left(z_{k}\right)-\frac{1}{4} R\left(z_{k}\right)+\sum_{n=1}^{\infty}\left[(-1)^{k+1} \lambda A\left(z_{k}\right)\right]^{-n} p_{n}\left(z_{k}\right) \tag{2.22}
\end{equation*}
$$

as $\lambda, A\left(z_{k}\right)$ tends to infinity, and

$$
\begin{equation*}
w_{k, 2}(s, \lambda ; \xi) \cong(-1)^{k+1} \lambda A\left(\zeta_{k}\right)-\frac{1}{4} R\left(\zeta_{k}\right)+\sum_{n=1}^{\infty}\left[(-1)^{k+1} \lambda A\left(\zeta_{k}\right)\right]^{-n} \boldsymbol{p}_{n}\left(\zeta_{k}\right) \tag{2.23}
\end{equation*}
$$

as $\lambda A\left(\zeta_{k}\right)$ tends to infinity. The estimates (2.18) and (2.21) depend on $\xi$ and they are not uniform over any interval which contains the points $a_{2 k-1}$ or $a_{2 k}$. The quantities $\lambda \mathcal{A}\left(z_{k}\right)$ and $\lambda \mathcal{A}\left(\zeta_{k}\right)$ tend to infinity either if $\lambda$ tends to infinity or if $s$ tends to infinity.

## 3. Main results

Let us denote by $\mathcal{B}(x)$ a branch of $P(x)^{\frac{1}{t}}$ such that

$$
\begin{equation*}
\mathcal{B}(x)>0 \tag{3.1}
\end{equation*}
$$

for large positive values of $x$. For real $\tau$, we shall put

$$
\begin{equation*}
\mathcal{A}(\tau-)=\lim _{t \rightarrow 0+} \mathcal{A}(\tau-i t), \quad \mathcal{A}(\tau+)=\lim _{t \rightarrow 0+} \mathcal{A}(\tau+i t) \tag{3.2}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{B}(\tau-)=\lim _{t \rightarrow 0+} \mathcal{B}(\tau-i t), \quad \mathcal{B}(\tau+)=\lim _{t \rightarrow 0+} \mathcal{B}(\tau+i t) \tag{3.3}
\end{equation*}
$$

Let

$$
\begin{equation*}
Q_{k, 1}(s, \lambda ; \xi)=w_{k, 1}(s, \lambda ; \xi)-(-1)^{k+1} \lambda \mathcal{A}\left(z_{k}\right)+\frac{1}{4} R\left(z_{k}\right) \tag{3.4}
\end{equation*}
$$

and $\quad Q_{k, 2}(s, \lambda ; \xi)=w_{k, 2}(s, \lambda ; \xi)-(-1)^{k+1} \lambda A\left(\zeta_{k}\right)+\frac{1}{4} R\left(\zeta_{k}\right)$,
where $z_{k}=z_{k}(s ; \xi)$ and $\zeta_{k}=\zeta_{k}(s ; \xi)$. Finally, put

$$
\begin{equation*}
\alpha_{k}=(-1)^{k} b_{1+m / 2}(a) \tag{3.6}
\end{equation*}
$$

We remark that the function $B(x)$ is single-valued and holomorphic in the $x$-plane cut along the interval $-\infty<x<a_{1}$ on the real axis. In (3.3), we use this holomorphic function $\mathcal{B}(x)$.

Our main result of this paper is the following theorem:
Theorem 1. For each fixed point $\xi$ in $I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$, where

$$
k= \begin{cases}1,2, \ldots,(m+1) / 2 & (m=\text { odd })  \tag{2.11}\\ 1,2, \ldots, m / 2 & (m=\text { even })\end{cases}
$$

we have

$$
\begin{align*}
f_{k}(\xi, \lambda) & =C_{k}(\lambda) \mathcal{B}(\xi-)^{-1} \exp \left[(-1)^{k+1} \lambda \int_{0}^{\xi} \mathcal{A}(\tau-) d \tau-\int_{0}^{+\infty} Q_{k, 2}(s, \lambda ; \xi) \xi_{k}(s ; \xi) d s\right] \\
f_{-k}(\xi, \lambda) & =C_{-k}(\lambda) \mathcal{B}(\xi+)^{-1} \exp \left[(-1)^{k+1} \lambda \int_{0}^{\xi} \mathcal{A}(\tau+) d \tau-\int_{0}^{+\infty} Q_{k, 1}(s, \lambda ; \xi) \dot{z}_{k}(s ; \xi) d s\right]  \tag{3.7}\\
f_{k}^{\prime}(\xi, \lambda) & =\left\{(-1)^{k+1} \lambda A(\xi-)-\frac{1}{4} R(\xi)+Q_{k, 2}(0, \lambda ; \xi)\right\} f_{k}(\xi, \lambda) \\
f_{-k}^{\prime}(\xi, \lambda) & =\left\{(-1)^{k+1} \lambda A(\xi+)-\frac{1}{4} R(\xi)+Q_{k, 1}(0, \lambda ; \xi)\right\} f_{-k}(\xi, \lambda)
\end{align*}
$$

where

$$
\begin{gather*}
C_{k}(\lambda)= \begin{cases}\left(\varrho \omega^{k}\right)^{-\frac{1}{2} m} \exp \left[(-1)^{k+1} \lambda \int_{0}^{+\infty}\left\{\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau-)\right\} d \tau\right] & (m=o d d) \\
\left(\varrho \omega^{k}\right)^{-\frac{1}{2} m-\lambda \alpha_{k}} \exp \left[(-1)^{k+1} \lambda \int_{0}^{+\infty}\left\{\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau-)\right\} d \tau\right] & (m=\text { even }),\end{cases}  \tag{3.8}\\
C_{-k}(\lambda)= \begin{cases}\left(\varrho \omega^{-k}\right)^{-\frac{1}{2} m} \exp \left[(-1)^{k+1} \lambda \int_{0}^{+\infty}\left\{\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau+)\right\} d \tau\right] \quad(m=o d d) \\
\left(\varrho \omega^{-k}\right)^{-\frac{1}{2 m-i \alpha x_{k}} \exp \left[(-1)^{k+1} \lambda \int_{0}^{+\infty}\left\{\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau+)\right\} d \tau\right]} \quad(m=\text { even }),\end{cases} \tag{3.9}
\end{gather*}
$$

and $s$ and $\tau$ are real variables, and

$$
\begin{equation*}
|\lambda| \geqslant M_{0}(\xi), \quad|\arg \lambda| \leqslant \theta_{0}(\xi) . \tag{3.10}
\end{equation*}
$$

Remark. In order to derive asymptotic representations for $f_{k}, f_{-k}, f_{k}^{\prime}$ and $f_{-k}^{\prime}$ as $\lambda$ tends to infinity in the domain (3.10), use the definitions (3.4) and (3.5) for $Q_{k, 1}$ and $Q_{k, 2}$ respectively, and use the asymptotic expansions (2.22) and (2.23) for $w_{k, 1}$ and $w_{k, 2}$ respectively. It was already mentioned that these expansions are not uniform over any interval which contains $a_{2 k-1}$ or $a_{2 k}$. Hence the corresponding asymptotic 16-672909 Acta mathematica 119. Imprimé le 9 février 1968.
representations for $f_{k}, f_{-k}, f_{k}^{\prime}$ and $f_{-k}^{\prime}$ are also not uniform over any interval containing $a_{2 k}$ or $a_{2 k-1}$.

This theorem will be proved in Section 8 of Chapter II.

## 4. Computation of eigenvalues of Problem ( $P$ )

Notice that, since $\omega^{m+2}=1$, the two solutions $f_{1+m / 2}$ and $f_{-1-m / 2}$ are identical. Furthermore, the solution $f_{0}$ is subdominant along the positive real axis, while the solution $f_{1+m / 2}$ is subdominant along the negative real axis. Therefore, eigenvalues of Problem ( $P$ ), when $m$ is even, are zeros of the Wronskian of $f_{0}$ and $f_{1+m / 2}$ :

$$
W(\lambda)=\left|\begin{array}{ll}
f_{0}(x, \lambda) & f_{1+m / 2}(x, \lambda)  \tag{4.1}\\
f_{0}^{\prime}(x, \lambda) & f_{1+m / 2}^{\prime}(x, \lambda)
\end{array}\right|
$$

We may compute large positive zeros of $W(\lambda)$ by using asymptotic evaluations of $f_{0}, f_{0}^{\prime}, f_{1+m / 2}$ and $f_{1+m / 2}^{\prime}$ as $\lambda$ tends to $+\infty$. However, a straightforward application of the theory of asymptotic solutions of ordinary differential equations with respect to parameters would yield asymptotic evaluations of $f_{0}, t_{0}^{\prime}, f_{1+m / 2}$ and $f_{1+m / 2}^{\prime}$ in different $x$-intervals that are disjoint. Hence such results cannot be used for the computation of zeros of $W(\lambda)$. This difficulty is due to the fact that the coefficient $P(x)$ of ( $E$ ) has zeros $a_{1}, \ldots, a_{m}$, These zeros of the coefficient $P(x)$ of $(E)$ are called transition points of ( $E$ ). For a general discussion and the history of the study of the problem of transition points see R. E. Langer [9], J. Heading [6] and W. Wasow [11].

In order to overcome such a difficulty, we need a more elaborate analysis. While in case $m=2$ the problem is reduced to the study of the parabolic cylinder function, there is no such well-known special function that is helpful in case $m>2$. For general discussion of the problem of two transition points see A. Erdélyi, M. Kennedy, and J. L. McGregor [2], R. E. Langer [10], and N. D. Kazarinoff [8]. In this paper, the function $y_{m}$ plays a role similar to that of the parabolic cylinder function in case $m=2$. Actually the functions $f_{k}$ and $f_{-k}$ were derived directly from the function $\mathcal{Y}_{m}$. We shall explain briefly how to evaluate $W(\lambda)$ asymptotically as $\lambda$ tends to $+\infty$. The proof will be given in detail later on in Chapter III. Let us put

$$
\begin{gather*}
\Delta_{k}(\lambda)=\left|\begin{array}{cc}
f_{k} & f_{-k} \\
f_{k}^{\prime} & f_{-k}^{\prime}
\end{array}\right| \quad(k=1,2, \ldots, m / 2)  \tag{4.2}\\
D_{k, 1}(\lambda)=\left|\begin{array}{cc}
f_{k} & f_{k+1} \\
f_{k}^{\prime} & f_{k+1}^{\prime}
\end{array}\right| \quad(k=0,1, \ldots, m / 2) \tag{4.3}
\end{gather*}
$$

and

$$
\begin{align*}
& D_{k, 2}(\lambda)=\left|\begin{array}{ll}
f_{-k} & f_{-k-1} \\
f_{-k}^{\prime} & f_{-k-1}^{\prime}
\end{array}\right| \quad(k=0,1, \ldots, m / 2)  \tag{4.4}\\
& \Delta_{k, 1}(\lambda)=\left|\begin{array}{ll}
f_{-k} & f_{k+1} \\
f_{-k}^{\prime} & f_{k+1}^{\prime}
\end{array}\right| \quad(k=1,2, \ldots,-1+m / 2)  \tag{4.5}\\
& \Delta_{k, 2}(\lambda)=\left|\begin{array}{ll}
f_{k} & f_{-k-1} \\
f_{k}^{\prime} & f_{-k-1}^{\prime}
\end{array}\right| \quad(k=1,2, \ldots,-1+m / 2) \tag{4.6}
\end{align*}
$$

First of all by using Theorem 1, we shall evaluate $\Delta_{k}(\lambda)$ at a fixed point of the interval $I_{2 k-1}$ for large positive values of $\lambda$. This will show that $f_{k}$ and $f_{-k}$ are linearly independent if $\lambda$ is positive and large. Let $\Phi_{k}(x, \lambda)$ be the matrix of these two independent solutions, i.e.

