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# On the Location of Zeros of Solutions of Non-homogeneous Linear Differential Equations. 

Steven B. Bank (*)


#### Abstract

For any non-homogeneous linear differential equation of arbitrary order, having rational functions for coefficients, we determine precisely, those rays $\operatorname{Arg} z=\phi$ which have the property that for any $\varepsilon>0$, there is a solution of the equation which has infinitely many zeros tending to $\infty$ in the sector $|\operatorname{Arg} z-\phi|<\varepsilon$. We remark that it is possible for there to be infinitely many such for a given equation.


## 1. Introduction.

For a non-homogeneous linear differential equation of arbitrary order $n$, having polynomial coefficients,

$$
\begin{equation*}
w^{(n)}+R_{n-1}(z) w^{(n-1)}+\ldots+R_{0}(z) w=Q(z) \tag{1.1}
\end{equation*}
$$

the following result was proved in [2] concerning the possible location of the zeros of the solutions of (1.1):

Theorem 1.1. Given an equation (1.1) where $n \geqslant 2$, and where $R_{0}(z), \ldots, R_{n-1}(z)$, and $Q(z)$ are polynomials with $Q(z) \not \equiv 0$. Then there exist a nonnegative integer $t$ and real numbers $s_{0}, s_{1}, \ldots, s_{t+1}$ with $-\pi=s_{0}<s_{1} \ldots<s_{t+1}=\pi$, such that for each $k \in\{0,1, \ldots, t\}$, one of the following two properties holds:
(a) For any $\theta \in\left(s_{k}, s_{k+1}\right)$ and any $\varepsilon>0$, there is a solution of (1.1) having infinitely many zeros in $|\operatorname{Arg} z-\theta|<\varepsilon$.
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(b) For any $\varepsilon>0$, any solution of (1.1) can have at most finitely many zeros in the closed sector $s_{k}+\varepsilon \leqslant \operatorname{Arg} z \leqslant s_{k+1}-\varepsilon$.

We remark that the results in [2] (which are reviewed in § 3 below for the reader's convenience) provide a method for the explicit calculation of the numbers $s_{k}$ from the equation, and also for determining which of the two possibilities $(a)$ or $(b)$ in Theorem 1.1 holds for a given $k$. (Examples produced in [2] show that for a given equation (1.1), it is possible for property ( $a$ ) to hold for some $k$, while property (b) holds for other $k$. For example, if we consider the third-order equation,

$$
\begin{equation*}
w^{\prime \prime \prime}+\left(z^{2}+3\right) w^{\prime \prime}+\left(2 z^{2}+z+3\right) w^{\prime}+\left(z^{2}+z+2\right) w=Q(z) \tag{1.2}
\end{equation*}
$$

where $Q(z)$ is any polynomial which is not identically zero, it is shown in $[2 ; \S 9]$ that the set $\left\{s_{0}, \ldots, s_{t+1}\right\}$ consists of $s_{0}=-\pi, s_{1}=-5 \pi / 6$, $s_{2}=-\pi / 2, s_{3}=-\pi / 6, s_{4}=\pi / 6, s_{5}=\pi / 2, s_{6}=5 \pi / 6$, and $s_{7}=\pi$, and that property ( $a$ ) holds for $k=0,1,5,6$, while property ( $b$ ) holds for $k=2,3,4$.) We also remark that analogous results to Theorem 1.1 hold for equations (1.1) having more general coefficients than polynomials, such as rational functions or algebraic functions (see[2; Theorem 6.1] or §3 below).

For a given equation (1.1) and a value $s_{k}$ (where $k \in\{1, \ldots, t\}$ ), it is obvious that if property ( $a$ ) holds for either $k$ or $k-1$, then for any $\varepsilon>0$, there is a solution of (1.1) having infinitely many zeros in $\left|\operatorname{Arg} z-s_{k}\right|<\varepsilon$. (This follows easily by applying property (a) to a ray $\operatorname{Arg} z=\theta$ where $\theta$ is sufficiently close to $s_{k}$.) However, if for this given $s_{k}$, property (b) holds for both $k$ and $k-1$, then for a given $\varepsilon>0$, no information is provided by Theorem 1.1 on whether or not there exists a solution of (1.1) which has infinitely many zeros in $\left|\operatorname{Arg} z-s_{k}\right|<\varepsilon$. If for some $\varepsilon>0$ no such solution exists, then $s_{k}$ is really extraneous since property (b) could then be improved to actually assert that for any $\varepsilon>0$, any solution of (1.1) can have at most finitely many zeros on $s_{k-1}+\varepsilon \leqslant \operatorname{Arg} z \leqslant s_{k+1}-\varepsilon$. In fact, it is easy (see §3) to give simple examples of equations (1.1) which possess such extraneous $s_{k}$. In this paper, we develop a simple method for determining whether or not a given $s_{k}$ is extraneous, and this method is given in Theorem 4.1 in §4. In order to develop this method, we must review how the $s_{k}$ are produced and this is done in $\S 3$. There are several reasons why methods in [2] can produce extraneous $s_{k}$. One reason is that we first deal with the homogeneous equation corresponding to (1.1), and we may find an $s_{k}$ for which the homogeneous equation has solutions $f \not \equiv 0$ having infinitely many zeros around $\operatorname{Arg} z=s_{k}$ but for which $s_{k}$ is extraneous for the non-homogeneous equation. (This is the case in the example in §3). Another reason for extraneous $s_{k}$ is that our asymptotic existence theo-
rem produces solutions having prescribed asymptotic behaviour in sectors, and the extraneous $s_{k}$ can occur as the boundary rays of these sectors. In this paper (see Theorem 6.1 in $\S 6$ below), we deal with this problem of boundary rays by developing a process of continuing the solutions over these rays, while at the same time preserving the asymptotic properties of the solutions in the enlarged domain. This process is based on Phragmen-Lindelöf principles.

It should be noted that the problem of extraneous rays arose earlier in [5] in connection with the oscillation properties of solutions of homogeneous linear differential equations of arbitrary order $n$, having polynomial coefficients,

$$
\begin{equation*}
w^{(n)}+R_{n-1}(w)^{(n-1)}+\ldots+R_{0}(z) w=0 . \tag{1.3}
\end{equation*}
$$

The problem for (1.3) is considerably simpler than for (1.1) because for homogeneous equations there cannot be the «mixture» of properties (a) and (b) as the following result[5; Theorem 1] shows:

Theorem 1.2. Given an equation (1.3) where $n \geqslant 2$ and where $R_{0}(z), \ldots, R_{n-1}(z)$ are polynomials. Then one of the following two properties holds:
(A) For any $\theta \in(-\pi, \pi]$ and any $\varepsilon>0$, there is a solution $f \not \equiv 0$ of (1.3) having infinitely many zeros in the sector $|\operatorname{Arg} z-\theta|<\varepsilon$.
(B) There exist a nonnegative integer $t$ and real numbers $s_{0}$, $s_{1}, \ldots, s_{t+1}$ with $-\pi=s_{0}<s_{1}<\ldots<s_{t+1}=\pi$, such that for any $\varepsilon>0$, any solution $f \not \equiv 0$ of (1.3) has at most finitely many zeros in the closed sector $s_{k}+\varepsilon \leqslant \operatorname{Arg} z \leqslant s_{k+1}-\varepsilon$ for each $k \in\{0,1, \ldots, t\}$.

It was shown in [1] that a simple example of an equation (1.3) which possesses property ( $A$ ) in Theorem 1.2 is,

$$
\begin{equation*}
w^{(n)}+z^{2} w^{\prime \prime}+z w^{\prime}+w=0 \quad \text { for any } n \geqslant 3 . \tag{1.4}
\end{equation*}
$$

However for second-order equations (1.3), a classical result (see Hille [7; § 5.6] or R. Nevanlinna [9; p. 345] or Wittich [16; p. 282]) shows that all second-order equations (1.3) possess property $(B)$ in Theorem 1.2.

For a given homogeneous equation (1.3), the results in [5] provided a method for determining which property $(A)$ or $(B)$ in Theorem 1.2 holds. If property $(B)$ holds, the results in [5] give a method for explicitly calculating the numbers $s_{1}, s_{2}, \ldots, s_{t}$ and will also determine if an $s_{k}$ is extraneous in the sense that for some $\varepsilon>0$, no solution $f \not \equiv 0$ of (1.3) has infinitely many zeros in $\left|\operatorname{Arg} z-s_{k}\right|<\varepsilon$. The problem of determin-
ing which $s_{k}$ (if any) are extraneous for (1.3) was solved in [5] by also using a continuation process for the solutions. However, for the reasons indicated earlier, the process in [5] for (1.3) is of a much simpler nature than that required in this paper for (1.1). For completeness, the results from[5] are reviewed in § 10 below.

## 2. Concepts from the Strodt theory [10].

In this paper, we will consider equations (1.1) whose coefficients are of a more general nature than just polynomials or rational functions. We will actually treat equations (1.1) whose coefficients belong to a certain type of function field (i.e. a logarithmic differential field of rank zero) consisting of functions which are analytic in some element of a neighborhood system of sectorial regions. This neighborhood system was introduced by W. Strodt in [10; §94], and is denoted $F(a, b)$. It has the following property (see [2; p. 269]):

Lemma 2.1. Let $V$ be an element of $F(a, b)$, and let $\varepsilon>0$ be arbitrary. Then there is a constant $R_{0}(\varepsilon)>0$ such that $V$ contains the set, $a+\varepsilon \leqslant \operatorname{Arg} z \leqslant b-\varepsilon,|z| \geqslant R_{0}(\varepsilon)$.

In a logarithmic differential field of rank zero over $F(a, b)$, every element $f(z)$ is admissible in $F(a, b)$ (i.e. is analytic in some element of $F(a, b))$ and if $f(z)$ is not identically zero, then $f(z)$ is asymptotically equivalent to a function of the form $c z^{\alpha}$ (where $c$ and $\alpha$ are constants with $c$ nonzero and $\alpha$ real) as $z$ tends to infinity over elements of $F(a, b)$. (The concept of a logarithmic differential field of rank zero over $F(a, b)$ was introduced by W. Strodt in [12; p. 244].) We are using the strong relations of asymptotic equivalence, $f \sim g$, and of asymptotic smallness, $f \ll g$, which were introduced by Strodt in [10; §13]. (The relation, $f \ll g$, is a generalization of the usual relation, $f=o(g)$. The strong relation of asymptotic equivalence defined in $[10 ; \S 13]$ is designed to ensure that if $M(z)$ is a nonconstant logarithmic monomial of rank $\leqslant p$ (i.e. a function of the form,

$$
\begin{equation*}
M(z)=K z^{a_{0}}(\log z)^{a_{1}} \ldots\left(\log _{p} z\right)^{a_{p}} \tag{2.1}
\end{equation*}
$$

for real $a_{j}$, and complex $K \neq 0$ ), then $f \sim M$ implies $f^{\prime} \sim M^{\prime}$ in $F(a, b)$ (see [10; §28]). As usual, $z^{\alpha}$ and $\log z$ will denote the principal branches of these functions on $|\operatorname{Arg} z|<\pi)$. If $f \sim M$ in $F(a, b)$ where $M$ is given by (2.1), then we will denote $a_{0}$ by $\delta_{0}(f), a_{1}$ by $\delta_{1}(f)$ etc. If $f \equiv 0$, we will set $\delta_{0}(f)=-\infty$.

The following two facts are proved in [6; p. 309] and[10; § 28] respectively:

Lemma 2.2. Let $f(z)$ be admissible in $F(a, b)$. Then:
(A) If $f \rightarrow 0$ in $F(a, b)$, then $z f^{\prime}(z) \rightarrow 0$ in $F(a, b)$.
(B) If $f \ll 1$ in $F(a, b)$, then $\theta_{j} f \ll 1$ in $F(a, b)$ for each $j=1,2, \ldots$ where $\theta_{j}$ denotes the operator $\theta_{j} f=z(\log z)(\log \log z) \ldots$ $\ldots\left(\log _{j-1} z\right) f^{\prime}(z)$.

