# ADJOINT QUASI-DIFFERENTIAL OPERATORS OF EULER TYPE 

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This paper treats linear quasi-differential operators of the form

$$
L[y]=\sum_{j=0}^{n} p_{0 j} y^{(j)}-\left(\sum_{j=0}^{n} p_{1 \nu} y^{(j)}-\left(\cdots-\left(\sum_{j=0}^{n} p_{m\lrcorner} y^{(j)}\right)^{\prime} \cdots\right)^{\prime}\right)^{\prime},
$$

based on an integrable $(m+1) \times(n+1)$ matrix function [ $p_{i j}$ ], ( $i=0, \cdots, m ; j=0, \cdots, n$ ), about which suitable regularity assumptions are made. Results obtained by Reid (Trans. Amer. Math. Soc. Vol. 85 (1957), pp. 446-461) are extended to operators of the type considered here.

A generalized Green's function for the system $\{L[y]=0$, $y \in \mathscr{D}\}$ is defined, where $\mathscr{D}$ is a linear subspace of the domain of $L$. Resolvent and deterministic properties of this function are presented, together with the relationship of such a generalized Green's function to the generalized Green's function for the associated adjoint system.

For a large class of two-point boundary problems in which the boundary conditions involve the characteristic parameter linearly it is shown that there exists a simultaneous canonical representation of the boundary conditions for a given problem and those of its adjoint; in particular, in the self-adjoint case this canonical representation has the form of boundary conditions and transversality conditions for a variational problem. Finally, these results are applied to a two-point boundary problem involving a differential operator of the type considered in the paper of Reid above.

Since an important example of an operator of the form of $L[y]$ is the Euler operator in the calculus of variations, we shall refer to such operators as quasi-differential operators of Euler type.

Section 2 gives a more precise description of the operator, and Section 3 is concerned with a discussion of its adjoint. In particular it is shown that if $\mathscr{D}_{0}$ is the class of functions $y$ in the domain of $L$ with the property that the functions $y, y^{\prime}, \cdots, y^{(n-1)}, \widetilde{y}_{m} \equiv \sum_{j=0}^{n} p_{m j} y^{(j)}$, $\widetilde{y}_{i} \equiv \sum_{j=0}^{n} p_{i j} y^{(3)}-\widetilde{y}_{i+1}^{\prime},(i=m-1, \cdots, 1)$, vanish at $\alpha$ and at $b$, and if $T_{0}$ is the restriction of $L$ to $\mathscr{D}_{0}$, then the adjoint operator $T_{0}^{*}$ is given by

[^0]$$
T_{0}^{*}[z]=L^{\grave{\chi}}[z] \equiv \sum_{i=0}^{m} \bar{p}_{i 0} z^{(i)}-\left(\sum_{i=0}^{m} \bar{p}_{i 1} z^{(\nu)}-\left(\cdots-\left(\sum_{i=0}^{m} \bar{p}_{i n} z^{(i)}\right)^{\prime} \cdots\right)^{\prime}\right)^{\prime} .
$$

Section 4 is a study of extensions of the operator $T_{0}$, and their adjoints. Section 5 is devoted to generalized Green's functions for Euler type quasi-differential systems and their adjoints, and extends the results of Elliott [3] and Reid [5] to the case where the number of linearly independent boundary conditions may differ from the order of the differential equation.

Section 6 is concerned with a certain class of two-point boundary problems in which the boundary conditions involve the characteristic parameter linearly. It is shown that there exists a simultaneous canonical representation of the boundary conditions for a given problm and those of its adjoint; in particular, in the self-adjoint case this canonical representation has the form of boundary conditions and transversality conditions for a variational problem.

Finally, $\S 7$ is devoted to an application of the results of $\S 6$ to a two-point boundary problem involving a differential operator of the type considered by Reid in [7].