Then

$$
\begin{gather*}
\Phi_{k}(x, \lambda)=\left(\begin{array}{ll}
f_{k}(x, \lambda) & f_{-k}(x, \lambda) \\
f_{k}^{\prime}(x, \lambda) & f_{-k}^{\prime}(x, \lambda)
\end{array}\right) .  \tag{4.7}\\
\Phi_{k}(x, \lambda)^{-1}=\Delta_{k}(\lambda)^{-1}\left(\begin{array}{cc}
f_{-k}^{\prime} & -f_{-k} \\
-f_{k}^{\prime} & f_{k}
\end{array}\right) . \tag{4.8}
\end{gather*}
$$

Hence if we put

$$
H_{k c}(\lambda)=\left(\begin{array}{rr}
\Delta_{k, 2}(\lambda) & D_{k, 2}(\lambda)  \tag{4.9}\\
-D_{k, 1}(\lambda) & -\Delta_{k, 1}(\lambda)
\end{array}\right) \quad(k=1, \ldots,-1+m / 2),
$$

we get

$$
\begin{equation*}
\Phi_{k}(x, \lambda)=\Delta_{k+1}(\lambda)^{-1} \Phi_{k+1}(x, \lambda) H_{k}(\lambda) \tag{4.10}
\end{equation*}
$$

and

$$
\begin{equation*}
\Phi_{1}(x, \lambda)=\left\{\Delta_{2}(\lambda) \Delta_{3}(\lambda) \ldots \Delta_{m / 2}(\lambda)\right\}^{-1} \Phi_{m / 2}(x, \lambda) H(\lambda) \tag{4.11}
\end{equation*}
$$

where

$$
\begin{equation*}
H(\lambda)=H_{-1+m / 2}(\lambda) \ldots H_{2}(\lambda) H_{1}(\lambda) \tag{4.12}
\end{equation*}
$$

On the other hand, we have

$$
\begin{align*}
\binom{f_{0}(x, \lambda)}{f_{0}^{\prime}(x, \lambda)} & =\Delta_{1}(\lambda)^{-1} \Phi_{1}(x, \lambda)\binom{D_{0,2}(\lambda)}{-D_{0,1}(\lambda)}  \tag{4.13}\\
\text { and } \quad\binom{f_{1+m / 2}(x, \lambda)}{f_{1+m / 2}^{\prime}(x, \lambda)} & =\Delta_{m / 2}(\lambda)^{-1} \Phi_{m / 2}(x, \lambda)\binom{-D_{\frac{1}{2} m, 2}(\lambda)}{D_{\frac{1}{2} m, 1}(\lambda)} . \tag{4.14}
\end{align*}
$$

Hence we get

$$
\begin{equation*}
W(\lambda)=\left\{\Delta_{1}(\lambda) \Delta_{2}(\lambda) \ldots \Delta_{\frac{1}{2} m}(\lambda)\right\}^{-1}\left[D_{\frac{1}{2} m, 1}(\lambda), D_{\frac{1}{2} m, 2}(\lambda)\right] H(\lambda)\binom{D_{0,2}(\lambda)}{-D_{0.1}(\lambda)} \tag{4.15}
\end{equation*}
$$

Therefore, large positive eigenvalues of Problem ( $P$ ) are zeros of

$$
\begin{equation*}
\left[D_{\frac{1}{2} m, 1}(\lambda), D_{\frac{1}{2} m, 2}(\lambda)\right] H(\lambda)\binom{D_{0.2}(\lambda)}{-D_{0.1}(\lambda)} \tag{4.16}
\end{equation*}
$$

We shall evaluate the quantities $D_{k, 1}$ and $D_{k, 2}$ at $x=\infty$ by using asymptotic representations (1.15) of $f_{k}, f_{-k}$ and their derivatives. On the other hand, by using Theorem 1, we shall evaluate the quantities $\Delta_{k, 1}$ and $\Delta_{k, 2}$ at a fixed point of $I_{2 k}$. In this manner, we shall prove the following theorem.

## Theorem 2. Let

$$
\lambda_{k}(\eta)=\frac{\left(\eta+\frac{1}{2}\right) \pi}{\int_{a_{2 k}}^{a_{2 k-1}} \sqrt{|P(\tau)|} d \tau}+\sum_{n=1}^{\infty} \mathcal{E}_{k, n} \eta^{-n} \quad(k=1,2, \ldots, m / 2)
$$

be formal power series in $\eta^{-1}$ with constant coefficients $\mathcal{E}_{k, n}$ which satisfy the formal equations

$$
\begin{gather*}
\lambda_{k} \int_{a_{2 k}}^{a_{2 k-1}} \sqrt{|P(\tau)|} d \tau=\left(\eta+\frac{1}{2}\right) \pi-(i / 2) \sum_{n=1}^{\infty} \lambda_{k}^{-2 n+1} \int_{v_{k}}\left[(-1)^{k+1} \mathcal{A}(x)\right]^{-2 n+1} p_{2 n-1}(x) d x  \tag{4.17}\\
(k=1,2, \ldots, m / 2)
\end{gather*}
$$

where $\gamma_{k}$ is a circle which encircles only $a_{2 k}$ and $a_{2 k-1}$, and the integration must be taken in the counterclockwise sense. Then there is a positive integer $\nu_{0}$ such that we can denote almost all positive eigenvalues of the boundary-value problem $(P)$ by $\lambda_{k, v}(k=1,2, \ldots, m / 2$; $\left.v=1+\nu_{0}, 2+\nu_{0}, \ldots\right)$ so that we have

$$
\lambda_{k, v}=\frac{\left(v+\frac{1}{2}\right) \pi}{\int_{a_{2 k}}^{a_{2 k-1}} \sqrt{|P(\tau)|} d \tau}+\sum_{n=1}^{N} \mathcal{E}_{k, n} v^{-n}+O\left(v^{-N-1}\right)
$$

as $\nu$ tends to $+\infty$, where $N$ may be any positive integer.

Remark. Previously M. V. Fedorjuk [4] claimed that the eigenvalues $\lambda_{k \nu}$ are determined asymptotically as $v$ tends to $+\infty$ by the formal equations

$$
\begin{equation*}
\lambda_{k, \nu} \int_{a_{2 k}}^{a_{2 k-1}} \sqrt{|P(\tau)|} d \tau=\left(v+\frac{1}{2}\right) \pi-\frac{i}{2} \sum_{n=1}^{\infty} \lambda_{k, \nu}^{n} \int_{\gamma_{k}}\left[(-1)^{k+1} \mathcal{A}(x)\right]^{-n} p_{n}(x) d s \tag{4.18}
\end{equation*}
$$

However, we felt that he gave only a sketch of the proof. Furthermore, since then, M. A. Evgrafov and M. V. Fedorjuk [3] explained again exactly the same method, and they gave only the first approximations of $\lambda_{k, v}$, i.e.

$$
\begin{align*}
& \text { SUBDOMINANT SOLUTIONS OF } y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0 \\
& \qquad \lambda_{k, \nu}=\frac{\left(v+\frac{1}{2}\right) \pi}{\int_{a_{2 k}}^{a_{2 k-1}} \sqrt{|P(\tau)|} d \tau}+o\left(\frac{1}{v}\right) . \tag{4.19}
\end{align*}
$$

In addition, they remarked that we need a more elaborate analysis to obtain higher order approximations. Therefore, we felt that it might be worth-while to present another proof for the results of M. V. Fedorjuk. By using an idea due to N. Fröman [5] we can prove that

$$
\begin{equation*}
\int_{\gamma_{k}} \mathcal{A}(x)^{-2 n} p_{2 n}(x) d x=0 \quad(n=1,2, \ldots) \tag{4.20}
\end{equation*}
$$

Hence Theorem 2 gives us the same results as those of M. V. Fedorjuk. We shall give a proof of (4.20) in Section 10 of Chapter III. The analysis of M. V. Fedorjuk is based on the computation of the so-called Stokes multipliers around a transition point after another. Our analysis is based on the computation of the Stokes multipliers around a pair of transition points $a_{2 k-1}$ and $a_{2 k}$ after another pair.

The author of the present paper wishes to express his thanks to Professor Masahiro Iwano for valuable advice.

## II. Proof of Theorem 1

## 5. Proof of Lemma 2 (Part I)

We shall prove Lemma 2 of Section 2 in Sections 5 and 6. First of all, we shall construct some particular trajectories of the autonomous system (2.13)

$$
d x / d t=i \mathcal{A}(\bar{x})
$$

where $t$ is a real independent variable and - denotes the complex conjugate, and prove the following lemma:

Lemma 6. For each $h=1,2, \ldots, m$ the system (2.13) has a unique trajectory $\mathcal{J}_{h}$

$$
\begin{equation*}
x=x_{h}(t) \tag{5.1}
\end{equation*}
$$

such that $x_{h}(t)$ is defined for $t_{h}<t<0$ if $h=2$ or $3 \bmod 4$ and $x_{h}(t)$ is defined for $0<t<t_{h}$ if $h=0$ or $1 \bmod 4$, where $\left|t_{h}\right|=+\infty$ if $m=1,2$, and $\left|t_{h}\right|<+\infty$ otherwise, and such that
(i) $\lim _{t \rightarrow 0} x_{h}(t)=a_{h}$;
(ii) $\operatorname{Im} x_{h}(t)>0$ for $t_{h}<t<0$ or $0<t<t_{h}$;
(iii) $\lim _{t \rightarrow t_{h}}\left|x_{h}(t)\right|=+\infty$;
(iv) there is an integer $k$ such that

$$
\lim _{t \rightarrow t_{h}} \arg x_{h}(t)=\frac{(2 k-1) \pi}{(m+2)}
$$

and that
(a)

$$
1 \leqslant k \leqslant(m+3) / 2 \quad(m=o d d), \quad 1 \leqslant l \leqslant 1+m / 2 \quad(m=\text { even }) ;
$$

(b)

$$
\begin{aligned}
& k \text { is even if } h=2 \text { or } 3 \bmod 4 ; \\
& k \text { is odd if } h=0 \text { or } 1 \bmod 4 .
\end{aligned}
$$

In Lemma 6, the integers $k$ are not exactly specified. If this lemma is proved, we can determine the exact values of the integers $k$ by using the fact that the trajectories $\mathcal{J}_{\boldsymbol{h}}$ cannot intersect each other. We can prove the following lemma.

Lemma 7. The trajectories $\mathfrak{J}_{h}$ satisfy the conditions

$$
\begin{gather*}
\lim _{t \rightarrow t_{h}} \arg x_{h}(t)=h \pi /(m+2) \text { if } h \text { is odd, }  \tag{5.2}\\
\lim _{t \rightarrow t_{h}} \arg x_{h}(t)=(h+1) \pi /(m+2) \text { if } h \text { is even. } \tag{5.3}
\end{gather*}
$$

Now we shall prove Lemma 6. Let

$$
\begin{aligned}
\mathcal{D}_{2 j-1} & =\left\{x ;\left|x-a_{2 j-1}\right|<\delta, \quad\left|\arg \left(x-a_{2 j-1}\right)\right|<\pi\right\}, \\
\mathcal{D}_{2 j} & =\left\{x ;\left|x-a_{2 j}\right|<\delta, \quad\left|\arg \left(x-a_{2 j}\right)-\pi\right|<\pi\right\},
\end{aligned}
$$

where $\delta$ is a small positive number. Then $\mathcal{D}_{h}$ are contained in $\Omega_{0}$ if $\delta$ is sufficiently small. Hence the functions

$$
\begin{equation*}
F_{h}(x)=\int_{a_{h}}^{x} \mathcal{A}(\tau) d \tau \quad(h=1,2, \ldots, m) \tag{5.4}
\end{equation*}
$$

are holomorphic in $\mathcal{D}_{h}$ respectively. Let $x=x(t)\left(t_{0} \leqslant t \leqslant t_{1}\right)$ be a trajectory of (2.13) which is contained in $D_{h}$ for some $h$. Then

$$
d F_{h}(x(t)) / d t=\mathcal{A}(x(t)) \quad d x(t) / d t=i|\mathcal{A}(x(t))|^{2} .
$$

Hence $\operatorname{Re} F_{h}(x(t))$ is constant. On the other hand, if a curve $x=x(\tau)\left(\tau_{0} \leqslant \tau \leqslant \tau_{1}\right)$ satisfies the conditions:
(i) $x(\tau) \in D_{h}$ for some $h$;
(ii) $d x(\tau) / d \tau \neq 0$ for $\tau_{0} \leqslant \tau \leqslant \tau_{1}$;
(iii) $\operatorname{Re} \boldsymbol{F}_{h}(x(\tau)) \equiv$ constant for $\tau_{0} \leqslant \tau \leqslant \tau_{1}$;
then this curve is a trajectory of (2.13). In fact

$$
\alpha(\tau)=-i d F_{h}(x(\tau)) / d \tau
$$

is real and different from zero for $\tau_{0} \leqslant \tau \leqslant \tau_{1}$. Let

$$
\begin{equation*}
t=\int_{\tau_{0}}^{\tau} \alpha(s)|\mathcal{A}(x(s))|^{-2} d s \quad\left(\tau_{0} \leqslant \tau \leqslant \tau_{1}\right) . \tag{5.5}
\end{equation*}
$$

Then the function $x=x(\tau(t))$ satisfies the differential equation (2.13), where $\tau(t)$ is the inverse function of (5.5).