We will write $f_{1} \approx f_{2}$ in $F(a, b)$ to mean that $f_{1} \sim c f_{2}$ for some nonzero constant $c$. An admissible function $g(z)$ in $F(a, b)$ is called trivial in $F(a, b)$ if $g \ll z^{-\alpha}$ in $F(a, b)$ for every $\alpha>0$. If $f \sim c z^{-1+d}$ in $F(a, b)$, where $c \neq 0$ and $d>0$, then the indicial function of $f$ is the function $I F(f, \phi)$ defined by,

$$
\begin{equation*}
I F(f, \phi)=\operatorname{Cos}(d \phi+\operatorname{Arg} c) \quad \text { for } a<\phi<b \tag{2.2}
\end{equation*}
$$

(It is obvious that $\operatorname{IF}(f, \phi)$ has at most finitely many zeros on $(a, b)$.) If $g$ is any admissible function in $F(a, b)$, we will denote by $\int g$ any primitive of $g$ in an element of $F(a, b)$. We will require the following two facts (see [2; p. 270]):

Lemma 2.3. Let $f \sim c z^{-1+d}$ in $F(a, b)$, where $c \neq 0$ and $d>0$. If $\left(a_{1}, b_{1}\right)$ is any subinterval of $(a, b)$ on which $\operatorname{IF}(f, \phi)<0$ (respectively, $\operatorname{IF}(f, \phi)>0$ ), then for all real $\alpha, \exp \int f \ll z^{\alpha} \quad$ (respectively, $\left.\exp \int f \gg z^{\alpha}\right)$ in $F\left(a_{1}, b_{1}\right)$.

Lemma 2.4. Let $\alpha=a+b i$ be a complex number. Then for any $\varepsilon>0$, we have $z^{a-\varepsilon} \ll z^{\alpha}$ and $z^{\alpha} \ll z^{a+\varepsilon}$ in $F(-\pi, \pi)$.

We will also require the following facts. The first is obvious and the second follows from [10; Lemma 30]:

Lemma 2.5. (a) If $b$ is a real number, then on $|\operatorname{Arg} z|<\pi$, we have $\left|z^{b i}\right| \leqslant e^{|b| \pi}$ and $\left|z^{b i}\right| \geqslant e^{-|b| \pi}$.
(b) If $f$ is a trivial function in $F(a, b)$, then $f^{\prime}$ is also a trivial function in $F(a, b)$.

## 3. Preliminaries and results from [2] and [5].

Given a homogeneous equation (1.3) where the $R_{j}(z)$ are functions which belong to a logarithmic differential field of rank zero over some
$F(a, b)$, we can follow the procedure developed in [5; pp.6-8] (or see [2; pp. 271-274]), and thereby determine for (1.3) the critical degree, the critical equation, the logarithmic set, the full factorization polynomial, the exponential set, and the transition set on $(a, b)$ as defined in [5]. As defined in [5; p. 8], if $M_{1}, M_{2}, \ldots M_{p}$ are the elements of the logarithmic set for (1.3), and if $\psi_{1}, \psi_{2}, \ldots \psi_{p}$ are admissible functions in some $F(c, d)$ such that for each $j$, the function $\psi_{j}$ solves (1.3) and satisfies $\psi_{j} \sim M_{j}$ in $F(c, d)$ for $j=1, \ldots, p$, then we will call $\left\{\psi_{1}, \ldots, \psi_{p}\right\}$ a complete logarithmic set of solutions of (1.3) in $F(c, d)$.

Using the above concepts, we recall the following result from [2; p. 274] or [5; p. 10-11]:

Theorem 3.2. Given the equation (1.3) where $n \geqslant 1$ and where the functions $R_{0}(z), \ldots, R_{n-1}(z)$ belong to a logarithmic differential field of rank zero over $F(a, b)$. When written in terms of the operator $\theta$ (where $\theta w=z w^{\prime}$ ) let (1.3) have the form $\Omega(w)=0$, where $\Omega(w)=$ $=\sum_{j=0}^{n} B_{j}(z) \theta^{j} w$. Let $p$ be the critical degree, let $F^{*}(\alpha)=0$ be the critical equation, and let $N_{1}, \ldots, N_{s}$ be the distinct elements (if any) of the exponential set for (1.3). Then there exist a nonnegative integer $d$, with $d \leqslant n-p$, and a set $\left\{V_{1}, \ldots, V_{d}\right\}$ of $d$ distinct functions such that all of the following hold:
(a) For each $j$, the function $V_{j}$ belongs to a logarithmic differential field of rank zero over $F(a, b)$, and there exists $k \in\{1, \ldots, s\}$ such that $V_{j} \sim N_{k}$ over $F(a, b)$.
(b) If $j \neq m$, then there exists a strictly positive real number $c=c(j, m)$ such that $V_{j}-V_{m} \approx z^{-1+c}$ over $F(a, b)$.
(c) For each $j \in\{1, \ldots, d\}$, the equation $\Omega_{j}(u)=0$, where,

$$
\begin{equation*}
\Omega_{j}(u)=\Omega\left(\left(\exp \int V_{j}\right) u\right) /\left(\exp \int V_{j}\right) \tag{3.8}
\end{equation*}
$$

has coefficients belonging to a logarithmic differential field of rank zero over $F(a, b)$, and has a strictly positive critical degree $t_{j}$;
(d) $t_{1}+\ldots+t_{d}=n-p$;
(e) if $k \in\{1, \ldots, s\}$ and if $N_{k}$ has multiplicity $m$ as a critical monomial of the full factorization polynomial for (1.3), then $\sum_{j \in J_{k}} t_{j}=m$,
where $J_{k}$ is the set of all $j \in\{1, \ldots, d\}$ for which $V_{j} \sim N_{k}$ over $F(a, b)$;
$(f)$ let $\Omega_{0}$ denote $\Omega$, and let $E_{1}$ denote the union (for $j$ in $\{0,1, \ldots, d\}$ ) of the transition sets of the equations $\Omega_{j}(u)=0$ on $(a, b)$, say $E_{1}=\left\{r_{1}, \ldots, r_{q}\right\}$ where $r_{1}<r_{2}<\ldots<r_{q}$. Let $r_{0}=a$ and $r_{q+1}=b$. (If $E_{1}$ is empty, set $q=0$.) Then in each of $F\left(r_{0}, r_{1}\right)$, $F\left(r_{1} r_{2}\right), \ldots, F\left(r_{q}, r_{q+1}\right)$ separately, the following hold:
(i) the equation $\Omega(w)=0$ possesses a complete logarithmic set of solutions $\left\{\phi_{1}, \ldots, \phi_{p}\right\}$, and for each $j \in\{1, \ldots, d\}$, the equation $\Omega_{j}(u)=0$ possesses a complete logarithmic set of solutions, $\left\{\phi_{j, 1}, \ldots, \phi_{j, t_{j}}\right\}$;
(ii) if we set $\Delta_{0}=\left\{\phi_{1}, \ldots, \phi_{p}\right\}$, and

$$
\begin{equation*}
\Delta_{j}=\left\{\left(\exp \int V_{j}\right) \phi_{j, 1}, \ldots,\left(\exp \int V_{j}\right) \phi_{j, t_{j}}\right\} \tag{3.9}
\end{equation*}
$$

for $j \in\{1, \ldots, d\}$ then the set $\Delta_{0} \cup \Delta_{1} \cup \ldots \cup \Delta_{d}$ is a fundamental set of solutions of (1.3).

REmark. The functions $V_{j}$ can be explicitly calculated from (1.3) (see[5; Theorem 3]), and hence the set $E_{1}$ can be explicitly calculated from (1.3).

The previous theorem can be used to discuss the oscillation properties of solutions of an equation (1.3). We now recall the following concepts which were introduced in [2; p. 276] for an equation (1.1) whose coefficients are admissible functions in some $F(a, b)$ :
(A) We say that (1.1) has the global oscillation property in $a<\operatorname{Arg} z<b$ if the following holds: For any $\theta \in(a, b)$ and any $\varepsilon>0$, there exist strictly positive constants $\delta$ and $K$, and a solution $f \not \equiv 0$ of (1.1), such that $\delta<\min \{\theta-a, b-\theta, \varepsilon\}$, and such that $f$ is analytic and has infinitely many zeros $z_{1}, z_{2}, \ldots$, with $\lim _{m \rightarrow \infty}\left|z_{m}\right|=+\infty$, on the region defined by $|\operatorname{Arg} z-\theta|<\delta$ and $|z|>K$.
(B) We say that (1.1) has the non-oscillation property in $a<\operatorname{Arg} z<b$, if for any solution $f \not \equiv 0$ of (1.1) which is admissible in $F(a, b)$, and any $\varepsilon>0$, there exists a constant $K=K(\varepsilon, f)>0$ such that $f$ is analytic and has no zeros on the set defined by,

$$
\begin{equation*}
a+\varepsilon \leqslant \operatorname{Arg} z \leqslant b-\varepsilon \quad \text { and } \quad|z|>K \tag{3.10}
\end{equation*}
$$

Using Theorem 3.2, we can now deduce the following oscillation result for equations (1.3) as stated in [2; p. 277]:

Theorem 3.3. Assume the hypothesis and notation of Theorem 3.2. Let $E_{2}$ denote the union (over all $j$ and $k$ in $\{1, \ldots, d\}$ with $j \neq k)$ of the sets of zeros on ( $a, b$ ) of the functions $I F\left(V_{j}-V_{k}, \theta\right)$. Let $E_{3}=E_{1} \cup E_{2}$ (where $E_{1}$ is defined in Part ( $f$ ) of Theorem 3.2), say $E_{3}=\left\{s_{1}, \ldots, s_{t}\right\}$, where $s_{1}<s_{2}<\ldots<s_{t}$. Set $s_{0}=a$ and $s_{t+1}=b$. (If $E_{3}$ is empty set $t=0$.) Then the following hold:
(A) If for some $j \in\{0,1, \ldots, d\}$, the critical equation of the equation $\Omega_{j}(u)=0$ possesses two distinct roots having the same real part, then (1.3) has the global oscillation property in $a<\operatorname{Arg} z<b$.
( $B$ ) If (1.3) does not satisfy the condition in Part ( $A$ ), then (1.3) possesses the non-oscillation property in $s_{k}<\operatorname{Arg} z<s_{k+1}$ for each $k \in\{0,1, \ldots, t\}$.

Remarks. (a) In view of the Remark after Theorem 3.2, it easily follows that the set $E_{3}=\left\{s_{1}, \ldots, s_{t}\right\}$ can be explicitly calculated from (1.3). It will be seen later that the numbers $s_{1}, \ldots, s_{t}$ in Theorem 1.1 for equation (1.1) are precisely the numbers $s_{1}, \ldots, s_{t}$ obtained in Theorem 3.3 for the corresponding homogeneous equation (1.3), and so they can be explicitly calculated from (1.1).
(b) It is easy to see that Theorem 1.2 follows immediately from Theorem 3.3 since we may take $(a, b)=(-\pi, \pi)$ when the coefficients of (1.3) are polynomials.

For non-homogeneous equations (1.1), the following result was proved in [2; p. 278]:

Theorem 3.4. Given the equation (1.1) where $n \geqslant 2$ and the $R_{j}(z)$ belong to a logarithmic differential field of rank zero over $F(a, b)$, and where $Q(z)$ is an admissible function in $F(a, b)$ for which there is a real number $\beta$ such that $Q(z) \approx z^{\beta}$ over $F(a, b)$. Using (3.1) and dividing the equation (1.1) through by a suitable real power of $z$, we can rewrite (1.1) in the form $\Omega(w)=\psi$, where $\Omega(w)$ is of the form (3.3), whose coefficients satisfy (3.4), and where for some complex $c \neq 0$ and some real number $\sigma$, we have $\psi \sim c z^{\sigma}$ over $F(a, b)$. We define a function $M(z)$ as follows: If $\sigma$ is not a root of the critical equation, $F^{*}(\alpha)=0$, of $\Omega(w)=0$, we set $M=\left(F^{*}(\sigma)\right)^{-1} c z^{\sigma}$. If $\sigma$ is a root, say of multiplicity $r$, of $F^{*}(\alpha)=$ $=0$, we set,

$$
\begin{equation*}
M(z)=K_{0}^{-1} c z^{\sigma}(\log z)^{r} \tag{3.11}
\end{equation*}
$$

where $K_{0}$ is the value at $\alpha=\sigma$ of the $r^{\text {th }}$ derivative of $F^{*}(\alpha)$. Let $r_{1}<r_{2}<\ldots<r_{q}$ be the points of the transition set for $\Omega(w)=0$ on $(a, b)$, and set $r_{0}=a$ and $r_{q+1}=b$ (If the transition set is empty, set $q=$ $=0)$. Then in each of $F\left(r_{0}, r_{1}\right), F\left(r_{1}, r_{2}\right), \ldots, F\left(r_{q}, r_{q+1}\right)$ separately, the equation (1.1) possesses a solution $w_{0} \sim M$.