The symbol $\mathfrak{c}_{n},(n=0,1,2, \cdots)$, will signify the class of complexvalued functions defined on the compact interval $[a, b]$ which have $n$ continuous derivatives. The set of functions $y$ in $\mathfrak{C}_{n-1}$ for which $y^{(n-1)}$ is a.c. (absolutely continuous) is denoted by $\mathfrak{H}_{n},(n=0,1,2, \cdots)$. In particular, $\mathfrak{\Im}_{0}$ and $\mathfrak{Y}_{0}$ will signify respectively the classes of continuous and Lebesgue integrable complex-valued functions defined on $[a, b]$. If $f$ and $g$ belong to $\mathfrak{N}_{0}$ and $f(x)=g(x)$ almost everywhere, we will simply write $f=g$. If $f$ is a complex-valued function on $[a, b]$, then $\bar{f}$ denotes the function with domain $[a, b]$ whose value at $x$ is the complex conjugate of $f(x)$. If $u$ and $v$ are functions on $[a, b]$ and $\bar{v} u \in \mathfrak{U}_{0}$, then we define ( $u, v$ ) as

$$
(u, v)=\int_{a}^{b} \bar{v} u
$$

Matrix notation will be used except where it is impracticable. If $M$ is a matrix, then the conjugate transpose of $M$ is denoted by $M^{*}$. Vectors are treated as matrices with one column. The symbols $E_{n}$ and $0_{m n}$ are used to represent the $n \times n$ identity matrix and the $m \times n$ zero matrix, respectively; the subscripts will be omitted when there is no danger of confusion.

A matrix function is said to be continuous, integrable, etc. whenever each of its elements possesses the specified property. If $A$ is an a.c. matrix function, then $A^{\prime}(x)$ signifies the matrix of derivatives at values for which these derivatives exist and the zero matrix elsewhere.
2. Description of the operator. Suppose that $\left[p_{i j}\right],(i=0, \cdots$,
$m \geqq 1 ; j=0, \cdots, n \geqq 1$ ), is an integrable $(m+1) \times(n+1)$ matrix function on a compact interval $[a, b]$ and that $p_{o n}$ and $p_{m o}$ are essentially bounded. For suitable $y$ in $\mathscr{Y}_{n}$ define functions $\widetilde{y}_{1}, \cdots, \widetilde{y}_{m}$ as follows:

$$
\begin{align*}
& \widetilde{y}_{m}(x)=\sum_{j=0}^{n} p_{m j}(x) y^{(j)}(x) \\
& \text { if } \widetilde{y}_{j+1} \in \mathfrak{U}_{1}, \text { then } \widetilde{y}_{i}(x)=\sum_{j=0}^{n} p_{i j}(x) y^{(j)}(x)-\widetilde{y}_{i+1}^{\prime}(x),  \tag{2.1}\\
& (i=m-1, \cdots, 1)
\end{align*}
$$

The class of functions $y$ in $\mathfrak{U}_{n}$ for which $\widetilde{y}_{1}, \cdots, \widetilde{y}_{m}$ are a.c. will be denoted by $\tilde{\mathfrak{A}}_{n}$. For convenience the vector functions ( $y^{(j-1)}$ ), $(j=1, \cdots, n)$, and $\left(\widetilde{y}_{i}\right),(i=1, \cdots, m)$, will be denoted by $\hat{y}$ and $\widetilde{y}$, respectively; the $(n+m)$-vector function ( $y, \cdots, y^{(n-1)}, \widetilde{y}_{1}, \cdots, \widetilde{y}_{m}$ ) will be represented by $\hat{y}$.

Denote by $L$ the operator with domain $\tilde{\mathfrak{A}}_{n}$ which is defined by

$$
\begin{equation*}
L[y]=\sum_{j=0}^{n} p_{o j} y^{(j)}-\widetilde{y}_{1}^{\prime} . \tag{2.2}
\end{equation*}
$$

The operator $L$ is a quasi-differential operator in the sense of Bôcher [1]; in particular, it is a generalization of the Euler operator in the calculus of variations and, as was stated in the introduction, it will be called a quasi-differential operator of the Euler type.