Now we shall study the function $\boldsymbol{F}_{h}$ in the neighborhood of $a_{h}$. It is easily seen that in the interval $a_{2 j-1}<x<a_{2(j-1)}$ we have

$$
\mathcal{A}(x) \begin{cases}>0 & \text { if } j \text { is odd, } \\ <0 \text { if } j \text { is even }\end{cases}
$$

where $a_{0}=+\infty$. Therefore

$$
\mathcal{A}(x)=\left\{\begin{array}{l}
(-1)^{j+1}\left(x-a_{2 j-1}\right)^{\frac{1}{2}}\left\{\phi_{2 j-1}+O\left(x-a_{2 j-1}\right)\right\} \quad \text { in } D_{2 j-1},  \tag{5.6}\\
(-1)^{j+1}\left(x-a_{2 j}\right)^{\frac{1}{2}}\left\{i \phi_{2 j}+O\left(x-a_{2 j}\right)\right\} \quad \text { in } D_{2 j}
\end{array}\right.
$$

where

$$
\arg \left[\left(x-a_{h}\right)^{\frac{1}{2}}\right]=\frac{1}{2} \arg \left(x-a_{h}\right)
$$

and

$$
\phi_{h}=\sqrt{\prod_{n \neq h}\left|a_{h}-a_{n}\right|}>0
$$

$\left.\begin{array}{ll}\text { Hence } & F_{2 j-1}=(-1)^{j+1}\left(x-a_{2 j-1}\right)^{\frac{3}{2}}\left\{\frac{2}{3} \phi_{2 j-1}+O\left(x-a_{2 j-1}\right)\right\} \text { in } \mathcal{D}_{2 j-1}, \\ & F_{2 j}=(-1)^{j+1}\left(x-a_{2 j}\right)^{\frac{1}{2}}\left\{\frac{2}{3} i \phi_{2 j}+O\left(x-a_{2 j}\right)\right\} \text { in } D_{2 j},\end{array}\right\}$
where $\arg \left[\left(x-a_{h}\right)^{\frac{1}{2}}\right]=\frac{3}{2} \arg \left(x-a_{h}\right)$. From (5.7) we can conclude that there are two trajectories in $\mathcal{D}_{h}$ which approach the point $a_{h}$. Two such trajectories can be given by
and

$$
\begin{array}{ll}
\theta=\frac{\pi}{3}+O(r) \quad \text { and } \quad \theta=-\frac{\pi}{3}+O(r) & \text { in } \mathcal{D}_{2 j-1} \\
\theta=\frac{2 \pi}{3}+O(r) \quad \text { and } \quad \theta=\frac{4 \pi}{3}+O(r) & \text { in } \mathcal{D}_{2 j}
\end{array}
$$

if we put $x-a_{h}=r e^{i \theta}$. Let us denote by $\mathcal{J}_{h}$ the one of these two trajectories in $\mathcal{D}_{h}$
which lies in the upper half-plane in the neighborhood of $a_{h}$. If $b_{h}$ is a point on the trajectory $\mathscr{J}_{h}$, this trajectory is determined by the solution of (2.13) which satisfies the initial condition $x(0)=b_{h}$. Let $x=y_{h}(t)$ be this unique solution. If we put

$$
y_{h}(t)-a_{h}=r(t) e^{i \theta(t)}
$$

then we get

$$
\begin{equation*}
\theta(t)=\frac{\pi}{3}+O(r(t)) \quad \text { if } h \text { is odd } \tag{5.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\theta(t)=\frac{2 \pi}{3}+O(r(t)) \quad \text { if } h \text { is even. } \tag{5.9}
\end{equation*}
$$

From (2.13), (5.6), (5.8) and (5.9), we derive

$$
\frac{d r(t)}{d t}= \begin{cases}(-1)^{j+1} r(t)^{\frac{1}{2}} \phi_{2 j-1}\{1+O(r(t))\} & \text { in } D_{2 j-1}, \\ (-1)^{j+1} r(t)^{\frac{1}{2}} \phi_{2 j}\{-1+O(r(t))\} & \text { in } D_{2 j} .\end{cases}
$$

Hence in $\mathcal{D}_{h}$ we have

$$
\frac{d r}{d t}\left\{\begin{array}{lll}
>0 & \text { if } & h=0 \text { or } 1 \bmod 4 \\
<0 & \text { if } & h=2 \text { or } 3 \bmod 4
\end{array}\right.
$$

On the other hand, since $d t / d r=O\left(r^{-\frac{1}{2}}\right)$, the trajectory $x=y_{h}(t)$ arrives at $a_{n}$ from $b_{h}$ within a finite time-interval. Thus we can construct the unique trajectories $\mathcal{T}_{h}$ which satisfy the conditions (i) and (ii) of Lemma 6 in the neighborhood of the point $a_{h}$.

Now we shall complete the proof of Lemma 6 for the case $h=3 \bmod 4$. Other cases can be treated in the same manner. Assume that the trajectory

$$
\begin{equation*}
x=x_{h}(t) \quad t_{h}<t<0 \tag{5.10}
\end{equation*}
$$

satisfies the conditions (i) and (ii) of Lemma 6. Furthermore, assume that this trajectory cannot be extended without violating the condition (ii). We want to prove that the trajectory (5.10) satisfies the condition (iii) of Lemma 6.

To do this, first of all, assume that there is a decreasing sequence of values $\hat{t}_{n}$ of $t$ such that $\lim _{n \rightarrow \infty} \hat{t}_{n}=t_{h}$, and that $\lim _{n \rightarrow \infty} x_{h}\left(\hat{t}_{n}\right)=\xi$ exists and $\operatorname{Im} \xi>0$. If, in this case, $\left|t_{h}\right|<+\infty$, then we have $\lim _{t \rightarrow t_{h}} x_{h}(t)=\xi$. This contradicts the assumption that trajectory (5.10) cannot be extended without violating the condition (ii) of Lemma 6. Hence we must have $t_{h}=-\infty$. Now assume that $t_{h}=-\infty$. Then since $\operatorname{Im} x_{h}(t)>0$ for $0>t>-\infty$ and $\operatorname{Im} x_{h}(0)=\operatorname{Im} a_{h}=0$, this trajectory must approach a limit cycle as $t$ tends to $-\infty$. Otherwise, we must have $\lim _{n \rightarrow \infty} x_{h}\left(t_{n}\right)=\xi=\infty$. This limit cycle must be in the upper half-plane. However, this is impossible, since there is no zero of $\mathcal{A}(\bar{\xi})$ in the upper half-plane. (Poincaré-Bendixson Theorem [1; pp. 389-403].)

$$
\text { SUBDOMINANT SOLUTIONS OF } y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0
$$

Secondly, assume that there is a decreasing sequence of values $\hat{t}_{n}$ of $t$ such that $\lim _{n \rightarrow+\infty} \hat{t}_{n}=t_{h}$, and that $\lim _{n \rightarrow+\infty} x_{h}\left(\hat{t}_{n}\right)=\xi$ exists and $\operatorname{Im} \xi=0$. Then we can prove that

$$
\begin{equation*}
\lim _{t \rightarrow t_{h}} x_{h}(t)=\xi \quad \text { and } \quad\left|t_{h}\right|<+\infty \tag{5.11}
\end{equation*}
$$

In fact, if $a_{2 p+1}<\xi<a_{2 p}$, then $\mathcal{A}(\xi)$ is real and different from zero. Hence the assertion (5.11) can be easily verified. In case when $\xi=a_{p}$, let us consider the function

$$
F_{p}(x)=\int_{a_{p}}^{x} A(\tau) d \tau
$$

in the upper half-plane $\Omega$. The function $F_{p}$ is holomorphic in $\Omega$ and $\operatorname{Re} F_{p}\left(x_{h}(t)\right)$ is constant for $\boldsymbol{t}_{h}<\boldsymbol{t}<\mathbf{0}$. On the other hand, $\boldsymbol{F}_{p}\left(x_{h}\left(\hat{t}_{n}\right)\right)$ tends to zero as $n$ tends to infinity. Hence

$$
\operatorname{Re} F_{p}\left(x_{h}(t)\right) \equiv 0 \quad \text { for } \quad t_{h}<t<0
$$

Therefore, in the neighborhood of $a_{p}$, the trajectory (5.10) must coincide with the trajectory $\mathcal{J}_{\mathfrak{p}}$. This implies (5.11). We shall now prove that $\xi$ is not in the interval $a_{2 p}<x<a_{2 p-1}$. To do this, assume that $a_{2 p}<\xi<a_{2_{p-1}}$. Then consider the function

$$
F_{2 p}(x)=\int_{a_{2 p}}^{x} \mathcal{A}(\tau) d \tau
$$

in $\Omega$. The function $F_{2 p}$ is holomorphic in $\Omega$ and $\operatorname{Re} F_{2 p}\left(x_{h}(t)\right)$ is constant for $t_{\boldsymbol{h}}<t<0$. It is easily seen that $\mathcal{A}(x)$ is purely imaginary for $a_{2 p}<x<a_{2 p-1}$. Hence $\operatorname{Re} F_{2 p}\left(x_{h}\left(\hat{l}_{n}\right)\right)$ tends to zero as $n$ tends to infinity, and we have

$$
\operatorname{Re} F_{2 p}\left(x_{h}(t)\right) \equiv 0 \quad \text { for } \quad t_{h}<t<0 .
$$

However, for sufficiently large $n$, it is easily verified that

$$
\operatorname{Re} F_{2 p}\left(x_{n}\left(\hat{t}_{n}\right)\right) \neq 0
$$

This is a contradiction. Thus we have proved (5.11). Moreover, we also have shown that

$$
\begin{equation*}
a_{2 p+1} \leqslant \xi \leqslant a_{2 p} \quad \text { for some } p . \tag{5.12}
\end{equation*}
$$

Now we shall derive a contradiction from the conditions (i), (ii), (5.11) and (5.12). In fact, since (5.10) is a trajectory of (2.13), then $x=\overline{x_{h}(-t)}\left(0<t<t_{h}\right)$ is also a trajectory of (2.13). Hence if the conditions (i), (ii), (5.11), and (5.12) are satisfied, we have a closed trajectory of (2.13). By modifying this closed trajectory at $a_{h}$ and $\xi$, if necessary, we can construct a simple closed curve in $\Omega_{0}$ so that $\arg \mathcal{A}(\bar{x})$ increases by
$2 \pi$ if $x$ goes along this closed curve counterclockwise. This means that the index of this closed curve with respect to the vector field $i \mathcal{A}(\bar{x})$ is one. (See, for example, [1; pp. 398-400].) This implies that

$$
\int d \log \mathcal{A}(x)=-2 \pi i
$$

if the integration is taken along this closed curve counterclockwise. However, every residue of $\mathcal{A}^{\prime}(x) / \mathcal{A}(x)$ is $\frac{1}{2}$. Hence $\int d \log \mathcal{A}(x)>0$. Thus we arrive at a contradiction. Consequently, we proved that

$$
\begin{equation*}
\lim _{t \rightarrow t_{h}}\left|x_{h}(t)\right|=+\infty \tag{5.13}
\end{equation*}
$$

In order to prove the condition (iv) of Lemma 6, we shall consider the equation (2.13) at $x=\infty$. We know that $\mathcal{A}(x)$ has the form

$$
\mathcal{A}(x)=x^{\frac{2}{2} m}\left\{1+\sum_{n=1}^{\infty} b_{h}(a) x^{-n}\right\}
$$

in the neighborhood of $x=\infty$. Let us define $F_{\infty}(x)$ by

$$
F_{\infty}(x)=\left\{\begin{array}{l}
x^{\frac{1}{2} m+1}\left\{\frac{2}{m+2}+\sum_{l=1}^{\infty} \frac{2}{m+2-2 l} b_{l}(a) x^{-l}\right\} \quad \text { if } m \text { is odd, } \\
x^{\frac{1}{2} m+1}\left\{\frac{2}{m+2}+\sum_{2 l \neq m+2} \frac{2}{m+2-2 l} b_{l}(a) x^{-l}\right\}+b_{1+m / 2}(a) \log x \quad \text { if } m \text { is even, }
\end{array}\right.
$$

where $\arg F_{\infty}(x)=0$ for large positive $x$. Then the function $F_{\infty}$ is holomorphic in $\Omega$. Furthermore, $\operatorname{Re} F_{\infty}\left(x_{h}(t)\right)$ is constant for $t_{h}<t<0$. In the neighborhood of $x=\infty$, any curve defined by $\operatorname{Re} F_{\infty}(x)=$ constant has one of the following forms

$$
\theta=(2 k-1) \pi /(m+2)+O\left(r^{-1}\right)
$$

if $x=r e^{i \theta}$ and $0<\theta<\pi$, where

$$
\mathbf{1} \leqslant k \leqslant(m+3) / 2 \quad(m=\text { odd }), \quad 1 \leqslant k \leqslant 1+m / 2 \quad(m=\text { even }) .
$$

On the other hand, along such a curve, we have

Hence

$$
\frac{d r}{d t}=r^{\frac{1}{2} m}\left\{\sin \left(\frac{m+2}{2} \theta\right)+O\left(r^{-1}\right)\right\}=r^{\frac{1}{2}}\left\{(-1)^{k+1}+O\left(r^{-1}\right)\right\}
$$

Moreover, $d t / d r=O\left(r^{-\frac{1}{2} m}\right)$. This completes the proof of Lemma 6 for the case $h=3 \bmod 4$.

## 6. Proof of Lemma 2 (Part II)

Now we shall construct the curve $\mathcal{C}_{k, 1}(\xi)$ of Lemma 2 for $\xi \in I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$ when $k$ is an even integer. The case when $k$ is odd can be treated in the same manner. Let $\Omega_{k}$ be the simply connected subdomain of the upper half-plane $\Omega$ which is bounded by $\mathcal{J}_{2 k}, I_{2_{k-1}}$, and $\mathcal{J}_{2 k-1}$. On the other hand, for each point $x_{0}$ in the interior of $\Omega$, we denote by $x=\omega\left(t, x_{0}\right)$ the trajectory of the system (2.13) such that $x_{0}=\omega\left(0, x_{0}\right)$.