For a non-homogeneous equation (1.1), the function $M(z)$ in (3.11) plays a key role in the oscillation properties of (1.1). Therefore, the following concepts were introduced in [2; p. 278]:

Under the hypothesis of Theorem 3.4, the number $\sigma$ will be called the principal exponent of (1.1). If $\sigma$ is a root of $F^{*}(\alpha)=0$, we will call $r$ the multiplicity of $\sigma$. If $\sigma$ is not a root of $F^{*}(\alpha)=0$, we will say that it has multiplicity $r=0$. Following the terminology of [10; §§67,69], the function $M(z)$ will be called the principal monomial of (1.1), and any solution $w(z)$ of (1.1) which satisfies $w \sim M$ over some $F(c, d)$, will be called a principal solution of (1.1) in $F(c, d)$.

Finally, as in [2; p. 280], the oscillation properties of the solutions of (1.1) can be deduced from Theorem 3.4 as follows:

TheOrem 3.5. Given an equation (1.1), where $n \geqslant 2$, where the coefficients $R_{j}(z)$ belong to a logarithmic differential field of rank zero over $F(a, b)$, and where $Q(z)$ is an admissible function in $F(a, b)$ for which there exists a real number $\beta$ such that $Q(z) \approx z^{\beta}$ over $F(a, b)$. Let $\sigma$ be the principal exponent of (1.1) with multiplicity $r$. Applying Theorem 3.2 to the homogeneous equation (1.3) corresponding to (1.1), let $p, F^{*}(\alpha), N_{1}, \ldots, N_{s}, V_{1}, \ldots, V_{d}$, and $\Omega_{0}, \ldots, \Omega_{d}$ be as defined in Theorem 3.2 and let $s_{0}, \ldots, s_{t+1}$ be as in Theorem 3.3. Then:
(A) Assume that at least one of the following two conditions is satisfied: (i) the equation $F^{*}(\alpha)=0$ possesses two distinct roots $\alpha$ and $\lambda$ such that $\operatorname{Re}(\alpha)=\operatorname{Re}(\lambda)$ and $\sigma<\operatorname{Re}(\alpha)$; (ii) the equation $F^{*}(\alpha)=0$ possesses a root $\alpha$ of multiplicity at least $r+1$ such that $\operatorname{Re}(\alpha)=\sigma$. Then, (1.1) possesses the global oscillation property in $a<\operatorname{Arg} z<b$.
(B) Let $k \in\{0,1, \ldots, t\}$, and assume that the following condition is satisfied: (iii) there exists $j \in\{1, \ldots, d\}$ such that $\operatorname{IF}\left(V_{j}, \theta\right)>0$ on $\left(s_{k}, s_{k+1}\right)$, and the critical equation of $\Omega_{j}(u)=0$ possesses at least two distinct roots having the same real part. Then equation (1.1) possesses the global oscillation property in $s_{k}<\operatorname{Arg} z<s_{k+1}$.
(C) Let $k \in\{0,1, \ldots, t\}$, and assume that neither condition (i) or (ii) in Part ( $A$ ) holds. Assume further that the following condition holds: (iv) for every $j \in\{1,2, \ldots, d\}$ such that the critical equation of $\Omega_{j}(u)=0$ has two distinct roots having the same real part, we have
$I F\left(V_{j}, \theta\right)<0$ on $\left(s_{k}, s_{k+1}\right)$. Then equation (1.1) possesses the non-oscillation property in $s_{k}<\operatorname{Arg} z<s_{k+1}$.

Remarks. (a). It is easy to see (see [2; p. 290]) that Theorem 1.1 follows from Theorem 3.5 when we take $(a, b)=(-\pi, \pi)$. We note that the points $s_{k}$ in Theorem 1.1 are the points $s_{k}$ given in Theorem 3.5.
(b) It is easy to give simple examples of equations (1.1) for which some of the points $s_{k}$ which are produced by Theorem 3.5 are extraneous in the sense described in $\S 1$. For example, we may take the simple equation,

$$
\begin{equation*}
w^{\prime \prime}+(1-i) w^{\prime}-i w=1 \tag{3.12}
\end{equation*}
$$

Of course, this equation can be explicitly solved and possesses the general solution,

$$
\begin{equation*}
w=i+c_{1} e^{-z}+c_{2} e^{i z} \quad \text { for arbitrary } c_{1}, c_{2} . \tag{3.13}
\end{equation*}
$$

Since $e^{-z}$ and $e^{i z}$ are trivial functions in $F(0, \pi / 2)$ (e.g., see Lemma 2.3), it follows that any solution of (3.12) satisfies $w \sim i$ on $F(0, \pi / 2)$. It follows from Lemma 2.1, that no solution of (3.12) can have infinitely many zeros on a sector $|\operatorname{Arg} z-(\pi / 4)|<\varepsilon$ where $\varepsilon<\pi / 4$. However, we now show that $\pi / 4$ is an $s_{k}$ produced by our method and so extraneous. To see this, we first use (3.1) to write (3.12) in the form $\Omega(w)=\psi$, where $\Omega(w)$ is of the form (3.3) whose coefficients satisfy (3.4). This form is,

$$
\begin{equation*}
z^{-2} \theta^{2} w+\left((1-i) z^{-1}-z^{-2}\right) \theta w-i w=1 . \tag{3.14}
\end{equation*}
$$

For the homogeneous equation corresponding to (3.14), the critical degree is clearly zero, and the full factorization polynomial $H(z, v)$ in (3.6) is,

$$
\begin{equation*}
H(z, v)=v^{2}+\left(1-i-z^{-1}\right) v-i . \tag{3.15}
\end{equation*}
$$

Using the remarks at the end of $\S 2$, we find that $H(z, v)$ has two critical monomials, $N_{1}=-1$ and $N_{2}=i$ and both have multiplicity 1. By definition, the set $\left\{N_{1}, N_{2}\right\}$ is the exponential set for the homogeneous equation corresponding to (3.14), and so in Theorem 3.2 (using Parts (a) and (e)), the set of functions $\left\{V_{1}, \ldots, V_{d}\right\}$ consists of two functions, $V_{1} \sim N_{1}$ and $V_{2} \sim N_{2}$ over $F(-\pi, \pi)$. Since $I F\left(V_{1}-V_{2}, \phi\right)$ vanishes at $\phi=\pi / 4$, it follows from Theorem 3.3 that $\pi / 4$ is an $s_{k}$ as we claimed. We note that by the same reasoning, $-3 \pi / 4$ is also an $s_{k}$, but our main result (which will be stated in the next section) will show that $-3 \pi / 4$ is
not extraneous because both of the functions $e^{-z}$ and $e^{i z}$ have large growth in $F(-\pi,-\pi / 2)$.

## 4. Main result.

We now state one of our main results. The proof will be concluded in §8.

Theorem 4.1. Given an equation (1.1), where $n \geqslant 2$, where the coefficients $R_{j}(z)$ belong to a logarithmic differential field of rank zero over $F(a, b)$, and where $Q(z)$ is an admissible function in $F(a, b)$ for which there exists a real number $\beta$ such that $Q(z) \approx z^{\beta}$ over $F(a, b)$. Applying Theorem 3.2 to the homogeneous equation (1.3) corresponding to (1.1), let $p, F^{*}(\alpha), N_{1}, \ldots, N_{s}, V_{1}, \ldots, V_{d}$, and $\Omega_{0}, \ldots, \Omega_{d}$ be as defined in Theorem 3.2, and let $s_{0}, \ldots, s_{t+1}$ be as in Theorem 3.3. Let $k \in\{1, \ldots, t\}$ and assume that equation (1.1) possesses the non-oscillation property in both $s_{k-1}<\operatorname{Arg} z<s_{k}$ and $s_{k}<\operatorname{Arg} z<s_{k+1}$ (see Theorem (3.5)). Then:
(A) Assume that at least one of the following two conditions is satisfied: (i) There exists $j \in\{1, \ldots, d\}$ such that $\operatorname{IF}\left(V_{j}, s_{k}\right)=0$; (ii) There exist distinct elements $j$ and $m$ in the set $\{1, \ldots, d\}$ such that $I F\left(V_{j}-V_{m}, s_{k}\right)=0$ while $I F\left(V_{j}, s_{k}\right)>0$ and $I F\left(V_{m}, s_{k}\right)>0$. Then, there is an admissible solution $f(z)$ of (1.1) in $F\left(s_{k-1}, s_{k+1}\right)$ such that for any $\varepsilon>0, f(z)$ possesses infinitely many zeros $z_{1}, z_{2}, \ldots$, with $\lim _{n \rightarrow \infty}\left|z_{n}\right|=+\infty$, which lie in the sector $\left|\operatorname{Arg} z-s_{k}\right|<\varepsilon$.
(B) Assume that neither condition (i) or (ii) in Part (A) holds. Then (1.1) has the non-oscillation property in $s_{k-1}<\operatorname{Arg} z<s_{k+1}$.

Remarks. 1) Theorem 4.1 gives a simple criterion for determining when an $s_{k}$, for $k \in\{1, \ldots, t\}$, is extraneous, namely when Part ( $B$ ) holds. In the case where the coefficients of (1.1) are rational functions (so that $(a, b)=(-\pi, \pi))$, one can ask if the negative real axis (i.e., $s_{t+1}=\pi$ ) is also extraneous in the sense described in $\S 1$. Clearly, Theorem 4.1 does not apply directly to $k=t+1$, but one can use it indirectly to test $s_{t+1}=\pi$ by the following simple device: We make the change of independent variable $\zeta=-z$ in (1.1) which converts the negative real axis into $\operatorname{Arg} \zeta=0$. We then determine whether the value 0 is an $s_{k}$ for the transformed equation (using Theorem 3.3), and, if so, we then use Theorems 3.5 and 4.1 to test it.
2) The method of proof of Theorem 4.1 will show that if either (i) or (ii) in Part (A) holds, then the sequence of zeros of the solution $f(z)$ in

Part (A) will have a strictly positive exponent of convergence (see[14; p. 50]).

## 5. Continuation theorems.

Definition 5.1. Given an equation (1.3) where we assume the hypothesis and notation of Theorem 3.2. Let $\left\{M_{0,1}, \ldots, M_{0, t_{0}}\right\}$ be the logarithmic set of $\Omega(w)=0$ (so that $t_{0}=p$ ) and for $j \in\{1, \ldots, d\}$, let $\left\{M_{j, 1}, \ldots, M_{j, t_{j}}\right\}$ denote the logarithmic set for the equation $\Omega_{j}(u)=0$. Setting $V_{0} \equiv 0$, and letting $V_{1}, \ldots, V_{d}$ be as in Theorem 3.2, the set of $n$ functions

$$
\begin{equation*}
M_{j, q} \exp \int V_{j} \quad \text { for } 0 \leqslant j \leqslant d, \quad 1 \leqslant q \leqslant t_{j}, \tag{5.1}
\end{equation*}
$$

will be called the asymptotic set for (1.3). (We note that the elements of the asymptotic set are admissible in $F(a, b)$.) If $\left\{G_{1}, \ldots, G_{n}\right\}$ denotes the asymptotic set for (1.3), and $\left\{f_{1}, \ldots, f_{n}\right\}$ is a fundamental set of solutions of (1.3) consisting of functions which are admissible in $F(c, d)$, where ( $c, d$ ) is a subset of ( $a, b$ ), and satisfying $f_{j} / G_{j} \rightarrow 1$ in $F(c, d)$ for $1 \leqslant j \leqslant n$, then we will call $\left\{f_{1}, \ldots, f_{n}\right\}$, a basic fundamental set for (1.3) in $F(c, d)$. (Thus, Theorem 3.2 asserts the existence of a basic fundamental set for (1.3) in each of $F\left(r_{0}, r_{1}\right), \ldots, F\left(r_{q}, r_{q+1}\right)$ separately, and hence in each $F\left(s_{k}, s_{k+1}\right)$ separately, for $0 \leqslant k \leqslant t$, where the $s_{k}$ are as in Theorem 3.3.)

Definition 5.2. Given an equation (1.1) where we assume the hypothesis and notation of Theorem 3.4. Let $M$ denote the principal monomial of (1.1). Any solution $w_{0}(z)$ of (1.1) which is admissible in $F(c, d)$, where $(c, d)$ is a subset of $(a, b)$, and which satisfies $w_{0} / M \rightarrow 1$ in $F(c, d)$ will be called a solution of principal type for (1.1) in $F(c, d)$. (Thus, as in Definition 5.1, the result in Theorem 3.4 asserts the existence of a solution of principal type for (1.1) in each $F\left(s_{k}, s_{k+1}\right)$ separately, for $0 \leqslant k \leqslant t$, where the $s_{k}$ are the points obtained by applying Theorem 3.3 to the homogeneous equation corresponding to (1.1).)

We now state our continuation theorems. The proofs will be given in §7.

Theorem 5.1. Given the equation (1.3) where $n \geqslant 1$ and where the functions $R_{0}(z), \ldots, R_{n-1}(z)$ belong to a logarithmic differential field of rank zero over $F(a, b)$. Let $s_{0}, \ldots, s_{t+1}$ be as in Theorem 3.3. Then for any $k \in\{1, \ldots, t\}$, the equation (1.3) possesses a basic fundamental set in $F\left(s_{k-1}, s_{k+1}\right)$.