Let $\tilde{\mathfrak{A}}_{n}^{0}$ be the collection of functions $y$ in $\tilde{\mathfrak{A}}_{n}$ for which $\bar{y}(a)=0=$ $\widehat{y}(b)$, and denote by $T_{0}$ the restriction of $L$ to $\tilde{\mathfrak{Y}}_{n}^{0}$. Suppose that $\mathscr{D}_{0}^{*}$ is the class of functions $z$ in $\mathfrak{N}_{0}$ which are essentially bounded and have the property that there exists a function $f_{z}$ in $\mathfrak{N}_{0}$ such that $(L[y], z)=\left(y, f_{z}\right)$ for all $y$ in $\widetilde{\mathfrak{N}}_{n}^{0}$.

A second operator $L^{\star}$ will now be defined. For suitable functions $z$ in $\mathfrak{U}_{m}$ define functions $\widetilde{z}_{1}, \cdots, \widetilde{z}_{n}$ as follows:

$$
\begin{align*}
& \widetilde{z}_{n}(x)=\sum_{i=0}^{m} \bar{p}_{i n}(x) z^{(i)}(x) ; \\
& \text { if } \widetilde{z}_{j+1} \in \mathfrak{\varkappa}_{1}, \text { then } \widetilde{z}_{j}(x)=\sum_{i=0}^{m} \bar{p}_{i j}(x) \boldsymbol{z}^{(i)}(x)-\widetilde{z}_{j+1}^{\prime}(x),  \tag{2.3}\\
& (j=n-1, \cdots, 1) .
\end{align*}
$$

The class of functions $z$ in $\mathfrak{N}_{m}$ for which $\widetilde{z}_{1}, \cdots, \widetilde{z}_{n}$ are a.c. will be denoted by $\widetilde{\mathfrak{A}}_{m}$. Let $L^{\star}$ be the operator with domain $\widetilde{\mathfrak{A}}_{m}$ defined by

$$
\begin{equation*}
L^{\text {凶 }}[z]=\sum_{i=0}^{m} \bar{p}_{i o} z^{(i)}-\widetilde{z}_{1}^{\prime} . \tag{2.4}
\end{equation*}
$$

If $z \in \widetilde{\mathfrak{A}}_{m}$, then $\tilde{z}$ and $\widetilde{z}$ will signify the vector functions $\left(z^{(i-1)}\right)$, $(i=1, \cdots, m)$, and $\left(\widetilde{z}_{j}\right),(j=1, \cdots, n)$, respectively. The $(m+n)$ vector function $\left(z, \cdots, z^{(m-1)}, \widetilde{z}_{1}, \cdots, \widetilde{z}_{n}\right)$ will be denoted by $\bar{z}$.

Except when a statement is made to the contrary, the following hypothesis will be assumed throughout this paper.

HYpothesis (H). The matrix $\left[p_{i j}(x)\right],(i=0, \cdots, m ; j=0, \cdots, n)$, is integrable and there exists an $\varepsilon>0$ such that $\left|p_{m n}(x)\right| \geqq \varepsilon$ almost everywhere on $[a, b]$. Moreover, $p_{0 n}$ and $p_{m 0}$ are essentially bounded and $p_{i n} p_{m n}^{-1} p_{m j}$ is integrable, $(i=1, \cdots, m-1 ; j=1, \cdots, n-1)$.

It is to be noted that if $y \in \widetilde{\mathfrak{A}}_{n}$ and $z \in \widetilde{\mathfrak{A}}_{m}$, then $L[y]$ and $L^{\text {㐫 }}[z]$ are integrable.