Let $x_{0} \in \Omega_{k}$, and assume that the trajectory $\omega\left(t, x_{0}\right)$ is defined for $t_{0}^{\prime}<t<t_{0}$, where $-t_{0}^{\prime}$ and $t_{0}$ may be $+\infty$. Then $\omega\left(t, x_{0}\right)$ has the following properties:

$$
\begin{equation*}
\lim _{t \rightarrow t_{0}} \omega\left(t, x_{0}\right)=\infty, \quad \lim _{t \rightarrow t_{0}} \arg \omega\left(t, x_{0}\right)=(2 k+1) \pi /(m+2) \tag{6.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{t \rightarrow t_{0}} \omega\left(t, x_{0}\right)=\infty, \quad \lim _{t \rightarrow t_{0}} \arg \omega\left(t, x_{0}\right)=(2 k-1) \pi /(m+2) \tag{6.2}
\end{equation*}
$$

In fact, the trajectory $\omega\left(t, x_{0}\right)$ should be in the interior of $\Omega_{k}$ for $t_{0}^{\prime}<t<t_{0}$. Then by using the Poincaré-Bendixson Theorem [1; pp. 389-403] and Lemma 7, we can prove (6.1) and (6.2).

Now for $\xi \in I_{2 k-1}$ we define $\tilde{\mathcal{C}}_{k, 1}(\xi)$ by

$$
x=\tilde{z}_{k}(s ; \xi)=\left\{\begin{array}{l}
\xi+i s \quad\left(0 \leqslant s \leqslant s_{1}\right),  \tag{6.3}\\
\omega\left(s-s_{1}, \tilde{z}_{k}\left(s_{1} ; \xi\right)\right) \quad\left(s_{1} \leqslant s \leqslant s_{2}\right), \\
\tilde{z}_{k}\left(s_{2} ; \xi\right)+\left(s-s_{2}\right) \exp \{2 k \pi i /(m+2)\} \quad\left(s_{2} \leqslant s<+\infty\right) .
\end{array}\right.
$$

(See Fig. 1.)
If $s_{1}$ and $1 / s_{2}$ are sufficiently small, the function

$$
F(t)=\operatorname{Re} \int_{0}^{t} \mathcal{A}\left(\tilde{z}_{k}(s ; \xi)\right) \dot{\tilde{z}}_{k}(s ; \xi) d s
$$

is nondecreasing. In fact,

$$
F(t)= \begin{cases}-\operatorname{Im}[\mathcal{A}(\xi+i t)], & \left(0 \leqslant t \leqslant s_{1}\right), \\ 0 & \left(s_{1} \leqslant t \leqslant s_{2}\right), \\ \operatorname{Re}\left[\mathcal{A}\left(\tilde{z}_{k}(t ; \xi)\right) \exp \{2 k \pi i /(m+2)\}\right] \quad\left(s_{2} \leqslant t<+\infty\right)\end{cases}
$$

Notice that, if $k$ is even, we have

$$
\mathcal{A}(x)=\left\{\begin{array}{rll}
\sqrt{|P(x)|} & \text { if } & x \in I_{2 k} \\
-i \sqrt{|\bar{P}(x)|} & \text { if } & x \in I_{2 k-1} \\
-\sqrt{|P(x)|} & \text { if } & x \in I_{2(k-1)}
\end{array}\right.
$$



Fig. 1.
where $\sqrt{|P(x)|}>0$. Hence, if $s_{1}$ is sufficiently small, we have

$$
-\operatorname{Im}[\mathcal{A}(\xi+i t)]>\frac{1}{2} \sqrt{|P(x)|}>0 \quad \text { for } \quad 0 \leqslant t \leqslant s_{1}
$$

On the other hand, if $s_{2}$ is sufficiently large, we have

$$
k \pi<\arg \left[\mathcal{A}\left(\tilde{z}_{k}(t ; \xi)\right) \exp \{2 k \pi i /(m+2)\}\right]<\left(k+\frac{1}{2}\right) \pi
$$

for $s_{2} \leqslant t<+\infty$. Hence

$$
\operatorname{Re}\left[\mathcal{A}\left(\tilde{z}_{k}(t ; \xi)\right) \exp \{2 k \pi i /(m+2)\}\right]>0
$$

Thus we get $\dot{F}(t) \geqslant 0$ for all $t$.
For $\xi \in I_{2 k}$, let $\tilde{\mathcal{C}}_{k, 1}(\xi)$ be defined hy

$$
x=\tilde{z}_{k}(s ; \xi)= \begin{cases}\omega(s, \xi) & \left(0 \leqslant s \leqslant s_{1}\right)  \tag{6.4}\\ \tilde{z}_{k}\left(s_{1} ; \xi\right)+\left(s-s_{1}\right) & \left(s_{1} \leqslant s \leqslant s_{2}\right) \\ \omega\left(s-s_{2}, \tilde{z}_{k}\left(s_{2} ; \xi\right)\right) & \left(s_{2} \leqslant s \leqslant s_{3}\right) \\ \tilde{z}_{c}\left(s_{3} ; \xi\right)+\left(s-s_{3}\right) \overline{A\left(\tilde{z}_{k}\left(s_{3} ; \xi\right)\right)} & \left(s_{3} \leqslant s \leqslant s_{4}\right) \\ \omega\left(s-s_{4}, \tilde{z}_{k}\left(s_{4} ; \xi\right)\right) & \left(s_{4} \leqslant s \leqslant s_{5}\right) \\ \tilde{z}_{k}\left(s_{5} ; \xi\right)+\left(s-s_{5}\right) \exp \{2 k \pi i /(m+2)\} & \left(s_{5} \leqslant s<+\infty\right)\end{cases}
$$



Fig. 2.
(See Fig. 2.) We can determine $s_{1}, s_{2}, s_{3}, s_{4}$, and $s_{5}$ so that the function

$$
\operatorname{Re} \int_{0}^{t} \mathcal{A}\left(\tilde{z}_{k}(s ; \xi)\right) \dot{z}(s ; \xi) d s
$$

is nondecreasing.
Finally, for $\xi \in I_{2(k-1)}$, let $\tilde{\mathcal{C}}_{k, 1}(\xi)$ be defined by

$$
x=\tilde{z}_{k}(s ; \xi)= \begin{cases}\omega(-s, \xi) & \left(0 \leqslant s \leqslant s_{1}\right)  \tag{6.5}\\ \tilde{z}_{k}\left(s_{1} ; \xi\right)-\left(s-s_{1}\right) & \left(s_{1} \leqslant s \leqslant s_{2}\right) \\ \omega\left(s_{2}-s, \tilde{z}_{k}\left(s_{2} ; \xi\right)\right) & \left(s_{2} \leqslant s \leqslant s_{3}\right) \\ \tilde{z}_{k}\left(s_{3} ; \xi\right)+\left(s-s_{3}\right) \bar{A}\left(\overline{\tilde{z}_{k}}\left(s_{3} ; \xi\right)\right) & \left(s_{3} \leqslant s \leqslant s_{4}\right) \\ \omega\left(s-s_{4}, \tilde{z}_{k}\left(s_{4} ; \xi\right)\right) & \left(s_{4} \leqslant s \leqslant s_{5}\right) \\ \tilde{z}_{k}\left(s_{5} ; \xi\right)+\left(s-s_{5}\right) \exp \{2 k \pi i /(m+2)\} & \left(s_{5} \leqslant s<+\infty\right)\end{cases}
$$

(See Fig. 3.)
Since the mapping $z=\int_{\xi}^{x} \mathcal{A}(t) d t$ is conformal, we can find $\mathcal{C}_{k, 1}(\xi)$ and $\theta_{0}(\xi)$ by modifying $\tilde{\mathcal{C}}_{k, 1}(\xi)$ in a suitable manner.


Fig. 3.

## 7. Proof of Lemma 4

Let $\quad \Omega(\delta)=\left\{x ; x \in \Omega,\left|x-a_{h}\right| \geqslant \delta>0,(h=1,2, \ldots, m)\right\}$,
where $\Omega$ is the upper half-plane. Then

$$
\begin{equation*}
|\mathcal{A}(x)| \geqslant \delta^{\frac{1}{2} m}, \quad|\arg \mathcal{A}(x)| \leqslant \gamma \quad \text { for } \quad x \in \Omega(\delta) \tag{7.2}
\end{equation*}
$$

where $\gamma$ is a positive constant. By using the Borel-Ritt Theorem [11; p. 47] we can construct a function $g(u, x)$ of two complex variables $(u, x)$ so that
(i) $\mathfrak{g}$ is holomorphic for

$$
\begin{equation*}
x \in \Omega(\delta), \quad|u| \geqslant u_{0}>0, \quad|\arg u| \leqslant \gamma+2 \pi \tag{7.3}
\end{equation*}
$$

(ii) $g$ satisfies the inequalities

$$
\begin{align*}
& \left|\boldsymbol{g}(u, x)-\sum_{n=1}^{N} u^{-n} \boldsymbol{p}_{n}(x)\right| \leqslant E_{N}(\delta)|u|^{-N-1} \\
& \left|\partial \mathfrak{g}(u, x) / \partial u-\sum_{n=1}^{N}(-n) u^{-n-1} \boldsymbol{p}_{n}(x)\right| \leqslant E_{N}(\delta)|u|^{-N-2}  \tag{7.4}\\
& \left|\partial \boldsymbol{g}(u, x) / \partial x-\sum_{n=1}^{N} u^{-n} \boldsymbol{p}_{n}^{\prime}(x)\right| \leqslant E_{N}(\delta)|u|^{-N-1}
\end{align*}
$$

for (7.3) and $N=1,2, \ldots$, where $E_{N}$ are positive constants depending of $\delta$.

$$
\text { SUBDOMINANT SOLUTIONS OF } y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0
$$

Put

$$
\begin{equation*}
q(x, \lambda)=g\left((-1)^{k+1} \lambda A(x), x\right) . \tag{7.5}
\end{equation*}
$$

Then we can prove the following lemma.
Lemma 8. The function $\mathfrak{q}(x, \lambda)$ is holomorphic in $(x, \lambda)$ for

$$
\begin{equation*}
x \in \Omega(\delta), \quad|\lambda| \geqslant u_{0} \delta^{-\frac{1}{2} m}, \quad|\arg \lambda| \leqslant \theta_{0}(\xi) . \tag{7.6}
\end{equation*}
$$

Moreover, $q$ satisfies the inequalities

$$
\left\{\begin{array}{l}
\left|\boldsymbol{q}(x, \lambda)-\sum_{n=1}^{N}\left[(-1)^{k+1} \lambda \mathcal{A}(x)\right]^{-n} \boldsymbol{p}_{n}(x)\right| \leqslant B_{N}(\delta)|\lambda A(x)|^{-N-1} \\
\left|\boldsymbol{q}^{\prime}(x, \lambda)-\sum_{n=1}^{N}\left[(-1)^{k+1} \lambda \mathcal{A}(x)\right]^{-n}\left(-\frac{1}{2} n R(x) \boldsymbol{p}_{n}(x)+\boldsymbol{p}_{n}^{\prime}(x)\right)\right| \leqslant B_{N}(\delta)|\lambda \mathcal{A}(x)|^{-N-1} \tag{7.7}
\end{array}\right\}
$$

for (7.6) and $N=1,2, \ldots$, where $B_{N}$ are positive constants depending on $\delta$.
In fact, if $(x, \lambda)$ is in the region (7.6), then the inequalities (7.2) imply that

$$
|\lambda \mathcal{A}(x)| \geqslant u_{0} \quad \text { and }\left|\arg \left[(-1)^{k+1} \lambda \mathcal{A}(x)\right]\right| \leqslant \gamma+2 \pi .
$$

Hence $q$ is holomorphic for (7.6). The first inequality of (7.7) is then derived from the first inequality of (7.4). On the other hand, the formula

$$
\mathfrak{q}^{\prime}(x, \lambda)=(-1)^{k+1} \lambda \mathcal{A}^{\prime}(x) \partial \mathfrak{g}\left((-1)^{k+1} \lambda \mathcal{A}(x), x\right) / \partial u+\partial \boldsymbol{g}\left((-1)^{k+1} \lambda A(x), x\right) / \partial x
$$

and the second and the third inequalities of (7.4) imply the second inequality of (7.7). Here we used the identity

$$
\mathcal{A}^{\prime}(x)=\frac{1}{2} \mathcal{A}(x) R(x)
$$

and the fact that $|R(x)|$ is bounded in $\Omega(\delta)$.
Let us transform the Riccati equation (2.3) by

$$
\begin{equation*}
w=(-1)^{k+1} \lambda A(x)+\frac{1}{4} R(x)+q(x, \lambda)+\lambda A(x) v . \tag{7.8}
\end{equation*}
$$

Then we have

$$
\begin{equation*}
v^{\prime}=L(x, \lambda)+\left((-1)^{k} 2 \lambda \mathcal{A}(x)-2 \boldsymbol{q}(x, \lambda)\right) v-\lambda \mathcal{A}(x) v^{2} \tag{7.9}
\end{equation*}
$$

where $\quad L(x, \lambda)=-(\lambda, A(x))^{-1}\left[q^{\prime}(x, \lambda)+\frac{1}{16} R(x)^{2}-\frac{1}{4} R^{\prime}(x)\right.$

$$
\begin{equation*}
\left.+2\left((-1)^{k+1} \lambda A(x)+\frac{1}{4} R(x)\right) q(x, \lambda)+q(x, \lambda)^{2}\right] . \tag{7.10}
\end{equation*}
$$

By using Lemma 8 and the definitions (2.9) and (2.10) of the functions $\boldsymbol{p}_{n}(x)$, we can prove that

$$
\begin{equation*}
|L(x, \lambda)| \leqslant H_{N}(\delta)|\lambda \mathcal{A}(x)|^{-N} \quad(N=1,2, \ldots) \tag{7.11}
\end{equation*}
$$

for (7.6), where $H_{N}$ are positive constants depending on $\delta$.