THEOREM 5.2. Given the equation (1.1) where $n \geqslant 2$ and the $R_{j}(z)$ belong to a logarithmic differential field of rank zero over $F(a, b)$, and where $Q(z)$ is an admissible function in $F(a, b)$ for which there is a real number $\beta$ such that $Q(z) \approx z^{\beta}$ over $F(a, b)$. Let $s_{0}, \ldots, s_{t+1}$ be the points obtained by applying Theorem 3.3 to the homogeneous equation corresponding to (1.1). Then, for any $k \in\{1, \ldots, t\}$, the equation (1.1) possesses a solution of principal type in $F\left(s_{k-1}, s_{k+1}\right)$.

## 6. Preliminary results for continuation theorems.

We will require several preliminary results. The first is a combination of several Phragmen-Lindelöf principles whose proofs can be found in [14; pp. 176-180].

Lemma 6.1. Let $f(z)$ be analytic and of finite order of growth in a closed sectorial region of the form $\alpha \leqslant \arg z \leqslant \beta,|z| \geqslant K$. Then there exists $\delta>0$ such that for any real numbers $c$ and $d$, with $\alpha \leqslant c<d \leqslant \beta$ and $d-c<\delta$, for which the limits,

$$
\begin{equation*}
L_{1}=\lim _{r \rightarrow+\infty} f\left(r e^{i c}\right) \quad \text { and } \quad L_{2}=\lim _{r \rightarrow+\infty} f\left(r e^{i d}\right) \tag{6.1}
\end{equation*}
$$

exist and are finite, the following conclusions hold: $L_{1}=L_{2}$, and $f(z) \rightarrow L_{1}$ as $z \rightarrow \infty$ in $c \leqslant \arg z \leqslant d$.

Lemma 6.2. Given the equation (1.1) where the coefficients $R_{j}(z)$ belong to a logarithmic differential field of rank zero over $F(a, b)$, and where either $Q(z) \equiv 0$ or $Q(z)$ is an admissible function in $F(a, b)$ for which there is a real number $\beta$ such that $Q(z) \approx z^{\beta}$ over $F(a, b)$. Let $f(z)$ be a solution of (1.1) which is admissible in $F(a, b)$. Then, for any real numbers $c$ and $d$, with $a<c<d<b$, there exists $K>0$ such that $f$ is analytic and of finite order in the closed sectorial region defined by, $c \leqslant \operatorname{Arg} z \leqslant d,|z| \geqslant K$.

Proof. For homogeneous equations (i.e., $Q(z) \equiv 0$ ) the conclusion was proved in [5; Lemma 8.2] by using the method developed in [3] for transforming the equation to the unit disk and applying the ValironWiman theory (see[3; p. 149]). The same proof also works in the nonhomogeneous case because the non-homogeneous term $Q(z)$, when transformed to the unit disk, becomes a function whose growth is negligible with respect to the growth of a solution of positive order.

Lemma 6.3. Let $f(z)$ be analytic and of finite order of growth in a closed sectorial region of the form $\alpha \leqslant \arg z \leqslant \beta,|z| \geqslant K$. Then, for any
real numbers $\alpha_{1}$ and $\beta_{1}$ with $\alpha<\alpha_{1}<\beta_{1}<\beta$, there exists $K_{1}>0$ such that $f^{\prime}(z)$ is analytic and of finite order on $\alpha_{1} \leqslant \arg z \leqslant \beta_{1},|z| \geqslant$ $\geqslant K_{1}$.

Proof. This follows immediately from Cauchy's formula for derivatives.

Lemma 6.4. Let $f(z)$ be an admissible function in $F(a, b)$ and let $c \in(a, b)$. Assume that there exist real numbers $a_{1}$ and $b_{1}$ with $a<a_{1}<c<b_{1}<b$ such that $f$ is analytic and of finite order of growth on the closed sectorial region $a_{1} \leqslant \operatorname{Arg} z \leqslant b_{1},|z| \geqslant K$, for some $K>0$. Assume that for some complex number $\sigma$, we have $f \rightarrow \sigma$ over $F(c, b)$. Assume also that for some positive integer $k$, there exist complex constants $c_{1}, \ldots, c_{k}$, and distinct real numbers $\lambda_{1}, \ldots, \lambda_{k}$ such that,

$$
\begin{equation*}
f=\sum_{j=1}^{k} c_{j} z^{i \lambda_{j}}\left(1+E_{j}\right)+o(1) \quad \text { over } F(a, c) \tag{6.2}
\end{equation*}
$$

where $E_{1}, \ldots, E_{k}$ are admissible functions in $F(a, c)$ which are all $o(1)$ in $F(a, c)$. Then the following two conclusions hold:
(A) $f \rightarrow \sigma$ over $F(a, b)$.
(B) If $\sigma=0$ then all $c_{j}=0$ for $1 \leqslant j \leqslant k$.

Proof. We first prove Part ( $B$ ) by induction on $k$. We assume that $\sigma=0$. If $k=1$ in (6.2), we set $g=f / z^{i \lambda_{1}}$. Since $z^{i \lambda_{1}}$ is bounded from below on $|\operatorname{Arg} z|<\pi$ by a nonzero constant (see Lemma 2.5), it follows from the hypothesis that $g \rightarrow 0$ over $F(c, b)$ and $g \rightarrow c_{1}$ over $F(a, c)$. Since $g$ is clearly of finite order on $a_{1} \leqslant \operatorname{Arg} z \leqslant b_{1}$, for $|z|$ sufficiently large, it easily follows from Lemmas 2.1 and 6.1 that $c_{1}=0$ proving that Part ( $B$ ) holds for $k=1$.

We now assume that $k>1$ and that Part $(B)$ holds for $k-1$. If $f$ satisfies (6.2), we set $g=f / z^{i \lambda_{k}}$, so that over $F(a, c)$, we have

$$
\begin{equation*}
g=c_{k}\left(1+E_{k}\right)+\sum_{j=1}^{k-1} c_{j} z^{i\left(\lambda_{j}-\lambda_{k}\right)}\left(1+E_{j}\right)+E \tag{6.3}
\end{equation*}
$$

where $E \rightarrow 0$ by Lemma 2.5. We now compute $z g^{\prime}$ from (6.3), and observe that $z E^{\prime}$ and all $z E_{j}^{\prime}$ tend to zero over $F(a, c)$ by Part (A) of Lemma 2.2. Since the functions $z^{i\left(\lambda_{j}-\lambda_{k}\right)}$ are all bounded
on $|\operatorname{Arg} z|<\pi$ by Lemma 2.5, the resulting relation has the form,

$$
\begin{equation*}
z g^{\prime}=\sum_{j=1}^{k-1} c_{j} i\left(\lambda_{j}-\lambda_{k}\right) z^{i\left(\lambda_{j}-\lambda_{k}\right)}\left(1+E_{j}\right)+o(1) \quad \text { on } F(a, c) \tag{6.4}
\end{equation*}
$$

Again by Lemma 2.5, we have $g \rightarrow 0$ over $F(c, b)$ so that by Lemma 2.2, we have $z g^{\prime} \rightarrow 0$ over $F(c, b)$. Finally, we note that $z g^{\prime}$ is of finite order on some closed region $a_{2} \leqslant \operatorname{Arg} z \leqslant b_{2},|z| \geqslant K_{1}$ where $a_{1}<a_{2}<c<b_{2}<b_{1}$, by Lemma 6.3, and we note also that the numbers $\lambda_{j}-\lambda_{k}$ are all distinct for $1 \leqslant j \leqslant k-1$. Thus, by the induction hypothesis, it follows from (6.4) that $c_{j} i\left(\lambda_{j}-\lambda_{k}\right)=0$ for $1 \leqslant j \leqslant k-1$. Since the $\lambda_{j}$ are distinct, we have $c_{j}=0$ for $1 \leqslant j \leqslant k-1$, and so from (6.2) we have

$$
\begin{equation*}
f=c_{k} z^{i \lambda_{k}}\left(1+E_{k}\right)+o(1) \quad \text { over } F(a, c) . \tag{6.5}
\end{equation*}
$$

Since we have proved the result for $k=1$, we can now conclude that $c_{k}=0$ also, proving Part (B) by induction.

Proof of Part (A). We set $f_{1}=f-\sigma$ so that $f_{1} \rightarrow 0$ over $F(c, b)$. We distinguish two cases.

Case I. The set $\left\{\lambda_{1}, \ldots, \lambda_{k}\right\}$ in (6.2) does not contain zero. Thus, writing $f_{1}$ as $f-\sigma z^{i 0}$ and using (6.2), we can apply Part (B) to $f_{1}$ to conclude that $\sigma=0$ and $c_{j}=0$ for $1 \leqslant j \leqslant k$. Thus $f_{1} \rightarrow 0$ in $F(a, c)$ also. Using Lemmas 2.1 and 6.1, it follows that there is a $\delta>0$ such that $f_{1} \rightarrow 0$ as $z \rightarrow \infty$ in $|\operatorname{Arg} z-c| \leqslant \delta / 3$ and thus $f_{1} \rightarrow 0$ over $F(c-(\delta / 3)$, $c+(\delta / 3)$ ). From [10; Lemma 97], we can conclude that $f_{1} \rightarrow 0$ over $F(a, b)$ and Part $(A)$ is proved in this case.

Case II. The set $\left\{\lambda_{1}, \ldots, \lambda_{k}\right\}$ in (6.2) contains zero, say, $\lambda_{m}=0$. Again we set $f_{1}=f-\sigma$, and we have from (6.2) that in $F(a, c)$,

$$
\begin{equation*}
f_{1}=\sum_{j \neq m} c_{j} z^{i \lambda_{j}}\left(1+E_{j}\right)+\left(c_{m}-\sigma\right) z^{i \lambda_{m}}+o(1) . \tag{6.6}
\end{equation*}
$$

Since $f_{1} \rightarrow 0$ over $F(c, b)$, we can conclude from Part (A) that $c_{j}=0$ for $j \neq m$ and $c_{m}=\sigma$. Thus $f_{1} \rightarrow 0$ over $F(a, c)$ also, and as in Case I, we again obtain $f_{1} \rightarrow 0$ over $F(a, b)$ proving Part (A) completely.

## 7. Proofs of Theorems 5.1 and 5.2.

We begin with some notation:
Definition 7.1. Let $M_{1}=z^{\alpha}(\log z)^{m}$ and $M_{2}=z^{\sigma}(\log z)^{n}$, where
$\alpha$ and $\sigma$ are complex numbers, and $m$ and $n$ are nonnegative integers. We write $M_{1} \cong M_{2}$ if $\operatorname{Re}(\alpha)=\operatorname{Re}(\sigma)$ and $m=n$. (Thus, $M_{1} \cong M_{2}$ if and only if $M_{1} / M_{2}=z^{i \lambda}$ for some real number $\lambda$ ).

Definition 7.2. Given a homogeneous equation (1.3) whose coefficients $R_{j}(z)$ belong to a logarithmic differential field of rank zero over $F(a, b)$. Let $q$ be the critical degree of (1.3), and let $\left\{P_{1}, \ldots, P_{q}\right\}$ be the logarithmic set for (1.3), say,

$$
\begin{equation*}
P_{j}=z^{\alpha_{j}}(\log z)^{n_{j}} 1 \quad \text { for } j=1, \ldots, q \tag{7.1}
\end{equation*}
$$

where $\alpha_{j}$ is complex and $n_{j}$ is a nonnegative integer. Let $\alpha_{j}=\sigma_{j}+i \lambda_{j}$ where $\sigma_{j}$ and $\lambda_{j}$ are real, and consider all the triples $A_{j}=\left(\sigma_{j}, n_{j}, \lambda_{j}\right)$. We arrange these triples in decreasing lexicographic order, say, $A_{j_{1}}, \ldots, A_{j_{q}}$ (so that for each $k$, we have $\sigma_{j_{k}} \geqslant \sigma_{j_{k+1}}$, and if equality holds, we have $n_{j_{k}} \geqslant n_{j_{k+1}}$, and if equality holds again, we have $\lambda_{j_{k}} \geqslant \lambda_{j_{k+1}}$. Since the triples are all distinct, this order is a total order). The corresponding $q$-tuple ( $P_{j_{1}}, \ldots, P_{j_{q}}$ ) will be called the ordered logarithmic system for (1.3). By partitioning this system by keeping together only those triples which have the same pairs $\left(\sigma_{j}, n_{j}\right)$ and if we set $Q_{m}=P_{j_{m}}$ for $m=1, \ldots, q$, then we can assume that the ordered logarithmic system has the form $\left(Q_{1}, \ldots, Q_{q}\right)$ and can be partitioned as follows:

$$
\begin{equation*}
\left(Q_{1}, Q_{2}, \ldots, Q_{d_{1}}\right), \quad\left(Q_{d_{1}+1},, \ldots, Q_{d_{2}}\right), \ldots,\left(Q_{d_{r}+1},, \ldots, Q_{q}\right) \tag{7.2}
\end{equation*}
$$

where by Definition 7.1 and Lemmas 2.4 and 2.5(a), we have on $F(-\pi, \pi)$,

$$
\begin{equation*}
Q_{1} \approx \ldots \cong Q_{d_{1}} \quad \text { and } \quad Q_{m}=o\left(Q_{1}\right) \quad \text { for } m>d_{1} \tag{7.3}
\end{equation*}
$$

and in general (setting $d_{r+1}=q$ ),

$$
\begin{equation*}
Q_{d_{j}+1} \cong \ldots \cong Q_{d_{j+1}} \quad \text { and } \quad Q_{m}=o\left(Q_{d_{j}+1}\right) \quad \text { for } m>d_{j+1} \tag{7.4}
\end{equation*}
$$

over $F(-\pi, \pi)$ for $1 \leqslant j \leqslant r$. We will call (7.2) the canonical partition of $\left(Q_{1}, \ldots, Q_{q}\right)$.

Proof of Theorem 5.1. We are given an equation (1.3), and we assume the notation in Theorems 3.2 and 3.3 when we apply these results to (1.3). Let $k \in\{1, \ldots, t\}$, and set $H_{m}=\exp \int V_{m}$ for $1 \leqslant m \leqslant d$, and set $H_{0} \equiv 1$. In each of $\left(s_{k-1}, s_{k}\right)$ and $\left(s_{k}, s_{k+1}\right)$, none of the indicial functions $\operatorname{IF}\left(V_{j}, \theta\right)$ or $\operatorname{IF}\left(V_{j}-V_{m}, \theta\right)$ for $j \neq m$, have any zeros and so must be of constant sign. It follows from Lemma 2.3 that if $j$ and $m$ are
distinct elements of $\{0,1, \ldots, d\}$ then
(7.5) either $H_{j} / H_{m}$ or $H_{m} / H_{j}$ is trivial in $F\left(s_{k-1}, s_{k}\right)$,
and
(7.6) either $H_{j} / H_{m}$ or $H_{m} / H_{j}$ is trivial in $F\left(s_{k}, s_{k+1}\right)$.

It easily follows (e.g., see[2; Lemma 7.2]) that there exist permutations $\left\{m_{0}, \ldots, m_{d}\right\}$ and $\left\{q_{0}, \ldots, q_{d}\right\}$ of $\{0,1, \ldots, d\}$ such that both of the following hold:

$$
\begin{align*}
& H_{m_{0}} \gg H_{m_{1}} \gg \ldots \gg H_{m_{d}} \quad \text { in } F\left(s_{k-1}, s_{k}\right) ;  \tag{7.7}\\
& H_{q_{0}} \gg H_{q_{1}} \gg \ldots \gg H_{q_{d}} \quad \text { in } F\left(s_{k}, s_{k+1}\right) . \tag{7.8}
\end{align*}
$$

Let $\Delta_{0}, \ldots, \Delta_{d}$ denote the sets of solutions in (3.9) of equation (1.3) in $F\left(s_{k}, s_{k+1}\right)$, and let $\bar{\Delta}_{0}, \ldots, \bar{\Delta}_{d}$ denote the corresponding sets in $F\left(s_{k-1}, s_{k}\right)$. We begin with $\Delta_{m_{0}}$, and we let ( $Q_{1},, \ldots, Q_{q}$ ) denote the ordered logarithmic system for the equation $\Omega_{m_{0}}(u)=0$, with canonical partition (7.2). By Theorem 3.2, there are admissible functions $\psi_{1}, \ldots, \psi_{q}$ in $F\left(s_{k}, s_{k+1}\right)$ and there are admissible functions $g_{1}, \ldots, g_{q}$ in $F\left(s_{k-1}, s_{k}\right)$ such that

$$
\begin{equation*}
\Delta_{m_{0}}=\left\{\psi_{1} H_{m_{0}}, \ldots, \psi_{q} H_{m_{0}}\right\} \quad \text { and } \quad \psi_{j} \sim Q_{j} \quad \text { in } F\left(s_{k}, s_{k+1}\right) \tag{7.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\bar{\Delta}_{m_{0}}=\left\{g_{1} H_{m_{0}}, \ldots, g_{q} H_{m_{0}}\right\} \quad \text { and } \quad g_{j} \sim Q_{j} \quad \text { in } F\left(s_{k-1}, s_{k}\right) \tag{7.10}
\end{equation*}
$$

We begin with solution $\psi_{1} H_{m_{0}}$ of (1.3). By basic existence theory (e.g.[15; Theorem 2.2]), this solution has an extension $G_{1}(z)$ which is admissible in $F(a, b)$. Hence in $F\left(s_{k-1}, s_{k}\right)$, the solution $G_{1}$ can be written as a linear combination of the elements of $\bar{\Delta}_{0} \cup \ldots \cup \bar{\Delta}_{d}$. Thus in view of (7.5), (7.7) and (7.10), there exist constants $c_{1}, \ldots, c_{q}$ such that,

$$
\begin{equation*}
G_{1} / H_{m_{0}}=\sum_{j=1}^{q} c_{j} g_{j}+T_{1} \quad \text { on } F\left(s_{k-1}, s_{k}\right) \tag{7.11}
\end{equation*}
$$

where $T_{1}$ is trivial in $F\left(s_{k-1}, s_{k}\right)$. Thus from (7.10) and (7.3), and Lemma 2.4 , we have,

$$
\begin{equation*}
G_{1} / Q_{1} H_{m_{0}}=\sum_{j=1}^{d_{1}} c_{j}\left(g_{j} / Q_{1}\right)+o(1) \quad \text { on } F\left(s_{k-1}, s_{k}\right) \tag{7.12}
\end{equation*}
$$

But from (7.2) and (7.10) clearly for $1 \leqslant j \leqslant d_{1}$, there is a real number $\lambda_{j}$ such that,

$$
\begin{equation*}
g_{j} / Q_{1}=z^{i \lambda_{j}}(1+o(1)) \quad \text { over } F\left(s_{k-1}, s_{k}\right), \tag{7.13}
\end{equation*}
$$

and the numbers $\lambda_{1}, \ldots, \lambda_{d_{1}}$ are all distinct since the functions $Q_{1}, \ldots, Q_{d_{1}}$ are distinct. In addition, we note that on $F\left(s_{k}, s_{k+1}\right)$, we have $G_{1} / Q_{1} H_{m_{0}} \rightarrow 1$ since $G_{1}=\psi_{1} H_{m_{0}}$ and $\psi_{1} \sim Q_{1}$. Finally, we note that $G_{1} / Q_{1} H_{m_{0}}$ is of finite order of growth on any sector $a_{1} \leqslant \operatorname{Arg} z<b_{1}$ where $a<a_{1}<b_{1}<b$, for $|z|$ sufficiently large, since $G_{1}$ has this property by Lemma 6.2. Thus, we can apply Lemma 6.4 to conclude that,

$$
\begin{equation*}
G_{1} / Q_{1} H_{m_{0}} \rightarrow 1 \quad \text { over } F\left(s_{k-1}, s_{k+1}\right) \tag{7.14}
\end{equation*}
$$

Since $Q_{1} H_{m_{0}}$ is in the asymptotic set for (1.3), we now have constructed the first element, namely $G_{1}$, in the basic fundamental set for (1.3) in $F\left(s_{k-1}, s_{k+1}\right)$. The same proof is obviously valid for the extension $G_{j}$ of $\psi_{j} H_{m_{0}}$ for each $j=1, \ldots, d_{1}$, so now have solutions $G_{1}, \ldots, G_{d_{1}}$ of (1.3) such that for $j=1, \ldots, d_{1}$

$$
\begin{equation*}
G_{j} / Q_{j} H_{m_{0}} \rightarrow 1 \quad \text { over } F\left(s_{k-1}, s_{k+1}\right) \tag{7.15}
\end{equation*}
$$

We now consider the solution $\psi_{d_{1}+1} H_{m_{0}}$, where $d_{1}$ is as in (7.2). As above, this solution has an extension $G_{d_{1}+1}$ to $F(a, b)$, which is of finite order of growth, and as in (7.11) we have,

$$
\begin{equation*}
G_{d_{1}+1} / H_{m_{0}}=\sum_{j=1}^{q} c_{j} g_{j}+T_{d_{1}+1} \quad \text { on } F\left(s_{k-1}, s_{k}\right), \tag{7.16}
\end{equation*}
$$

where $T_{d_{1}+1}$ is trivial. Dividing (7.16) by $Q_{1}$, we obtain (7.12) with $G_{1}$ replaced by $G_{d_{1}+1}$. Of course, (7.13) is still valid. But now, $G_{d_{1}+1} / Q_{1} H_{m_{0}}$ tends to zero over $F\left(s_{k}, s_{k+1}\right)$ by (7.9) and (7.3), since it agrees with $\psi_{d_{1}+1} / Q_{1}$. Thus by Lemma 6.4, we must have $c_{j}=0$ for $1 \leqslant$ $\leqslant j \leqslant d_{1}$. Thus, when we divide (7.16) by $Q_{d_{1}+1}$, we are now back in the situation we had in (7.12), and we can argue as before to show $G_{d_{1}+1} / Q_{d_{1}+1} H_{m_{0}} \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$. Similarly, all the extensions $G_{j}$ of $\psi_{j} H_{m_{0}}$ for $d_{1}+1 \leqslant j \leqslant d_{2}$ yield solutions satisfying $G_{j} / Q_{j} H_{m_{0}} \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$. If $r>1$ in (7.2), we consider the extension $G_{d_{2}+1}$ of $\psi_{d_{2}+1} H_{m_{0}}$. As above $G_{d_{2}+1} / H_{m_{0}}$ can be written as a linear combination (7.16). As above, by dividing by $Q_{1}$ and using Lemma 6.4, we obtain $c_{j}=0$ for $1 \leqslant j \leqslant d_{1}$. Then dividing the relation by $Q_{d_{1}+1}$, we obtain using Lemma 6.4 that $c_{j}=0$ for $j \leqslant d_{2}$. Then we are back in the situ-
ation in (7.12) and we show that $G_{d_{2}+1} / Q_{d_{2}+1} H_{m_{0}} \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$. In this way, we obtain solutions $G_{1}, \ldots, G_{q}$, which are admissible and satisfy

$$
\begin{equation*}
G_{j} / Q_{j} H_{m_{0}} \rightarrow 1 \quad \text { in } F\left(s_{k-1}, s_{k+1}\right) \tag{7.17}
\end{equation*}
$$

We also note that $\left\{G_{1}, \ldots, G_{q}\right\}$ together with the union of the sets $\bar{\Delta}_{m_{j}}$ for $j=1, \ldots, d$, form a fundamental set for (1.3) in $F\left(s_{k-1}, s_{k}\right)$. This can be seen as follows: A dependence relation would obviously take the form,

$$
\begin{equation*}
\sum_{j=1}^{q} c_{j}\left(G_{j} / H_{m_{0}}\right)+T_{2}(z) \equiv 0 \tag{7.18}
\end{equation*}
$$

where $T_{2}$ is trivial in $F\left(s_{k-1}, s_{k}\right)$. Dividing (7.18) by $Q_{1}$, and noting that $G_{j} / Q_{1} H_{m_{0}}$ is of the form $z^{i \lambda_{j}}(1+o(1))$ for $1 \leqslant j \leqslant d_{1}$, we can apply Lemma 6.4 to the zero function (where we take $(a, b)=\left(s_{k-1}, s_{k}\right)$ and $c$ to be any point in $\left(s_{k-1}, s_{k}\right)$ ) to yield $c_{j}=0$ for $1 \leqslant j \leqslant d_{1}$. We then divide (7.18) by $Q_{d_{1}+1}$ and repeat the argument. In this way we obtain $c_{j}=0$ for $1 \leqslant j \leqslant q$ and independence is now clear since the union of the sets $\bar{\Delta}_{m_{j}}$ consists of linearly independent solutions.