For each $\xi$ in $I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$ there exists a positive constant $\delta(\xi)$ such that the curve $\mathcal{C}_{k, 1}(\xi)$ is contained in $\Omega(\delta(\xi))$. We have

$$
\begin{equation*}
z_{k}(s ; \xi) \in \Omega(\delta(\xi)) \quad \text { for } \quad 0 \leqslant s<+\infty \tag{7.12}
\end{equation*}
$$

Hence we have

$$
\begin{equation*}
\left|L\left(z_{k}(s ; \xi), \lambda\right)\right| \leqslant H_{N}(\delta(\xi))\left|\lambda \mathcal{A}\left(z_{k}(s ; \xi)\right)\right|^{-N} \quad(N=1,2, \ldots) \tag{7.13}
\end{equation*}
$$

for $0 \leqslant s<+\infty$ and $|\lambda| \geqslant u_{0} \delta(\xi)^{-\frac{1}{2} m}$, and $|\arg \lambda| \leqslant \theta_{0}(\xi)$.
Assume that a function

$$
\begin{equation*}
v=v(s ; \lambda ; \xi) \tag{7.14}
\end{equation*}
$$

satisfies the differential equation

$$
\begin{equation*}
d v / d s=\dot{z}_{k}(s ; \xi)\left[L(z, \lambda)+\left((-1)^{k} 2 \lambda \mathcal{A}(z)-2 q(z, \lambda)\right) v-\lambda \mathcal{A}(z) v^{2}\right] \tag{7.15}
\end{equation*}
$$

where $z=z_{k}(s ; \xi)$, and the inequalities
hold for $\quad 0 \leqslant s<+\infty, \quad|\lambda| \geqslant M_{0}(\xi), \quad|\arg \lambda| \leqslant \theta_{0}(\xi)$,
where $M_{N}$ are positive constants depending on $\xi$. Then

$$
w_{k, 1}(s, \lambda ; \xi)=(-1)^{k+1} \lambda \mathcal{A}\left(z_{k}(s ; \xi)\right)-\frac{1}{4} R\left(z_{k}(s ; \xi)\right)+\mathfrak{q}\left(z_{k}(s ; \xi), \lambda\right)+\lambda A\left(z_{k}(s ; \xi)\right) v(s, \lambda ; \xi)
$$

satisfies all the conditions of Lemma 4. Hence we shall construct a solution (7.14) of (7.15) which satisfies the conditions (7.16). To do this let us reduce the equation (7.15) to an integral equation
$v(s, \lambda ; \xi)$
$=\int_{+\infty}^{s}\left[L(z, \lambda)-2 \boldsymbol{q}(z, \lambda) v(\sigma, \lambda ; \xi)-\lambda \mathcal{A}(z) v(\sigma, \lambda ; \xi)^{2}\right] \exp \left\{(-1)^{k} 2 \lambda \int_{0}^{s} \mathcal{A}(z) \dot{z} d \tau\right\} \dot{z} d \sigma$,
where $z=z_{k}(\sigma ; \xi), z=z_{k}(\tau ; \xi)$ We shall solve this integral equation using successive approximations.

Let us denote by $F(s, \lambda, v ; \xi)$ the right-hand member of (7.18). Then successive approximations are defined by

$$
\left.\begin{array}{l}
v_{0}(s, \lambda ; \xi)=0 \quad(0 \leqslant s<+\infty)  \tag{7.19}\\
v_{n}(s, \lambda ; \xi)=F\left(s, \lambda, v_{n-1} ; \xi\right) \quad(n \geqslant 1) .
\end{array}\right\}
$$

In order to prove that the approximations are well defined for (7.17) if $M_{0}$ is suf-

$$
\text { SUBDOMINANT SOLUTIONS OF } y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0
$$

ficiently large, we remark that the real part of $(-1)^{k} \lambda \int_{0}^{t} \mathcal{A}(z) \dot{z} d s$ is nondecreasing, where $z=z_{k}(s ; \xi)$. Hence

$$
\begin{equation*}
\left|\exp \left\{(-1)^{k} 2 \lambda \int_{0}^{s} \mathcal{A}(z) \dot{z} d \tau\right\}\right| \leqslant 1 \quad \text { for } \quad s \leqslant \sigma<+\infty . \tag{7.20}
\end{equation*}
$$

On the other hand, we have

$$
\begin{equation*}
|L(z, \lambda)| \leqslant H_{N+2}(\delta(\xi))|\lambda A A(z)|^{-N-2} \tag{7.21}
\end{equation*}
$$

and

$$
\begin{equation*}
|q(z, \lambda)| \leqslant\left|p_{1}(z)\right||\lambda, A(z)|^{-1}+B_{1}(\delta(\xi))|\lambda A(z)|^{-2} \tag{7.22}
\end{equation*}
$$

for $0 \leqslant s<+\infty$ and $|\lambda| \geqslant u_{0} \delta(\xi)^{-\frac{1}{2} m}, \quad|\arg \lambda| \leqslant \theta_{0}(\xi)$. Therefore, if

$$
\begin{equation*}
|v(s, \lambda ; \xi)| \leqslant K\left|\lambda \not A\left(z_{k}(s ; \xi)\right)\right|^{-N} \tag{7.23}
\end{equation*}
$$

for (7.17) and some positive number $K$ and some positive integer $N \geqslant 4$, we have

$$
\begin{align*}
|F(s, \lambda, v ; \xi)| \leqslant & \int_{s}^{+\infty}\left[\left(H_{N+2}(\delta)+2 B_{1}(\delta) K\right)|\lambda A(z)|^{-N-2}\right. \\
& \left.+2 K\left|p_{1}(z)\right||\lambda A(z)|^{-N-1}+K^{2}|\lambda A(z)|^{-2 N+1}\right]|\dot{z}| d \tau \tag{7.24}
\end{align*}
$$

where $z=z_{k}(\tau ; \xi)$.
Notice that, for large values of $s$, we have

$$
z_{k}(s ; \xi)=z_{0}+s \exp (2 k \pi i /(m+2)),
$$

where $z_{0}$ is some constant. Hence, for large values of $s$, we get

$$
\frac{1}{2} s \leqslant\left|z_{k}(s ; \xi)\right| \leqslant 2 s
$$

On the other hand, if $x \in \Omega$ and $|x|$ is sufficiently large, we have

$$
\frac{1}{2}|x|^{\frac{1}{2} m} \leqslant|\mathcal{A}(x)| \leqslant 2|x|^{\frac{1}{2} m} .
$$

Moreover, (2.6) and (2.9) imply that

$$
p_{1}(x)=O\left(x^{-2}\right)
$$

as $x$ tends to $\infty$. Thus, we can conclude that $F$ is well defined for (7.17). Furthermore, if $M_{0}$ is sufficiently large, the function $F$ satisfies the inequality

$$
\begin{equation*}
|F(s, \lambda, v ; \xi)| \leqslant K\left|\lambda A\left(z_{k}(s ; \xi)\right)\right|^{-N} \tag{7.25}
\end{equation*}
$$

for (7.17). Now fix $N$ and $K$, and choose $M_{0}$. Then the successive approximations $v_{n}$ can be well defined for (7.17) and they satisfy the inequalities 17-672909 Acta mathematica 119. Imprimé le 9 février 1968.

$$
\begin{equation*}
\left|v_{n}(s, \lambda ; \xi)\right| \leqslant K \mid \lambda A\left(z_{k}(s ; \xi)| |^{-N}\right. \tag{7.26}
\end{equation*}
$$

for (7.17). The uniform convergence of the sequence $\left\{v_{n}\right\}$ can be proved in a similar manner. Thus we get a solution (7.14) of (7.18) which satisfies the condition (7.16) for a particular $N$. Now it is still necessary to prove that we can choose a constant $M_{0}(\xi)$ independent of $N$. Since $\mathcal{A}(x) \neq 0$ in $\Omega(\delta)$, this can be accomplished by choosing $M_{N}(\xi)$ in a suitable manner. This completes the proof of Lemma 4.

## 8. Proof of Theorem 1

In this section we shall evaluate the quantities $f_{-k}(x, \lambda)$ and $f_{-k}^{\prime}(x, \lambda)$ for each fixed $x$ in $I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$. The quantities $f_{k}$ and $f_{k}^{\prime}$ can be treated in the same manner.

Assume that $\xi$ is given in $I_{2 k} \cup I_{2 k-1} \cup I_{2(k-1)}$. For $0 \leqslant s<+\infty,|\lambda| \geqslant M_{0}(\xi)$, $|\arg \lambda| \leqslant \theta_{0}(\xi)$, let us define a function $G(s, \lambda ; \xi)$ by

$$
G=\left[f_{-k}^{\prime}(z, \lambda)-w_{k, 1}(s, \lambda ; \xi) f_{-k}(z, \lambda)\right] \exp \left\{\int_{0}^{s} w_{k, 1}(\sigma, \lambda ; \xi) \dot{z}_{k}(\sigma ; \xi) d \sigma\right\}
$$

where $z=z_{k}(s ; \xi)$. Then it is easily verified that $d G / d s=0$ except at points of discontinuity of $\dot{z}_{k}(s ; \xi)$. Hence $G$ is piecewise constant with respect to $s$. Since $G$ is continuous in s, this function $G$ must be constant with respect to $s$. This constant can be computed by letting $s$ tend to $+\infty$. First of all, from Lemma 2 we derive

$$
\lim _{s \rightarrow+\infty} z_{k}(s ; \xi)=\infty, \quad \lim _{s \rightarrow+\infty} \arg z_{k}(s ; \xi)=2 \pi k /(m+2)
$$

Since $f_{-k}$ is subdominant for $\lambda>0$ in the sector (1.17-( $\left.-k\right)$ )

$$
|\arg x-2 \pi k /(m+2)|<\pi /(m+2),
$$

then, if $\theta_{0}(\xi)>0$ is small and $|\arg \lambda| \leqslant \theta_{0}(\xi)$, we have

$$
\lim _{s \rightarrow+\infty} f_{-k}\left(z_{k}, \lambda\right)=0 \quad \text { and } \lim _{s \rightarrow+\infty} f_{-k}^{\prime}\left(z_{k}, \lambda\right)=0
$$

Furthermore, if we use the asymptotic representation (1.15-(-k)) of $f_{-k}$ and the inequality (2.18) of Lemma 4, we can prove that

$$
\lim _{s \rightarrow+\infty} w_{k, 1}(s, \lambda ; \xi) f_{-k}\left(z_{k}, \lambda\right)=0
$$

On the other hand, it is easily verified that

$$
\lim _{s \rightarrow+\infty} \arg \left[\mathcal{A}\left(z_{k}\right) \dot{z}_{k}(s ; \xi)\right]=k \pi
$$

since $z_{k}(s ; \xi)=z_{0}+s \exp (2 \pi k i /(m+2))$ for large positive values of $s$, where $z_{0}$ is a certain constant. Thus we get

$$
\lim _{s \rightarrow+\infty} \exp \left\{\int_{0}^{s} w_{k, 1}(\sigma, \lambda ; \xi) \dot{z}_{k}(\sigma ; \xi) d \sigma\right\}=0
$$

Therefore, we have $G(s, \lambda ; \xi)=0$ identically. This in turn implies that
for

$$
\begin{gather*}
f_{-k}^{\prime}(z, \lambda)-w_{k .1}(s, \lambda ; \xi) f_{-k}(z, \lambda)=0 \\
0 \leqslant s<+\infty, \quad|\lambda| \geqslant M_{0}(\xi), \quad|\arg \lambda| \leqslant \theta_{0}(\xi), \tag{8.1}
\end{gather*}
$$

where $z=z_{k}(s ; \xi)$. Consider now the function

$$
H(s, \lambda ; \xi)=f_{-k}\left(z_{k}, \lambda\right) \exp \left\{-\int_{0}^{s} w_{k, 1}(\sigma, \lambda ; \xi) \dot{z}_{k}(\sigma ; \xi) d \sigma\right\}
$$

It is easily verified that $d H / d s=0$ except at points of discontinuity of $\dot{z}_{k}(s ; \xi)$. However $H$ is continuous in $s$. Thus we conclude that $H(s, \lambda ; \xi)$ is independent of $s$. Let us denote by $C(\lambda ; \xi)$ the values of $H$. Then we get

$$
\begin{aligned}
& f_{-k}\left(z_{k}(s ; \xi), \lambda\right)=C(\lambda ; \xi) \exp \left\{\int_{0}^{s} w_{k, 1}(\sigma, \lambda ; \xi) \dot{z}_{k}(\sigma ; \xi) d \sigma\right\} \\
& f_{-k}^{\prime}\left(z_{k}(s ; \xi), \lambda\right)=C(\lambda ; \xi) w_{k, 1}(s, \lambda ; \xi) \exp \left\{\int_{0}^{s} w_{k, 1}(\sigma, \lambda ; \xi) \dot{z}_{k} d \sigma\right\}
\end{aligned}
$$

for (8.1). By putting $s=0$, we get
for

$$
\begin{align*}
& f_{-k}(\xi, \lambda)=C(\lambda ; \xi) \\
& f_{-k}^{\prime}(\xi, \lambda)=C(\lambda, \xi) w_{k, 1}(0, \lambda ; \xi) \\
& |\lambda| \geqslant M_{0}(\xi), \quad|\arg \lambda| \leqslant \theta_{0}(\xi) \tag{8.2}
\end{align*}
$$

We can evaluate the quantity $C(\lambda ; \xi)$ if we compute $\lim _{s \rightarrow+\infty} H(s, \lambda ; \xi)$. To do this, first of all, we shall rewrite the asymptotic representation (1.15-(-k)) of the function $f_{-k}$ using the following identities;

$$
\varrho \omega^{k} \mathcal{A}_{m}\left(\varrho \omega^{k} x, \varrho \omega^{k} a\right)=\left\{\begin{array}{l}
(-1)^{k} \lambda \mathcal{A}_{m}(x, a) \quad(m=\mathrm{odd}), \\
(-1)^{k} \lambda \mathcal{A}_{m}(x, a)+\lambda \alpha_{k}\left\{\frac{1}{x+\varrho^{-1} \omega^{-k}}-\frac{1}{x+1}\right\} \quad(m=\text { even }) .
\end{array}\right.
$$

Let

$$
\gamma(\lambda)=\left\{\begin{array}{l}
\left(\varrho \omega^{-k}\right)^{-\frac{1}{2} m} \quad(m=\text { odd }), \\
\left(\varrho \omega^{-k}\right)^{-\frac{1}{2} m-\lambda \alpha_{k}} \quad(m=\text { even }) .
\end{array}\right.
$$

Then the function $f_{-k}$ has the asymptotic representation

$$
\begin{equation*}
f_{-k}=\gamma(\lambda) x^{-\frac{t^{m}}{}}\left\{1+O\left(x^{-\frac{1}{2}}\right)\right\} \exp \left\{(-1)^{k+1} \lambda \int_{0}^{x} \mathcal{A}_{m}(t, a) d t\right\} \tag{8.3}
\end{equation*}
$$

uniformly on each compact set of the $\lambda$-plane as $x$ tends to infinity in the sector

$$
\begin{equation*}
\left|\arg \varrho+\arg x-\frac{2 \pi k}{(m+2)}\right| \leqslant \frac{3 \pi}{(m+2)}-\delta_{0} \tag{8.4}
\end{equation*}
$$

where the path of integration should be taken in the sector (8.4). In deriving (8.3), we used the identity

$$
\left\{\frac{\varrho \omega^{-k} x+1}{x+1}\right\}^{\lambda \alpha_{k}}=\exp \left[\lambda \alpha_{k} \int_{0}^{x}\left\{\frac{1}{t+\varrho^{-1} \omega^{k}}-\frac{1}{t+1}\right\} d t\right]
$$

Secondly, we remark that

$$
\mathcal{B}(\xi+) / \mathcal{B}(z)=\exp \left[-\frac{1}{4} \int_{0}^{s} R\left(z_{k}(\sigma ; \xi)\right) \dot{z}_{k}(\sigma ; \xi) d \sigma\right]
$$

where $z=z_{k}(s ; \xi)$, and that the integral

$$
\int_{0}^{+\infty} Q_{k, 1}(s, \lambda ; \xi) \dot{z}_{k}(s ; \xi) d s
$$

exists for (8.2).
Finally, we shall write the integral

$$
\int_{0}^{s} \mathcal{A}\left(z_{k}(\sigma ; \xi)\right) \dot{z}_{k}(\sigma ; \xi) d \sigma
$$

in the form

$$
\int_{0}^{s} \mathcal{A}\left(z_{k}(\sigma ; \xi)\right) \dot{z}_{k}(\sigma ; \xi) d \sigma=\int_{0}^{z_{k}(s ; \xi)} \mathcal{A}(t) d t-\int_{0}^{\xi} \mathcal{A}(\tau+) d \tau
$$

where the first integral is taken along a path in the sector (8.4), and the second integral is taken along the real axis. Notice that

$$
\int_{0}^{\infty}\left\{\mathcal{A}_{m}(t, a)-\mathcal{A}(t)\right\} d t
$$

$$
\text { SUBDOMINANT SOLUTIONS OF } y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0
$$

exists if the integral is taken along a path in the sector (8.4). Furthermore, it is easily verified that this integral is equal to

$$
\int_{0}^{+\infty}\left\{\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau+)\right\} d \tau
$$

which is taken along the real axis. This completes the proof of Theorem 1 for $f_{-k}$.

## III. Proof of Theorem 2

## 9. Wronskians

Hereafter we assume that $m$ is even, and that

$$
\begin{equation*}
k=1,2, \ldots, \frac{1}{2} m . \tag{9.1}
\end{equation*}
$$

Let us fix a point $\xi_{h}$ in each interval $I_{h}$, where $h=1,2, \ldots, m-1$.
Lemma 9. If we denote by $\Delta_{k}(\lambda)$ the Wronskian of $f_{k}$ and $f_{-k}$, i.e.

$$
\Delta_{k}(\lambda)=\left|\begin{array}{ll}
f_{k}(x, \lambda) & f_{-k}(x, \lambda) \\
f_{k}^{\prime}(x, \lambda) & f_{-k}^{\prime}(x, \lambda)
\end{array}\right|
$$

we find that

$$
\begin{equation*}
\Delta_{k}(\lambda)=2 i \varrho^{1-2 \lambda \alpha_{k}} \exp \left\{(-1)^{k+1} \lambda L_{k}+O\left(\lambda^{-1}\right)\right\} \tag{9.2}
\end{equation*}
$$

as $\lambda$ tends to $\infty$ in the sector $|\arg \lambda| \leqslant \theta_{0}=\min _{h} \theta_{0}\left(\xi_{h}\right)$, where

$$
\begin{align*}
& L_{k}=2 \sum_{j=1}^{k-1}(-1)^{j+1} \int_{a_{2 j+1}}^{a_{2 j}} \sqrt{|P(\tau)|} d \tau+\int_{0}^{a_{1}}[\mathcal{A}(\tau+)+\mathcal{A}(\tau-)] d t \\
&+\int_{0}^{+\infty}\left[\left(\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau+)\right)+\left(\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau-)\right)\right] d \tau \tag{9.3}
\end{align*}
$$

and $\sqrt{|P(\tau)|} \geqslant 0$, and $\tau$ is a real variable.
Remark. This lemma verifies that $f_{k}$ and $f_{-k}$ are linearly independent for large $\lambda$ in the sector $|\arg \lambda| \leqslant \theta_{0}$.

In order to prove this lemma, let us put $x=\xi_{2 k-1}$ in (4.2). Then we derive

$$
\begin{equation*}
\Delta_{k}(\lambda)=f_{k}(\xi, \lambda) f_{-k}(\xi, \lambda)\left[(-1)^{k+1} \lambda\{\mathcal{A}(\xi+.)-\mathcal{A}(\xi-)\}+O\left(\lambda^{-1}\right)\right] \tag{9.4}
\end{equation*}
$$

from Lemmas 4 and 5 and Theorem 1, where $\xi=\xi_{2 k-1}$. Now, using Theorem 1, we get $f_{k}(\xi, \lambda) f_{-k}(\xi, \lambda)$
$=C_{k}(\lambda) C_{-k}(\lambda) \mathcal{B}(\xi-)^{-1} \mathcal{B}(\xi+)^{-1} \exp \left[(-1)^{k+1} \lambda \int_{0}^{\xi}\{\mathcal{A}(\tau+)+\mathcal{A}(\tau-)\} d \tau+O\left(\lambda^{-1}\right)\right]$
and

$$
\begin{align*}
& C_{k}(\lambda) C_{-k}(\lambda) \\
& \quad=\varrho^{-\frac{1}{2} m-2 \lambda_{k}} \exp \left[(-1)^{k+1} \lambda \int_{0}^{+\infty}\left\{\left(\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau+)\right)+\left(\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau-)\right)\right\} d \tau\right] . \tag{9.6}
\end{align*}
$$

On the other hand, it is easily verified that, for $\xi=\xi_{2 k-1}$, we have

$$
\left.\begin{array}{l}
\mathcal{A}(\xi+)=(-1)^{k+1} i \sqrt{|P(\xi)|},  \tag{9.7}\\
\mathcal{A}(\xi-)=(-1)^{k+1}(-i) \sqrt{|P(\xi)|}, \\
\mathcal{B}(\xi-) \mathcal{B}(\xi+)=\sqrt{|P(\xi)|} .
\end{array}\right\}
$$

Therefore

$$
\mathcal{B}(\xi-)^{-1} \mathcal{B}(\xi+)^{-1}(\mathcal{A}(\xi+)-\mathcal{A}(\xi-))=(-1)^{k+1} 2 i .
$$

Finally, since
and

$$
\mathcal{A}(\tau+)=\mathcal{A}(\tau-)=\mathcal{A}(\tau)=(-1)^{j} \sqrt{\left|P_{(\tau)}\right|} \text { for } \tau \in I_{2 j}
$$

we find that

$$
A(\tau+)=-A(\tau-) \quad \text { for } \quad \tau \in I_{2 j-1}
$$

$$
\int_{\xi}^{a_{1}}[\mathcal{A}(\tau+)+\mathcal{A}(\tau-)] d \tau=2 \sum_{j=1}^{k-1}(-1)^{j} \int_{a_{2 j+1}}^{a_{2 j}} \sqrt{|P(\tau)|} d \tau
$$

Hence

$$
\int_{0}^{\xi}[\mathcal{A}(\tau+)+\mathcal{A}(\tau-)] d \tau=\int_{0}^{a_{1}}[\mathcal{A}(\tau+)+\mathcal{A}(\tau-)] d \tau-2 \sum_{j=1}^{k-1}(-1)^{j} \int_{a_{2 j+1}}^{a_{2 j}} \sqrt{|P(\tau)|} d \tau .
$$

This completes the proof of Lemma 9.
Lemma 10. If we denote by $D_{k, 1}(\lambda)$ and $D_{k, 2}(\lambda)$ the Wronskians of $f_{k}, f_{k+1}$ and $f_{-k}, f_{-k-1}$ respectively, i.e.

$$
D_{k, 1}(\lambda)=\left|\begin{array}{cc}
f_{k} & f_{k+1} \\
f_{k}^{\prime} & f_{k+1}^{\prime}
\end{array}\right|, \quad D_{k, 2}(\lambda)=\left|\begin{array}{ll}
f_{-k} & f_{-k-1} \\
f_{-k}^{\prime} & f_{-k-1}^{\prime}
\end{array}\right|
$$

then we find that

$$
D_{k, 1}(\lambda)=2 \varrho \omega^{k-m / 4} \omega^{\lambda \alpha_{k}}, \quad D_{k, 2}(\lambda)=2 \varrho \omega^{-k+m / 4} \omega^{-\lambda \alpha_{k}}
$$

We shall compute $D_{k, 1}(\lambda)$. The Wronskian $D_{k .2}(\lambda)$ is computed in the same manner. To compute $D_{k, 1}$, we use the asymptotic representations (1.15) and the identity

$$
\varrho \omega^{k} \mathcal{A}_{m}\left(\varrho \omega^{k} x, \varrho \omega^{k} a\right)=(-1)^{k} \lambda \mathcal{A}_{m}(x, a)+\lambda \alpha_{k}\left\{\frac{1}{x+\varrho^{-1} \omega^{-k}}-\frac{1}{x+1}\right\} .
$$

We get

$$
\begin{gathered}
\text { SUBDOMINANT SOLUTIONS OF } y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0 \\
\begin{aligned}
D_{k, 1}(\lambda) & =\left(\varrho \omega^{k}\right)^{-\frac{1}{2} m}\left(\varrho \omega^{k+1}\right)^{-\frac{1}{m}} \lambda\left|\begin{array}{c}
1 \\
(-1)^{k+1}(-1)^{k+2}
\end{array}\right| \lim _{x \rightarrow \infty}\left(\frac{\varrho \omega^{k} x+1}{x+1}\right)^{-\lambda \alpha_{k}}\left(\frac{\varrho \omega^{k+1} x+1}{x+1}\right)^{\lambda \alpha_{k}} \\
= & (-1)^{k} 2 \varrho \omega^{-m(2 k+1) / 4} \omega^{\lambda \alpha_{k}}=2 \varrho \omega^{k-m / 4} \omega^{\lambda \alpha_{k}},
\end{aligned}
\end{gathered}
$$

where $x$ tends to infinity along the direction $\arg x=-2 \pi k / m+2$. Here we used the identity $\alpha_{k+1}=-\alpha_{k}$.