We now proceeded by induction. We assume that $r$ is a nonnegative integer less than $d$, and we assume that we have constructed solutions $G_{j, l}$ of (1.3) for $j=m_{0}, \ldots, m_{r}$ and $1 \leqslant l \leqslant t_{j}$, which are admissible in $F\left(s_{k-1}, s_{k+1}\right)$ and satisfy

$$
\begin{equation*}
G_{j, l} / M_{j, l} H_{j} \rightarrow 1 \quad \text { over } F\left(s_{k-1}, s_{k+1}\right), \tag{7.19}
\end{equation*}
$$

(where the $M_{j . l}$ are as in (5.1)) and have the property that together with the sets $\bar{\Delta}_{m_{j}}$ for $j=r+1, \ldots, d$, form a fundamental set for (1.3) in $F\left(s_{k-1}, s_{k}\right)$.

We use induction on $r$, and for ease of notation, let $\left(Q_{1}, \ldots, Q_{q}\right)$ denote the ordered logarithmic system for $\Omega_{m_{r+1}}(u)=0$, with canonical partition (7.2). For this equation, let $\psi_{1}, \ldots, \psi_{q}$ and $g_{1}, \ldots, g_{q}$ be the complete logarithmic sets of solutions in $F\left(s_{k}, s_{k+1}\right)$ and $F\left(s_{k-1}, s_{k}\right)$ respectively, with $\psi_{j} \sim Q_{j}$ in $F\left(s_{k}, s_{k+1}\right)$ and $g_{j} \sim Q_{j}$ in $F\left(s_{k-1}, s_{k}\right)$. We let $\Gamma$ denote the subset of $\left\{m_{0}, \ldots, m_{r}\right\}$ consisting of those $m_{j}$ for which $H_{m_{r+1}} \ll H_{m_{J}}$ in $F\left(s_{k}, s_{k+1}\right)$. For ease of notation, we denote the elements of $\Gamma$ by $m_{\alpha}, m_{\beta}, \ldots$ where $\alpha<\beta<\ldots$ Let $\Phi$ denote $\left\{m_{0}, \ldots, m_{r}\right\}-\Gamma$, so that by the total ordering (7.8), we have $H_{m_{j}} \ll H_{m_{r+1}}$ in $F\left(s_{k}, s_{k+1}\right)$ for $m_{j} \in \Phi$. As before, the solution $\psi_{1} H_{m_{r+1}}$ has an extension $G_{1}$ to $F(a, b)$, and so can be written as a linear combi-
nation of the $G_{j, l}$ for $j=m_{0}, \ldots, m_{r}$, and the elements of $\bar{\Delta}_{m_{j}}$ for $j \geqslant r+$ +1 , in $F\left(s_{k-1}, s_{k}\right)$. Setting $U_{j, l}=G_{j, l} / H_{j}$ (so that $U_{j, l} / M_{j, l} \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$ by (7.19)) this combination for $G_{1}$ in $F\left(s_{k-1}, s_{k}\right)$ takes the form,

$$
\begin{equation*}
G_{1}=W_{m_{0}} H_{m_{0}}+\ldots+W_{m_{r}} H_{m_{r}}+W H_{m_{r+1}}+\sum_{j=r+2}^{d} W_{m_{j}} H_{m_{j}} \tag{7.20}
\end{equation*}
$$

where for $j \leqslant r, W_{m_{j}}$ is a linear combination of the functions $U_{m_{j, l}}$, while $W$ is a linear combination of $g_{1}, \ldots, g_{q}$, and where for $j>r+1, W_{m_{j}}$ is a linear combination of a complete logarithmic set of solutions of $\Omega_{m_{j}}(u)=0$. We now set,

$$
\begin{equation*}
G_{1}^{*}=G_{1}-\sum_{m_{j} \in \Phi} W_{m_{j}} H_{m_{j}} \quad \text { on } F\left(s_{k-1}, s_{k+1}\right) \tag{7.21}
\end{equation*}
$$

so we have in $F\left(s_{k-1}, s_{k}\right)$,

$$
\begin{equation*}
G_{1}^{*}=W_{m_{\alpha}} H_{m_{\alpha}}+W_{m_{\beta}} H_{m_{\beta}}+\ldots+W H_{m_{r+1}}+\sum_{j=r+2}^{d} W_{m_{j}} H_{m_{j}} \tag{7.22}
\end{equation*}
$$

On $F\left(s_{k-1}, s_{k}\right.$ ), we have from (7.22), (7.5) and (7.7) that $G_{1}^{*} / H_{m_{\alpha}}=$ $=W_{m_{\alpha}}^{*}+T$ where $T=o\left(z^{c}\right)$ for all $c<0$. However, on $F\left(s_{k}, s_{k+1}\right)$, it follows from (7.21) and the definitions of $\Gamma, \Phi$ and $G_{1}$ that $G_{1}^{*} / H_{m_{a}}=o\left(z^{c}\right)$ for every $c<0$. This is the same situation that we had in (7.16) and we argue the same way to prove all coefficients in $W_{m_{\alpha}}$ are zero. That is, we canonically partition the ordered logarithmic set for $\Omega_{m_{\alpha}}(u)=0$ as in (7.2), say ( $D_{1}, D_{2}, \ldots$ ), $\left(D_{e}, \ldots\right), \ldots$, and we use Lemma 6.4 successively on $G_{1}^{*} / D_{1} H_{m_{\alpha}}, G_{1}^{*} / D_{e} H_{m_{\alpha}}$, etc. to conclude that all coefficients in $W_{m_{\alpha}}$ are zero. We then do the same for $W_{m_{\beta}}$ and so on, so we obtain from (7.22) that

$$
\begin{equation*}
G_{1}^{*}=W H_{m_{r+1}}+\sum_{j=r+2}^{d} W_{m_{j}} H_{m_{j}} \quad \text { on } F\left(s_{k-1}, s_{k}\right) \tag{7.23}
\end{equation*}
$$

Again from (7.21), we have $G_{1}^{*} / Q_{1} H_{m_{r+1}} \rightarrow 1$ in $F\left(s_{k}, s_{k+1}\right)$ while from (7.23), we have on $F\left(s_{k-1}, s_{k}\right)$,

$$
\begin{equation*}
G_{1}^{*} / Q_{1} H_{m_{r+1}}=\sum_{j=1}^{q} c_{j}\left(g_{j} / Q_{1}\right)+T_{2} \tag{7.24}
\end{equation*}
$$

where $T_{2}$ is trivial. This is the same situation as we had in (7.12) and by using Lemma 6.4 again, we obtain $G_{1}^{*} / Q_{1} H_{m_{r+1}} \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$. By replacing $Q_{1}$ in this argument by $Q_{j}$ for any $j=$
$=1, \ldots, d_{1}$, we obtain a solution $G_{j}^{*}$ of (1.3) such that $G_{j}^{*} / Q_{j} H_{m_{r+1}} \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$.

We then consider the extension $G_{d_{1}+1}$ of $\psi_{d_{1}+1} H_{m_{r+1}}$ to $F(a, b)$, so that $G_{d_{1}+1}$ is given by a linear combination as in (7.20). We then subtract $\sum_{m_{j} \in \Phi} W_{m_{j}} H_{m_{j}}$ from $G_{d_{1}+1}$, and call the resulting function $G_{d_{1}+1}^{*}$ which is given by the right side of (7.22) in $F\left(s_{k-1}, s_{k}\right)$. As above, we show $W_{m_{\alpha}}=0, W_{m_{\beta}}=0$, etc. so that $G_{d_{1}+1}^{*}$ is given by the right side of (7.23). This is the same situation as in (7.16), and by dividing the relation (7.23) by $Q_{1} H_{m_{r+1}}$ and using Lemma 6.4, we find that all coefficients of $g_{1}, \ldots, g_{d_{1}}$ in $W$ are zero. Then dividing (7.23) by $Q_{d_{1}+1}$, we find as in (7.16) that by Lemma 6.4 we obtain $G_{d_{1}+1}^{*} / Q_{d_{1}+1} H_{m_{r+1}} \rightarrow 1$ in $F\left(s_{k-1}, s_{k+1}\right)$. We continue this process, and we have now constructed solutions $G_{j, l}$ of (1.3) for $j=m_{r+1}$ and $1 \leqslant l \leqslant t_{j}$ such that

$$
\begin{equation*}
G_{j, l} / M_{j, l} H_{j} \rightarrow 1 \quad \text { over } F\left(s_{k-1}, s_{k+1}\right) . \tag{7.25}
\end{equation*}
$$

To complete the induction, we must show that these $G_{j, l}$, for $j=m_{0}, \ldots, m_{r+1}$, together with $\bar{\Delta}_{m_{j}}$ for $j \geqslant r+2$ form a fundamental set for (1.3) in $F\left(s_{k-1}, s_{k}\right)$. This is easily proved exactly as was done for the case $r=0$ in (7.18). If we have a dependence relation, we divide it by $H_{m_{0}}$ yielding a relation (7.18), which was shown to imply $c_{j}=0$ for $1 \leqslant j \leqslant q$. We then divide the dependence relation by $H_{m_{1}}$ and the same argument shows that all coefficients of the functions $G_{m_{1}, l}$ will be zero. We then continue and we obtain the desired conclusion that all coefficients in the dependence relation are zero.

## 8. Proof of Theorem 5.2.

Let $\sigma$ be the principal exponent of (1.1) with multiplicity $r$, so that the principal monomial $M(z)$ of (1.1) is given by (3.11). By Theorem 3.4, the equation (1.1) possesses a principal solution $w_{1}$ in $F\left(s_{k-1}, s_{k}\right)$ and a principal solution $w_{2}$ in $F\left(s_{k}, s_{k+1}\right)$, so,

$$
\begin{equation*}
w_{1} \sim M \text { in } F\left(s_{k-1}, s_{k}\right) \text { and } w_{2} \sim M \text { in } F\left(s_{k}, s_{k+1}\right) . \tag{8.1}
\end{equation*}
$$

For the homogeneous equation (1.3) corresponding to (1.1) we adopt the notation developed in the proof of Theorem 5.1. In particular, we assume (7.5)-(7.8) hold, and we let $G_{j, l}$ ) be the elements in the basic fundamental set for (1.3) in $F\left(s_{k-1}, s_{k+1}\right)$ satisfying the relation (7.19). We also note that from Lemma 2.3, in each of $F\left(s_{k-1}, s_{k}\right)$ and
$F\left(s_{k}, s_{k+1}\right)$ separately, we have,
(8.2) for each $m=1, \ldots, d$, either $H_{m}$ or $1 / H_{m}$ is trivial.

By basic existence theory, the solution $w_{2}$ has an extension $D_{2}$ which is admissible (and of finite order of growth by Lemma 6.2) in $F(a, b)$. Hence on $F\left(s_{k-1}, s_{k}\right)$ we can write (as in (7.20)),

$$
\begin{equation*}
D_{2}=w_{1}+W_{m_{0}} H_{m_{0}}+W_{m_{1}} H_{m_{1}}+\ldots+W_{m_{d}} H_{m_{d}} \quad \text { on } F\left(s_{k-1}, s_{k}\right) \tag{8.3}
\end{equation*}
$$

where $W_{m_{j}}$ is a linear combination of the functions $U_{m_{j}, l}=G_{m_{j}, l} / H_{m_{j}}$ for $l=1, \ldots, t_{m_{j}}$. We now distinguish two cases.

Case $I$. The critical degree of (1.3) is zero. In this case the term in (8.3) for which $m_{j}=0$, vanishes identically. From (8.2), we distinguish two subcases. First, if all $H_{m_{j}}$ (for $\left.m_{j} \neq 0\right)$, are trivial in $F\left(s_{k-1}, s_{k}\right)$, then $D_{2} / M \rightarrow 1$ in both $F\left(s_{k}, s_{k+1}\right)$ and in $F\left(s_{k-1}, s_{k}\right)$ (by (8.3)). It now follows from Lemmas 6.1 and 2.1, and [10; Lemma 97], that $D_{2} / M \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$, and so $D_{2}$ is the desired solution. Second, if for some $m_{j} \neq 0$, the function $H_{m_{j}}$ is not trivial in $F\left(s_{k-1}, s_{k}\right)$, then referring to (7.7), let $q$ be the largest element of $\{0,1, \ldots, d\}$ for which $m_{q} \neq 0$ and $H_{m_{q}}$ is not trivial in $F\left(s_{k-1}, s_{k}\right)$. Then, from (8.3), we can write,

$$
\begin{equation*}
D_{2}=w_{1}+W_{m_{0}} H_{m_{0}}+\ldots+W_{m_{q}} H_{m_{q}}+T_{1} \quad \text { on } F\left(s_{k-1}, s_{k}\right) \tag{8.4}
\end{equation*}
$$

where $T_{1}=o\left(z^{c}\right)$ for all $c<0$. We now let $\Phi$ denote the subset of $\left\{m_{0}, \ldots, m_{q}\right\}$ consisting of those $m_{j}$ for which $H_{m_{j}}$ is trivial in $F\left(s_{k}, s_{k+1}\right)$, and we let $\Gamma$ denote $\left\{m_{0}, \ldots, m_{q}\right\}-\Phi$. Let $D_{2}^{*}$ denote the function $D_{2}-\sum_{m_{j} \in \Phi} W_{m_{j}} H_{m_{j}}$ on $F\left(s_{k-1}, s_{k+1}\right)$, so we have,

$$
\begin{equation*}
D_{2}^{*}=w_{1}+W_{m_{\alpha}} H_{m_{\alpha}}+W_{m_{\beta}} H_{m_{\beta}}+\ldots+T_{1} \quad \text { on } F\left(s_{k-1}, s_{k}\right), \tag{8.5}
\end{equation*}
$$

where we are denoting the elements of $\Gamma$ by $m_{\alpha}, m_{\beta}, \ldots$, where $\alpha<\beta<\ldots$. Of course, on $F\left(s_{k}, s_{k+1}\right), D_{2}^{*} / M \rightarrow 1$, so that $D_{2}^{*} / H_{m_{\alpha}}=$ $=o\left(z^{c}\right)$ for all $c<0$. This is exactly the same situation we had in (7.21) and (7.22), and, as shown there, we find $W_{m_{\alpha}}=0, W_{m_{\beta}}=0, \ldots$, and we obtain from (8.5) that $D_{2}^{*} / M \rightarrow 1$ over $F\left(s_{k-1}, s_{k}\right)$ also. As in the first subcase, $D_{2}^{*}$ is then the desired solution.