Lemma 11. Let
and

$$
\begin{align*}
& \Delta_{k, 1}(\lambda)=\left|\begin{array}{ll}
f_{-k}(x, \lambda) & f_{k+1}(x, \lambda) \\
f_{-k}^{\prime}(x, \lambda) & f_{k+1}^{\prime}(x, \lambda)
\end{array}\right| \\
& \Delta_{k, 2}(\lambda)=\left|\begin{array}{ll}
f_{k}(x, \lambda) & f_{-k-1}(x, \lambda) \\
f_{k}^{\prime}(x, \lambda) & f_{-k-1}^{\prime}(x, \lambda)
\end{array}\right| . \\
& \Delta_{k, 1}(\lambda)=J_{l}(\lambda) D_{k, 1}(\lambda)+E_{k, 1}(\lambda),  \tag{9.8}\\
& \Delta_{k, 2}(\lambda)=D_{k, 2}(\lambda) / J_{k}(\lambda)+E_{k, 2}(\lambda),  \tag{9.9}\\
& J_{k}(\lambda)=(-1)^{k}\left(\omega^{2 k}\right)^{\frac{{ }^{2}}{m+\lambda} \alpha_{k}} \exp \left\{M_{k}(\lambda)\right\},  \tag{9.10}\\
& M_{k}(\lambda)=(-1)^{k+1} 2 i \lambda \sum_{j=1}^{k}(-1)^{j} \int_{a_{2 j}}^{a_{2 j-1}} \sqrt{|P(\tau)|} d \tau \\
& +\int_{0}^{+\infty} Q_{k, 2}\left(s, \lambda ; \xi_{2 k}\right) \xi_{k}\left(s ; \xi_{2 k}\right) d s-\int_{0}^{+\infty} Q_{k, 1}\left(s, \lambda ; \xi_{2 k}\right) \dot{z}_{k}\left(s ; \xi_{2 k}\right) d s, \tag{9.11}
\end{align*}
$$

Then
and
where
and for $\operatorname{Re} \lambda \geqslant M_{0}, \quad|\operatorname{Im} \lambda| \leqslant r_{0}$, we have

$$
\begin{equation*}
\left|E_{k, n}(\lambda)\right| \leqslant K_{0}\left|\varrho \exp \left\{-2 \lambda \int_{\xi_{2 k}}^{a_{2 k}} \sqrt{|P(\tau)|} d t\right\}\right| \quad(n=1,2) \tag{9.12}
\end{equation*}
$$

if $K_{0}, M_{0}$ and $1 / r_{0}$ are sufficiently large positive numbers.
We shall prove this lemma in a number of steps. First of all, we shall prove the following lemma.

Lemma 11. We have the identity

$$
\begin{equation*}
f_{-k}\left(\xi_{2 k}, \lambda\right)=J_{k}(\lambda) f_{k}\left(\xi_{2 k}, \lambda\right) \tag{9.13}
\end{equation*}
$$

In fact, by using Theorem 1, we get

$$
\begin{equation*}
\frac{f_{-k}(\xi, \lambda)}{f_{k}(\xi, \lambda)}=\frac{C_{-k}(\lambda) \mathcal{B}(\xi+)^{-1} \exp \left[(-1)^{k+1} \lambda \int_{0}^{\xi} \mathcal{A}(\tau+) d \tau-\int_{0}^{+\infty} Q_{k, 1}(s, \lambda ; \xi) \dot{z}_{k}(s ; \xi) d s\right]}{C_{k}(\lambda) \mathcal{B}(\xi-)^{-1} \exp \left[(-1)^{k+1} \lambda \int_{0}^{\xi} \mathcal{A}(\tau-) d \tau-\int_{0}^{+\infty} Q_{k .2}(s, \lambda ; \xi) \dot{\zeta}_{k}(s ; \xi) d s\right]} \tag{9.14}
\end{equation*}
$$

where $\xi=\xi_{2 k}$. Furthermore,

$$
\int_{0}^{+\infty}\left[\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau+)\right] d \tau-\int_{0}^{+\infty}\left[\mathcal{A}_{m}(\tau, a)-\mathcal{A}(\tau-)\right] d \tau=-\int_{0}^{+\infty}[\mathcal{A}(\tau+)-\mathcal{A}(\tau-)] d \tau
$$

On the other hand, it is easily seen that
and

$$
\mathcal{A}(\tau+)=\mathcal{A}(\tau-)=\mathcal{A}(\tau)=(-1)^{j} \sqrt{|P(\tau)|} \quad \text { for } \quad \tau \in I_{2 j}
$$

$$
\mathcal{A}(\tau+)=-\mathcal{A}(\tau-)=(-1)^{j-1} i \sqrt{|P(\tau)|} \text { for } \tau \in I_{2 j-1}
$$

Hence

$$
\int_{\xi_{2 k}}^{+\infty}[\mathcal{A}(\tau+)-\mathcal{A}(\tau-)] d \tau=2 i \sum_{j=1}^{k}(-1)^{j-1} \int_{a_{2 j}}^{a_{2 j-1}} \sqrt{|P(\tau)|} d \tau .
$$

Notice also that

$$
\arg \mathcal{B}\left(\xi_{2 k}+\right)=2 k(\pi / 4)=\frac{1}{2} k \pi, \quad \arg \mathcal{B}\left(\xi_{2 k}-\right)=2 k(-\pi / 4)=-\frac{1}{2} k \pi .
$$

Therefore

$$
\mathcal{B}\left(\xi_{2 k}-\right) / \mathcal{B}\left(\xi_{2 k}+\right)=(-1)^{k} .
$$

This completes the proof of Lemma 12.

Lemma 13. We have the inequality

$$
\begin{equation*}
\left|Q_{k, 1}\left(0, \lambda ; \xi_{2 k}\right)-Q_{k, 2}\left(0, \lambda ; \xi_{2 k}\right)\right| \leqslant K_{1}\left|\lambda \exp \left\{-2 \lambda \int_{\xi}^{a_{2 k}} \sqrt{|P(\tau)|} \cdot d \tau\right\}\right|, \tag{9.15}
\end{equation*}
$$

where $K_{1}$ is a suitable positive constant.
In order to prove this lemma, we shall again compute the Wronskian $\Delta_{k}(\lambda)$ of $f_{k}$ and $f_{-k}$ by putting $x=\xi_{2 k}$ in (4.2). Since $\mathcal{A}(\xi+)=\mathcal{A}(\xi-)$ for $\xi \in I_{2 k}$, we have

$$
\begin{equation*}
\Delta_{k}(\lambda)=\left\{\boldsymbol{Q}_{k, \mathbf{1}}(0, \lambda ; \xi)-\boldsymbol{Q}_{k, 2}(0, \lambda ; \xi)\right\} f_{k}(\xi, \lambda) f_{-k}(\xi, \lambda) \tag{9.16}
\end{equation*}
$$

where $\xi=\xi_{2 k}$. Hence

$$
\begin{equation*}
\boldsymbol{Q}_{k, 1}\left(0, \lambda ; \xi_{2 k}\right)-Q_{k, 2}\left(0, \lambda ; \xi_{2 k}\right)=\Delta_{k}(\lambda) / f_{k}\left(\xi_{2 k}, \lambda\right) f_{-k}\left(\xi_{2 k}, \lambda\right) \tag{9.17}
\end{equation*}
$$

On the other hand, from the formula (9.4) we derive

$$
\begin{aligned}
& \text { SUBDOMINANT SOLUTIONS OF } y^{\prime \prime}-\lambda^{2}\left(x-a_{1}\right)\left(x-a_{2}\right) \ldots\left(x-a_{m}\right) y=0 \\
& \begin{aligned}
&\left|Q_{k, 1}\left(0, \lambda ; \xi_{2 k}\right)-Q_{k, 2}\left(0, \lambda ; \xi_{2 k}\right)\right| \leqslant K|\lambda|\left|\frac{f_{k}\left(\xi_{2 k-1}, \lambda\right) f_{-k}\left(\xi_{2 k-1}, \lambda\right)}{f_{k}\left(\xi_{2 k}, \lambda\right) f_{-k}\left(\xi_{2 k}, \lambda\right)}\right| \\
&=|\lambda| \exp \left[(-1)^{k+1} \lambda \int_{\xi_{2 k}}^{\xi_{2 k-1}}(\mathcal{A}(\tau+)+\mathcal{A}(\tau-)) d \tau+O(1)\right] \mid
\end{aligned}
\end{aligned}
$$

where $K$ is a suitable positive constant. Since
and

$$
\mathcal{A}(\tau+)=-\mathcal{A}(\tau-) \text { for } \tau \in I_{2 k-1}
$$

$$
\mathcal{A}(\tau+)=\mathcal{A}(\tau-)=(-1)^{k} V \sqrt{|P(\tau)|} \quad \text { for } \quad \tau \in I_{2 k}
$$

we have

$$
\int_{\xi_{2 k}}^{\xi_{2 k-1}}(\mathcal{A}(\tau+)+\mathcal{A}(\tau-)) d \tau=\int_{\xi_{2 k}}^{a_{2 k}}(\mathcal{A}(\tau+)+\mathcal{A}(\tau-)) d \tau=(-1)^{k} 2 \int_{\xi_{2 k}}^{a_{2 k}} \sqrt{|P(\tau)|} d \tau
$$

This completes the proof of Lemma 13.

## Lemma 14. We find that

where

$$
\begin{equation*}
f_{-k}^{\prime}\left(\xi_{2 k}, \lambda\right)=J_{k}(\lambda) f_{k}^{\prime}\left(\xi_{2 k}, \lambda\right) \quad\left\{1+F_{k}(\lambda)\right\} \tag{9.18}
\end{equation*}
$$

$$
\left|F_{k}(\lambda)\right| \leqslant K_{2}\left|\exp \left\{-2 \lambda \int_{\xi_{2 k}}^{a_{2 k}} \sqrt{|P(\tau)|} d \tau\right\}\right|
$$

and $K_{2}$ is a suitable positive constant.
In fact, we have

$$
f_{-k}^{\prime}(\xi, \lambda)=\left[(-1)^{k+1} \lambda \mathcal{A}(\xi+)-\frac{1}{4} R(\xi)+Q_{k, 1}(0, \lambda ; \xi)\right] J_{k}(\lambda) f_{k}(\xi, \lambda)
$$

where $\xi=\xi_{2 k}$. Here we used Theorem 1 and Lemma 12. Then using $\mathcal{A}\left(\xi_{2 k}+\right)=\mathcal{A}\left(\xi_{2 k}-\right)$ and Lemma 13, we can prove (9.18) without any difficulty.

Now we shall prove (9.8). The formula (9.9) can be proved in the same manner. Letting $x=\xi_{2 k}$ in (4.5) and using Lemmas 12 and 14, we get

$$
\Delta_{k, 1}(\lambda)=J_{k}(\lambda) D_{k, 1}(\lambda)-J_{k}(\lambda) f_{k}^{\prime}\left(\xi_{2 k}, \lambda\right) f_{k+1}\left(\xi_{2 k}, \lambda\right) F_{k}(\lambda)
$$

It is easily seen that $J_{k}(\lambda)$ is bounded with respect to $\lambda$ for $\operatorname{Re} \lambda \geqslant M_{0},|\operatorname{Im} \lambda| \leqslant r_{0}$. On the other hand, Theorem 1 implies that

$$
f_{k}^{\prime}(\xi, \lambda) f_{k+1}(\xi, \lambda)=\left\{(-1)^{k+1} \lambda \mathcal{A}(\xi-)-\frac{1}{4} R(\xi)+Q_{k, 2}(0, \lambda ; \xi)\right\} f_{k}(\xi, \lambda) f_{k+1}(\xi, \lambda)
$$

and

$$
\left.f_{k}(\xi, \lambda) f_{k+1} \xi, \lambda\right)=C_{k}(\lambda) C_{k+1}(\lambda) \mathcal{B}(\xi-)^{-2} \exp \left[O\left(\lambda^{-1}\right)\right]=O\left(\varrho^{-\frac{1}{2} m}\right)
$$

where $\xi=\xi_{2 k}$. This completes the proof of Lemma 11.