Case II. The critical degree of (1.3) is not zero. In this case, the term in (8.3) with $m_{j}=0$ need not vanish. For definiteness, let $l$ denote the element of $\{0,1, \ldots, d\}$ for which $m_{l}=0$ (and so $H_{m_{l}} \equiv 1$ ). Again we distinguish two subcases. First, assume that for all $m_{j} \neq 0$, the function $H_{m_{j}}$ is trivial in $F\left(s_{k-1}, s_{k}\right)$. Of course, by (7.7), then $l=0$, and so
by (8.3), we have

$$
\begin{equation*}
D_{2}=w_{1}+W_{m_{0}}+T_{2} \quad \text { on } F\left(s_{k-1}, s_{k}\right), \tag{8.6}
\end{equation*}
$$

where $T_{2}=o\left(z^{c}\right)$ for all $c<0$. Let $\left(P_{1}, \ldots, P_{q}\right)$ be the ordered logarithmic system for (1.3), where the $P_{j}$ are given by (7.1). Then, of course, for some constants $c_{j}$, we have

$$
\begin{equation*}
W_{m_{0}}=\sum_{j=1}^{q} c_{j} P_{j}\left(1+E_{j}\right) \quad \text { where } \quad E_{j} \rightarrow 0 \text { over } F\left(s_{k-1}, s_{k+1}\right) \tag{8.7}
\end{equation*}
$$

We can write $W_{m_{0}}=u_{1}+u_{2}+u_{3}$, where $u_{1}$ consists of the terms in (8.7) where $\operatorname{Re}\left(\alpha_{j}\right)<\sigma$, where $u_{2}$ consists of the terms where $\operatorname{Re}\left(\alpha_{j}\right)=$ $=\sigma$, and where $u_{3}$ consists of the terms where $\operatorname{Re}\left(\alpha_{j}\right)>\sigma$. By Lemma 2.4, we have from (8.6) that,

$$
\begin{equation*}
D_{2}=w_{1}+u_{2}+u_{3}+o\left(w_{1}\right) \quad \text { in } F\left(s_{k-1}, s_{k}\right) \tag{8.8}
\end{equation*}
$$

Let $\left(P_{1}, \ldots, P_{d_{1}}\right),\left(P_{d_{1}+1}, \ldots, P_{d_{2}}\right), \ldots$, be those portions of the canonical partition of $\left(P_{1}, \ldots, P_{q}\right)$ consisting of those $P_{j}$ for which $\operatorname{Re}\left(\alpha_{j}\right)>\sigma$. Now $D_{2} / P_{1} \rightarrow 0$ in $F\left(s_{k}, s_{k+1}\right)$ while from (8.8) we have in $F\left(s_{k-1}, s_{k}\right)$,

$$
\begin{equation*}
D_{2} / P_{1}=\sum_{j=1}^{d_{1}} c_{j}\left(P_{j} / P_{1}\right)\left(1+E_{j}\right)+o(1) . \tag{8.9}
\end{equation*}
$$

Since $P_{j} / P_{1}=z^{i \lambda_{j}}$ for some real $\lambda_{j}$ if $1 \leqslant j \leqslant d_{1}$, we can conclude from Lemma 6.4 that $c_{j}=0$ for $1 \leqslant j \leqslant d_{1}$. We can then repeat the argument using $P_{d_{1}+1}$ etc. instead of $P_{1}$, and we conclude that $u_{3} \equiv 0$. Thus,

$$
\begin{equation*}
D_{2}=w_{1}+u_{2}+o\left(w_{1}\right) \quad \text { in } F\left(s_{k-1}, s_{k}\right) \tag{8.10}
\end{equation*}
$$

We note that if $u_{2} \equiv 0$, then $D_{2}$ is the desired solutions as in Case $I$. Thus we may assume that there are $P_{j}$ for which $\operatorname{Re}\left(\alpha_{j}\right)=\sigma$, and we let $\left(P_{\alpha}, \ldots, P_{e_{1}}\right),\left(P_{e_{1}+1}, \ldots, P_{e_{2}}\right), \ldots$, be those portions of the canonical partition having $\operatorname{Re}\left(\alpha_{j}\right)=\sigma$. Noting that within each of these portions, the exponent $n_{j}$ of $\log z$ in (7.1) is the same, let these exponents in these portions be $\beta_{1}>\beta_{2}>\ldots$ respectively. We note first that if $\beta_{1}<r$ (where $r$ is as in (3.11)), then by Lemma 2.5, clearly $u_{2}=o\left(w_{1}\right)$ in $F\left(s_{k-1}, s_{k}\right)$ and again by (8.10), the desired solution is $D_{2}$ as before. Thus we may assume that $\beta_{1} \geqslant r$. Assume first that $\beta_{1}>r$, and let $u_{4}$ be the sum of the terms in $u_{2}$ corresponding to $\beta_{1}$ (that is, corresponding to ( $\left.P_{\alpha}, \ldots, P_{e_{1}}\right)$ ). From (8.10), we obtain,

$$
\begin{equation*}
D_{1} / P_{\alpha}=\sum_{j \alpha}^{e_{1}} c_{j}\left(P_{j} / P_{\alpha}\right)\left(1+E_{j}\right)+o(1) \quad \text { in } F\left(s_{k-1}, s_{k}\right) \tag{8.11}
\end{equation*}
$$

But since $\beta_{1}>r$, clearly $D_{2} / P_{\alpha} \rightarrow 0$ in $F\left(s_{k}, s_{k+1}\right)$, so from Lemma 6.4, we conclude $c_{j}=0$ for $j=\alpha, \ldots, e_{1}$. Thus $u_{4} \equiv 0$. If $\beta_{2}>r$, we apply the same argument, and we show that all terms in $u_{2}$ corresponding to those $\beta_{j}$ for which $\beta_{j}>r$, are all zero. Thus we may write $u_{2}=u_{5}+u_{6}$, where $u_{5}$ consists of those terms (if any) where $\beta_{j}=r$, and $u_{6}$ contains the terms where $\beta_{j}<r$. Clearly $u_{6}=o\left(w_{1}\right)$ in $F\left(s_{k-1}, s_{k}\right)$ so we have,

$$
\begin{equation*}
D_{2}=w_{1}+u_{5}+o\left(w_{1}\right) \quad \text { in } F\left(s_{k-1}, s_{k}\right) . \tag{8.12}
\end{equation*}
$$

If $u_{5} \equiv 0$, we are done as in Case I. Hence we may assume that in $F\left(s_{k-1}, s_{k}\right)$,

$$
\begin{equation*}
u_{5}=\sum_{j \in J} c_{j} z^{\sigma+i \lambda_{j}}(\log z)^{r}\left(1+E_{j}\right) \quad \text { where } E_{j} \rightarrow 0 \tag{8.13}
\end{equation*}
$$

for some nonempty set $J$, and some distinct real numbers $\lambda_{j}$. We note that $\lambda_{j} \neq 0$ if $c_{j} \neq 0$ since $\lambda_{j}=0$ would imply that $z^{\sigma}(\log z)^{r}$ belongs to the logarithmic set for (1.3). But then $\sigma$ would be a root of the critical equation of (1.3) having multiplicity at least $r+1$ which contradicts the definition of $r$. Thus we have in $F\left(s_{k-1}, s_{k}\right)$ from (8.12) and (8.13),

$$
\begin{equation*}
D_{2} / M=\left(w_{1} / M\right)+\sum_{j \in J} c_{j} z^{i \lambda_{j}}\left(1+E_{j}\right)+o(1), \tag{8.14}
\end{equation*}
$$

where we note that $w_{1} / M=z^{i 0}(1+E)$, where $E \rightarrow 0$, since $w_{1} \sim M$ in $F\left(s_{k-1}, s_{k}\right)$. Since $D_{2} / M \rightarrow 1$ over $F\left(s_{k}, s_{k+1}\right)$ (since $D_{2} \equiv w_{2}$ ), we can conclude from Lemma 6.4 that $D_{2} / M \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$. Thus $D_{2}$ is the desired solution.

The second subcase in Case II, is the case where not all $H_{m_{j}}$ (for $m_{j} \neq 0$ ) are trivial in $F\left(s_{k-1}, s_{k}\right)$. Thus clearly from (7.7), $l>0$. Let $\Phi$ denote the subset of $\left\{m_{0}, \ldots, m_{l-1}\right\}$ consisting of those $m_{j}$ for which $H_{m_{j}}$ is trivial in $F\left(s_{k}, s_{k+1}\right)$, and let $\Gamma$ denote $\left\{m_{0}, \ldots, m_{l-1}\right\}-\Phi$. Letting $D_{2}^{*}$ denote the function $D_{2}-$ $-\sum_{m_{j} \in \Phi} W_{m_{j}} H_{m_{j}}$, we have from (8.3),

$$
\begin{equation*}
D_{2}^{*}=w_{1}+W_{m_{l}}+W_{m_{\alpha}} H_{m_{\alpha}}+W_{m_{\beta}} H_{m_{\beta}}+\ldots+T_{3} \quad \text { on } F\left(s_{k-1}, s_{k}\right), \tag{8.15}
\end{equation*}
$$

where $T_{3}=o\left(z^{c}\right)$ for all $c<0$, and where we are denoting the elements of $\Gamma$ by $m_{\alpha}, m_{\beta}, \ldots$ where $\alpha<\beta<\ldots$. This is identical to the situation in (8.5), and we showed that $W_{m_{\alpha}} \equiv 0, W_{m_{\beta}} \equiv 0, \ldots$. Thus, we have

$$
\begin{equation*}
D_{2}^{*}=w_{1}+W_{m_{l}}+T_{3} \quad \text { on } F\left(s_{k-1}, s_{k}\right) . \tag{8.16}
\end{equation*}
$$

Since $D_{2}^{*} / M \rightarrow 1$ over $F\left(s_{k}, s_{k+1}\right)$, this is exactly the same situation (with $D_{2}$ replaced by $D_{2}^{*}$ ) as we had in (8.6), and the same argument
shows that $D_{2}^{*} / M \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$. Thus $D_{2}^{*}$ is the desired solution, and the proof is complete.

## 9. Proof of Theorem 4.1.

We assume the hypothesis of the theorem, and we adopt the notation developed in the previous two proofs. In particular, we let $G_{j}, l$ be the elements of the basic fundamental set for (1.3) in $F\left(s_{k-1}, s_{k+1}\right)$, satisfying (7.19) (and where $H_{m}=\exp \int V_{m}$ and $H_{0} \equiv 1$ ). Let $W$ denote the solution of principal type of (1.1) in $F\left(s_{k-1}, s_{k+1}\right)$, so that $W / M \rightarrow 1$ where $M$ is the principal monomial (3.11).

Proof of Part ( $A$ ). Assume condition (i) in Part ( $A$ ) holds. For this $j$, the function $f=W-G_{j, 1}$ is a solution of (1.1), and in view of (7.19), the equation $f(z)=0$ can be written in the form,

$$
\begin{equation*}
z^{\alpha}(\log z)^{\beta}\left(1+E_{1}\right) \exp \int V_{j}=1 \tag{9.1}
\end{equation*}
$$

where $E_{1} \rightarrow 0$ over $F\left(s_{k-1}, s_{k+1}\right)$, and where $z^{\alpha}(\log z)^{\beta}$ is $M_{j, 1} / M$. This is the same equation as was obtained in [5; Formula (10.5)], and as shown in [5], since $\operatorname{IF}\left(V_{j}, s_{k}\right)=0$, the solution $f(z)$ satisfies the conclusion of Part (A).

Now assume (ii) holds, and consider the solution $f(z)$ of (1.1) defined by $W-G_{j, 1}-G_{m, 1}$. The equation $f(z)=0$ can be written in the form,

$$
\begin{equation*}
\left(U_{j, 1} /\left(\left(W \exp \int-V_{m}\right)-U_{m, 1}\right)\right) \exp \int\left(V_{j}-V_{m}\right)=1 \tag{9.2}
\end{equation*}
$$

where $U_{j, 1}=G_{j, 1} / H_{j}$ and $U_{m, 1}=G_{m, 1} / H_{m}$. Since $U_{j, 1} / M_{j, 1} \rightarrow 1$ and $U_{m, 1} / M_{m, 1} \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$ by (7.19), and since $W \exp \int-V_{m}=$ $=o\left(U_{m, 1}\right)$ by Lemmas 2.3 and 2.4, it is clear that (9.2) has the form (9.1) with $V_{j}$ replaced by $V_{j}-V_{m}$, and where $z^{\alpha}(\log z)^{\beta}$ is $-M_{j, 1} / M_{m, 1}$. Thus $f$ has the desired property as before. This proves Part (A).

Proof of Part ( $B$ ). We now assume that neither (i) nor (ii) in Part (A) holds. We will show that for any admissible solution $f$ of (1.1), there is an element of $F\left(s_{k-1}, s_{k+1}\right)$ on which $f$ has no zeros. This will yield the conclusion of Part (B) by Lemma 2.1. Letting $U_{j, l}$ denote $G_{j, l} / H_{j}$,
so that $U_{j, l} / M_{j, l} \rightarrow 1$ over $F\left(s_{k-1}, s_{k+1}\right)$, the solution $f$ can be written,

$$
\begin{equation*}
f=W+W_{0}+W_{1} \exp \int V_{1}+\ldots+W_{d} \exp \int V_{d} \tag{9.3}
\end{equation*}
$$

where $W_{j}$ is a linear combination of the functions $U_{j, l}$. We now distinguish two cases.

Case I. For each $j \in\{1, \ldots, d\}$ for which $W_{j} \not \equiv 0$, we have $\operatorname{IF}\left(V_{j}, s_{k}\right)<0$. Since $\operatorname{IF}\left(V_{j}, \theta\right)$ has no zeros on $\left(s_{k-1}, s_{k}\right)$ or $\left(s_{k}, s_{k+1}\right)$, we must have $\operatorname{IF}\left(V_{j}, \theta\right)<0$ on $\left(s_{k-1}, s_{k+1}\right)$. Thus by Lemmas 2.3 and 2.4, we can write (9.3) as,

$$
\begin{equation*}
f=W+W_{0}+T_{1} \tag{9.4}
\end{equation*}
$$

where $T_{1}=o\left(z^{-c}\right)$ for any $c>0$ over $F\left(s_{k-1}, s_{k+1}\right)$. If $W_{0} \equiv 0$, then (9.4) shows that $f$ has no zeros on some element of $F\left(s_{k-1}, s_{k+1}\right)$, since $f / M \rightarrow 1$, and so we are done. If $W_{0} \neq 0$, we observe that by the hypothesis of the theorem, clearly neither of the conditions (i) nor (ii) of Part (A) of Theorem 3.5 can hold (or we would have global oscillation). This is the same hypothesis as we had in [2; Lemma 7.5], and the same arguments shows that $W+W_{0}$ is either of the form $W(1+o(1))$ or $K_{0} U_{0, l}(1+o(1))$ in $F\left(s_{k-1}, s_{k+1}\right)$ for some $l \in\left\{1, \ldots, t_{0}\right\}$ and some nonzero constant $K_{0}$. In either case, (9.4) shows that $f$ has no zeros on some element of $F\left(s_{k-1}, s_{k+1}\right)$. Thus we are done in Case I.

Case II. There exists $j \in\{1, \ldots, d\}$ for which $W_{j} \not \equiv 0$ and $I F\left(V_{j}, s_{k}\right) \geqslant 0$. Let $J$ denote the set of all $j \in\{1, \ldots, d\}$ with this property. In view of the assumption that condition (i) does not hold, we must have $\operatorname{IF}\left(V_{j}, s_{k}\right)>0$ for $j \in J$, and hence we have $\operatorname{IF}\left(V_{j}, \theta\right)>0$ on ( $s_{k-1}, s_{k+1}$ ) since $\operatorname{IF}\left(V_{j}, \theta\right)$ has constant sign on both $\left(s_{k-1}, s_{k}\right)$ and $\left(s_{k}, s_{k+1}\right)$.

Subcase $A$. $J$ has only one element, $m$. Then the relation (9.3) takes the form,

$$
\begin{equation*}
f=W+W_{0}+W_{m} \exp \int V_{m}+T_{2} \tag{9.5}
\end{equation*}
$$

where $T_{2}=o\left(z^{-c}\right)$ for all $c>0$ over $F\left(s_{k-1}, s_{k+1}\right)$. In view of the hypothesis that (1.1) does not have the global oscillation property in $s_{k}<\operatorname{Arg} z<s_{k+1}$, we have by Part (B) of Theorem 3.5, that the critical equation of $\Omega_{m}(u)=0$ does not possess two distinct roots having the same real part. In view of Lemma 2.4, one term in the linear combina-
tion $W_{m}$ dominates the rest, so that for some $l \in\left\{1, \ldots, t_{m}\right\}$ we have for some nonzero constant $K_{O}$, the relation

$$
\begin{equation*}
f=K_{O} M_{m, l}\left(\exp \int V_{m}\right)(1+o(1)) \quad \text { in } F\left(s_{k-1}, s_{k+1}\right) \tag{9.6}
\end{equation*}
$$

which shows that $f$ has no zeros on some element of $F\left(s_{k-1}, s_{k+1}\right)$,
Subcase B. $J$ has more than one element. Since we are assuming (ii) does not hold, if $j$ and $m$ are distinct elements of $J$, then $I F\left(V_{j}-V_{m}, s_{k}\right) \neq 0$. Since $\operatorname{IF}\left(V_{j}-V_{m}, \theta\right)$ has constant sign on each of $\left(s_{k-1}, s_{k}\right)$ and $\left(s_{k}, s_{k+1}\right)$ by definition of the $s_{k}$, it follows that $I F\left(V_{j}-V_{m}, \theta\right)$ has constant sign on $\left(s_{k-1}, s_{k+1}\right)$ for any two distinct elements $j$ and $m$ in $J$. Thus by Lemma 2.3, the set $Y=\left\{\exp \int V_{j}: j \in J\right\}$ has the property that the ratio of two distinct elements of $Y$ is either trivial or its reciprocal is trivial in $F\left(s_{k-1}, s_{k+1}\right)$. It easily follows that for some element $m \in J$, we have

$$
\begin{equation*}
\left(\exp \int V_{j}\right) /\left(\exp \int V_{m}\right) \quad \text { is trivial in } F\left(s_{k-1}, s_{k+1}\right) \tag{9.7}
\end{equation*}
$$

for any $j \in J-\{m\}$. It easily follows from (9.3) that

$$
\begin{equation*}
f=W_{m}\left(\exp \int V_{m}\right)(1+o(1)) \quad \text { over } F\left(s_{k-1}, s_{k+1}\right) \tag{9.8}
\end{equation*}
$$

However, by the same reasoning as in Subcase $A, W_{m}=K_{O} M_{m, l^{\circ}}$ $\cdot(1+o(1))$ for some $l$ and now it is clear that $f$ has no zeros on some element of $F\left(s_{k-1}, s_{k+1}\right)$. This concludes the proof.

## 10. The homogeneous case [5].

For a homogeneous equation (1.3) whose coefficients belong to a logarithmic differential field of rank zero over $F(a, b)$, we known from Theorem 3.3 that either (1.3) has the global oscillation property in $a<\operatorname{Arg} z<b$, or it possesses the non-oscillation property in all the sectors $s_{k}<\operatorname{Arg} z<s_{k+1}$ for $k=0,1, \ldots, t$. In the latter case, we can ask whether a given $s_{k}$ (for $k \in\{1, \ldots, t\}$ ) is extraneous in the sense that there is an $\varepsilon>0$ such that for any solution $f \not \equiv 0$ of (1.3), $f$ has no zeros on a set of the form, $\left|\operatorname{Arg} z-s_{k}\right|<\varepsilon,|z| \geqslant K(f)$ for some $K(f)>0$. The question of determining whether an $s_{k}$ is extraneous for (1.3) was completely solved in [5; §7] as follows: Under the given hypothesis, Theorem 3.3 shows that the condition in Part ( $A$ ) of Theorem 3.3 does
not hold. It easily follows from Lemmas 2.3 and 2.4, that in each of $F\left(s_{k-1}, s_{k}\right)$ and $F\left(s_{k}, s_{k+1}\right)$ separately, the whole asymptotic set $\left\{G_{1}, \ldots, G_{n}\right\}$ for (1.3) (see $\S 5$ ) is totally ordered by the relation «<<». That is, there exist permutations $\left\{m_{1}, \ldots, m_{n}\right\}$ and $\left\{q_{1}, \ldots, q_{n}\right\}$ of $\{1, \ldots, n\}$ such that

$$
\begin{equation*}
G_{q_{1}}>G_{q_{2}} \gg \ldots \gg G_{q_{n}} \quad \text { in } F\left(s_{k}, s_{k+1}\right), \tag{10.1}
\end{equation*}
$$

and

$$
\begin{equation*}
G_{m_{1}} \gg G_{m_{2}} \gg \ldots \gg G_{m_{n}} \quad \text { in } F\left(s_{k-1}, s_{k}\right) . \tag{10.2}
\end{equation*}
$$

It is proved in $[5 ; \S 7]$ that $s_{k}$ is extraneous if and only if $G_{q_{j}}=G_{m_{j}}$ for all $j=1, \ldots, n$. Hence an $s_{k}$ is not extraneous if and only if the ordering of the asymptotic set changes as we pass the ray $\operatorname{Arg} z=s_{k}$.

## REFERENCES

[1] S. Bank, A note on the location of complex zeros of solutions of linear differential equations, Complex Variables, 12 (1989), pp. 159-167. Research announcement, Bull. Amer. Math. Soc. (New Series), 18 (1988), pp. 35-38.
[2] S. Bank, On the oscillation of solutions of non-homogeneous linear differential equations, Analysis, 10 (1990), pp. 265-293.
[3] S. BANK, A note on the rate of growth of solutions of algebraic differential equations in sectors, J. London Math. Soc., (2), 1 (1969), pp. 145-154.
[4] S. Bank, On the instability theory of differential polynomials, Ann. Math. Pura Appl., 74 (1966), pp. 84-112.
[5] S. BANK, On oscillation, continuation, and asymptotic expansions of solutions of linear differential equations, Rend, Sem. Mat. Univ. Padova, 85 (1991), pp. 1-25.
[6] E. W. Chamberlain, The univalence of functions asymptotic to nonconstant logarithmic monomials, Proc. Amer. Math. Soc., 17 (1966), pp. 302-309.
[7] E. Hille, Ordinary Differential Equations in the Complex Domain, Wi-ley-Interscience, New York (1976).
[8] E. Hille, Analytic Function Theory, Volume II, Chelsea, New York (1973).
[9] R. Nevanlinna, Über Riemannsche Flächen mit endlich vielen Windungspunkten, Acta Math., 58 (1932), pp. 295-373.
[10] W. Strodt, Contributions to the asymptotic theory of ordinary differential equations in the complex domain, Mem. Amer. Math. Soc., No., 13 (1954).
[11] W. Strodt, Principal solutions of ordinary differential equations in the complex domain, Mem. Amer. Math. Soc., No., 26 (1957).
[12] W. Strodt, On the algebraic closure of certain partially ordered fields, Trans. Amer. Math. Soc., 105 (1962), pp. 229-250.
[13] W. Strodt - R. W. Wright, Asymptotic behavior of solutions and adjunction fields for nonlinear first-order differential equations, Mem. Amer. Math. Soc., No., 109 (1971).
[14] E. C. Titchmarsh, Theory of Functions, Oxford University Press, London (1939).
[15] W. Wasow, Asymptotic Expansions for Ordinary Differential Equations, Wiley, New York (1965).
[16] H. Witтich, Eindeutige Lösungen der Differentialgeichung $w^{\prime}=R(z, w)$, Math. Z., 74 (1960), pp. 278-288.

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