## 10. Proof of Theorem 2

We shall use the notations which were introduced in Section 4. We have already shown that large positive eigenvalues of Problem ( $P$ ) are zeros of the function (4.16). Using Lemmas 10 and 11, and the formula (4.9), we find that

$$
H_{1}(\lambda)\binom{D_{0.2}(\lambda)}{-D_{0.1}(\lambda)}=J_{1}^{-1}\left(D_{0.2}-J_{1} D_{0.1}\right)\binom{D_{1.2}}{-J_{1} D_{1.1}}+\varrho^{2}\binom{U_{1}(\lambda)}{V_{1}(\lambda)}
$$

where

$$
\left|U_{1}(\lambda)\right|+\left|V_{1}(\lambda)\right|=O\left(\exp \left[-2 \lambda \int_{\xi_{2}}^{a_{2}} \sqrt{|P(\tau)|} d \tau\right]\right)
$$

for $\operatorname{Re} \lambda \geqslant M_{0},|\operatorname{Im} \lambda| \leqslant r_{0}$. In the same manner we have

$$
\begin{aligned}
& \qquad H_{2} H_{1}\binom{D_{0.2}}{-D_{0.1}}=\left(J_{1} J_{2}\right)^{-1}\left(D_{0.2}-J_{1} D_{0.1}\right)\left(D_{1.2}-J_{1} J_{2} D_{1,1}\right)\binom{D_{2.2}}{-J_{2} D_{2,1}}+\varrho^{3}\binom{U_{2}(\lambda)}{V_{2}(\lambda)} \\
& \text { where } \quad\left|U_{2}(\lambda)\right|+\left|V_{2}(\lambda)\right|=O\left(\sum_{j=1}^{2} \exp \left[-2 \lambda \int_{\xi_{2 j}}^{a_{2 j}} \sqrt{|P(\tau)|} d \tau\right]\right)
\end{aligned}
$$

for $\operatorname{Re} \lambda \geqslant M_{0},|\operatorname{Im} \lambda| \leqslant r_{0}$. In general, we have

$$
\begin{align*}
& H(\lambda)\binom{D_{0.2}(\lambda)}{-D_{0.1}(\lambda)}=\left\{\prod_{h-1}^{\frac{1}{m-1}} J_{h}(\lambda)\right\}^{-1} \prod_{k=1}^{\frac{1}{m} m-1}\left(D_{k-1,2}-J_{k-1} J_{k} D_{k-1.1}\right)\binom{D_{\frac{1}{2} m-1,2}}{-J_{\frac{1}{2} m-1} D_{\frac{1}{2} m-1,1}} \\
& +\varrho^{\frac{1}{2} m}\binom{U_{\frac{1}{2} m-1}(\lambda)}{V_{\frac{1}{2} m-1}(\lambda)}, \tag{10.1}
\end{align*}
$$

where

$$
\begin{equation*}
\left|U_{\frac{1}{2} m-1}(\lambda)\right|+\left|V_{\frac{1}{2} m-1}(\lambda)\right|=O\left(\sum_{j=1}^{\frac{1}{2} m-1} \exp \left[-2 \lambda \int_{\xi_{2 j}}^{a_{2 j}} \sqrt{|P(\tau)|} d \tau\right]\right) \tag{10.2}
\end{equation*}
$$

for $\operatorname{Re} \lambda \geqslant M_{0},|\operatorname{Im} \lambda| \leqslant r_{0}$. Thus we have proved that almost all positive eigenvalues of Problem ( $P$ ) must satisfy the equation

$$
\begin{equation*}
\left(D_{\frac{1}{2} m, 1} D_{\frac{1}{2} m-1,2}-J_{\frac{1}{\frac{1}{2} m-1}} D_{\frac{1}{2} m-1,1} D_{\frac{1}{2} m, 2}\right) \prod_{k=1}^{\frac{1}{\frac{1}{m}-1}}\left(D_{k-1,2}-J_{k-1} J_{k} D_{k-1,1}\right)=\lambda U(\lambda) \tag{10.3}
\end{equation*}
$$

where $J_{0}(\lambda) \equiv 1$ and

$$
\begin{equation*}
U(\lambda)=O(\exp (-\beta \lambda)) \tag{10.4}
\end{equation*}
$$

as $\lambda$ tends to $\infty$ in the strip $\operatorname{Re} \lambda \geqslant 0,|\operatorname{Im} \lambda| \leqslant r_{0}$, the quantity $\beta$ being a suitable positive constant.

Now by using Lemma 10 and the formula (9.10), we shall rewrite the equation (10.3). In fact, we get

$$
D_{k, 1}(\lambda)=\omega^{2 \lambda \alpha_{k}+2 k-m / 2} D_{k, 2}(\lambda)
$$

Hence, if we use the identity $\alpha_{k+1}=-\alpha_{k}$, we get

$$
\begin{equation*}
D_{0.2}(\lambda)-J_{1}(\lambda) D_{0.1}(\lambda)=D_{0.2}(\lambda)\left\{1+\exp \left[M_{1}(\lambda)\right]\right\} \tag{10.5}
\end{equation*}
$$

and

$$
\begin{align*}
& D_{k-1,2}(\lambda)-J_{k-1}(\lambda) J_{k}(\lambda) D_{k-1.1}(\lambda) \\
& \quad=D_{k-1,2}(\lambda)\left\{1+\exp \left[M_{k-1}(\lambda)+M_{k}(\lambda)\right]\right\} \quad(k=2,3, \ldots,-1+m / 2) \tag{10.6}
\end{align*}
$$

In the same manner, if we use the identities

$$
\alpha_{\frac{1}{2} m}=-\alpha_{\frac{1}{2} m-1}, \quad\left(\omega^{m+2}\right)^{\lambda \pi_{\frac{1}{2} m}}=\exp \left[2 \pi i \lambda \alpha_{\frac{1}{2} m}\right], \quad\left(\omega^{m+2}\right)^{\frac{1}{m}}=(-1)^{\frac{1}{\frac{1}{2}} m}
$$

we have

$$
\begin{align*}
& D_{\frac{1}{2} m, 1}(\lambda) D_{\frac{1}{2} m-1,2}(\lambda)-J_{\frac{1}{\frac{1}{2} m-1}}(\lambda) D_{\frac{1}{2} m-1,2}(\lambda) D_{\frac{1}{2} m, 2}(\lambda) \\
& \quad=D_{\frac{1}{2} m, 1}(\lambda) D_{\frac{1}{2} m-1,2}(\lambda)\left\{1+\exp \left[-2 \pi i \lambda \alpha_{\frac{1}{2} m}+M_{\frac{1}{2} m-1}(\lambda)\right]\right\} \tag{10.7}
\end{align*}
$$

Thus we have proved that almost all positive eigenvalues must satisfy the equation

$$
\begin{align*}
& \left\{1+\exp \left[-2 \pi i \lambda \alpha_{\frac{1}{3} m}+M_{\frac{1}{2} m-1}\right]\right\}\left(1+\exp \left[M_{1}\right]\right) \prod_{k=2}^{\frac{1}{2} m-1}\left(1+\exp \left[M_{k-1}+M_{k}\right]\right)=V(\lambda),  \tag{10.8}\\
& \text { where } \quad V(\lambda)=O(\exp (-\beta \lambda))
\end{align*}
$$

as $\lambda$ tends to $\infty$ in the strip $\operatorname{Re} \lambda \geqslant 0,|\operatorname{Im} \lambda| \leqslant r_{0}$.
Finally, using the definitions (9.11) of the quantities $M_{k}(\lambda)$, we shall derive Theorem 2. It is easily seen that

$$
M_{1}(\lambda)=-2 i \lambda \int_{a_{2}}^{a_{1}} \sqrt{|P(\tau)|} d \tau+\int_{0}^{+\infty} Q_{1.2}\left(s, \lambda ; \xi_{2}\right) \dot{\zeta}_{1}\left(s ; \xi_{2}\right) d s-\int_{0}^{+\infty} Q_{1,1}\left(s, \lambda ; \xi_{2}\right) \dot{z}_{1}\left(s ; \xi_{2}\right) d s
$$

and

$$
\begin{aligned}
M_{k-1}(\lambda)+ & M_{k}(\lambda)=-2 i \lambda \int_{a_{2 k}}^{a_{2 k-1}} \sqrt{|P(\tau)|} d \tau+\int_{0}^{+\infty} a_{k-1,2}\left(s, \lambda ; \xi_{2(k-1)}\right) \dot{\xi}_{k-1}\left(s ; \xi_{2(k-1)}\right) d s \\
& -\int_{0}^{+\infty} a_{k-1,1}\left(s, \lambda ; \xi_{2(k-1)}\right) \dot{z}_{k-1}\left(s ; \xi_{2(k-1)}\right) d s+\int_{0}^{+\infty} Q_{k, 2}\left(s, \lambda ; \xi_{2 k}\right) \dot{\xi}_{k}\left(s ; \xi_{2 k}\right) d s \\
& -\int_{0}^{+\infty} Q_{k, 1}\left(s, \lambda ; \xi_{2 k}\right) \dot{z}_{k}\left(s ; \xi_{2 k}\right) d s \quad(k=2,3 \ldots,-1+m / 2)
\end{aligned}
$$

On the other hand, notice that

$$
2 \pi i b_{1+m / 2}(a)=\int \mathcal{A}(t) d t
$$

where the integration should be taken along a large circle in the counterclockwise sense. Hence we have

$$
2 \pi i b_{1+m / 2}(a)=\int_{a_{m}}^{a_{1}}[\mathcal{A}(\tau-)-\mathcal{A}(\tau+)] d \tau=2 i \sum_{j=1}^{\frac{1}{j} m}(-1)^{j} \int_{a_{2 j}}^{a_{2} j-1} \sqrt{|P(\tau)|} d \tau .
$$

Therefore

$$
-2 \pi i \lambda \alpha_{m / 2}=2 \pi i(-1)^{1+m / 2} \lambda b_{1+m / 2}(a)=2 i(-1)^{1+m / 2} \lambda \sum_{j=1}^{\frac{t}{2}}(-1)^{j} \int_{a_{2 j}}^{a_{2 j-1}} \sqrt{|P(\tau)|} d \tau
$$

Thus we have

$$
\begin{aligned}
&-2 \pi i \lambda \alpha_{m / 2}+M_{-1+m / 2}(\lambda) \\
&=-2 i \lambda \int_{a_{m}}^{a_{m-1}} \sqrt{|P(\tau)|} d \tau+\int_{0}^{+\infty} Q_{\frac{1}{2} m-1,2}\left(s, \lambda ; \xi_{m-2}\right) \dot{\xi}_{\frac{1}{2} m-1}\left(s ; \xi_{m-2}\right) d s \\
&-\int_{0}^{+\infty} Q_{\frac{1}{2} m-1,1}\left(s, \lambda ; \xi_{m-2}\right) \dot{z}_{\frac{1}{2} m-1}\left(s ; \xi_{m-2}\right) d s
\end{aligned}
$$

For each $k=1,2, \ldots,-1+m / 2$, let $\Gamma_{k}$ be a circle which encircles the points $a_{1}, a_{2}, \ldots, a_{2 k}$, while the points $a_{2 k+1}, \ldots, a_{m}$ lie outside the curve $\Gamma_{k}$. Then by using Lemmas 4 and 5 and the definitions (3.4) and (3.5) of $Q_{k, 1}$ and $Q_{k, 2}$, we find that

$$
\begin{aligned}
M_{1}(\lambda) \cong & -2 i \lambda \int_{a_{1}}^{a_{2}} \sqrt{|P(\tau)|} \dot{d} \tau+\sum_{n=1}^{\infty} \lambda^{-n} \int_{\Gamma_{1}} \mathcal{A}(x)^{-n} p_{n}(x) d x \\
M_{k-1}(\lambda)+M_{k}(\lambda) \cong & -2 i \lambda \int_{a_{2 k}}^{a_{2 k-1}} \sqrt{|P(\tau)|} d \tau+\sum_{n=1}^{\infty} \lambda^{-n}\left\{\int_{\Gamma_{k}}\left[(-1)^{k+1} A(x)\right]^{-n} p_{n}(x) d x\right. \\
& \left.+(-1)^{n} \int_{\Gamma_{k-1}}\left[(-1)^{k+1} \mathcal{A}(x)\right]^{-n} p_{n}(x) d x\right\} \quad(k=2,3, \ldots,-1+m / 2)
\end{aligned}
$$

and
$-2 \pi i \lambda \alpha_{m / 2}+M_{-1+m / 2}(\lambda) \cong-2 i \lambda \int_{a_{m}}^{a_{m-1}} \sqrt{|P(\tau)|} d \tau+\sum_{n=1}^{\infty} \lambda^{-n} \int_{\Gamma_{\frac{1}{2} m-1}}\left[(-1)^{m / 2} \mathcal{A}(x)\right]^{-n} p_{n}(x) d x$
as $\lambda$ tends to $\infty$, where $\cong$ indicates an asymptotic expansion in powers of $\lambda^{-1}$ in the usual sense, and every integral should be taken in the counterclockwise sense. Then if we use (4.20), Theorem 2 can easily be proved. In fact,

$$
\begin{aligned}
& \int_{\Gamma_{k}} \mathcal{A}(x)^{-n} p_{n}(x) d x=0 \text { for even } n \\
& \int_{\Gamma}-\int_{\Gamma_{k-1}}=\int_{\gamma_{k}}, \quad \int_{\Gamma_{\frac{1}{2} m-1}}=-\int_{\gamma_{\frac{2}{2} m}}
\end{aligned}
$$

Finally we shall prove (4.20). Put
and

$$
\begin{aligned}
& r(x, \lambda)=\lambda A(x)+\sum_{n=1}^{\infty}(\lambda A(x))^{-2 n+1} p_{2 n-1}(x) \\
& s(x, \lambda)=-\frac{1}{4} R(x)+\sum_{n=1}^{\infty}(\lambda \mathcal{A}(x))^{-2 n} p_{2 n}(x) .
\end{aligned}
$$

Then the formal series $r+s$ is a formal solution of the Riccati equation (2.3). Since $r$ contains only odd powers of $\lambda$, we get

$$
\begin{equation*}
r^{\prime}+2 r s=0 . \tag{10.10}
\end{equation*}
$$

This formal identity was obtained by N. Fröman [5]. It can be written as

$$
\begin{equation*}
s=-\frac{1}{2}(\log r)^{\prime} \tag{10.11}
\end{equation*}
$$

Hence we get the formal identity

$$
\begin{equation*}
\int_{\gamma_{k}} s(x, \lambda) d x=0 \quad(\bmod \pi i), \tag{10.12}
\end{equation*}
$$

which completes the proof of (4.20).

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