Let $\mathscr{A}_{1}(x), \mathscr{A}_{2}(x), \mathscr{A}_{3}(x)$, and $\mathscr{A}_{4}(x)$ be $m \times n, m \times m, n \times n$, and $n \times m$ matrices, respectively, defined as follows:

$$
\begin{aligned}
& \mathscr{A}_{1}(x)=\left[p_{i j}(x)-p_{i n}(x) p_{m n}^{-1}(x) p_{m j}(x)\right], \\
& (i=0, \cdots, m-1 ; j=0, \cdots, n-1), \\
& \mathscr{A}_{2}(x)=\left[\begin{array}{cc}
0_{1_{m-1}} & p_{0 n}(x) p_{m n}^{-1}(x) \\
-E_{m-1} & p_{i n}(x) p_{m n}^{-1}(x)
\end{array}\right], \quad(i=1, \cdots, m-1), \\
& \mathscr{A}_{3}(x)=\left[\begin{array}{cc}
0_{n-11} & -E_{n-1} \\
p_{m n}^{-1}(x) p_{m 0}(x) & p_{m n}^{-1}(x) p_{m j}(x)
\end{array}\right], \quad(j=1, \cdots, n-1), \\
& \mathscr{A}_{4}(x)=\left[\begin{array}{ll}
0_{n-1 m-1} & 0_{n-11} \\
0_{1 m-1} & -p_{m n}^{-1}(x)
\end{array}\right] .
\end{aligned}
$$

If $f$ and $g$ belong to $\mathfrak{N}_{0}$, then the equation $L[y]=f$ is equivalent to the following system in the vector functions $\widehat{y}=\left(\widehat{y}_{i}\right),(i=1, \cdots, n)$, and $\widetilde{y}=\left(\widetilde{y}_{j}\right), \quad(j=1, \cdots, m)$ :

$$
\begin{align*}
& \hat{y}^{\prime}+\mathscr{A}_{3} \hat{y}+\mathscr{A}_{4} \widetilde{y}=0, \\
& \widetilde{y}^{\prime}-\mathscr{A}_{1} \hat{y}-\mathscr{A}_{2} \widetilde{y}=-f e^{(m, 1)} ; \tag{2.5}
\end{align*}
$$

and the equation $L^{\grave{\alpha}}[z]=g$ is equivalent to the following system in the vector functions $\check{z}=\left(\check{z}_{j}\right),(j=1, \cdots, m)$, and $\widetilde{z}=\left(\widetilde{z}_{j}\right),(i=1, \cdots, n)$ :

$$
\begin{align*}
& \check{z}^{\prime}+\mathscr{A}_{2}^{*} * \check{z}+\mathscr{A}_{4}^{*} \widetilde{z}=0, \\
& \widetilde{z}^{\prime}-\mathscr{A}_{1} * \check{z}-\mathscr{A}_{3} * \widetilde{z}=-g e^{(n, 1)}, \tag{2.6}
\end{align*}
$$

where $e^{(k, 1)},(k=1,2,3, \cdots)$, is used to denote the $k$-dimensional vector whose first coordinate is one, and whose remaining coordinates are zero. If $\mathscr{J}$ is the $(m+n) \times(m+n)$ matrix

$$
\mathscr{J}=\left[\begin{array}{cc}
0_{m n} & -E_{m}  \tag{2.7}\\
E_{n} & 0_{n m}
\end{array}\right]
$$

and $\mathscr{A}$ is the $(m+n) \times(m+n)$ matrix function defined by

$$
\mathscr{A}(x)=\left[\begin{array}{ll}
\mathscr{A}_{1}(x) & \mathscr{A}_{2}(x) \\
\mathscr{A}_{3}(x) & \mathscr{A}_{4}(x)
\end{array}\right],
$$

then (2.5) and (2.6) may be written as

$$
\begin{equation*}
\mathscr{L}[\hat{y}] \equiv \mathscr{J} \hat{y}^{\prime}+\mathscr{A} \hat{y}=f e^{(m+n, 1)} \tag{2.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathscr{L}^{\star}[\bar{z}] \equiv-\mathscr{J}^{*} \bar{z}+\mathscr{A}^{*} \bar{z}=g e^{(m+n, 1)} \tag{2.9}
\end{equation*}
$$

respectively.
Theorems on existence and uniqueness of solutions of $L[y]=f$ and $L^{\star}[z]=g$ follow from corresponding theorems for the respective first order systems (2.8) and (2.9). It also follows that $y \in \tilde{\mathbb{N}}_{n}$ if and only if there exists an integrable function $f$ such that $y$ is the first coordinate of a vector function $\widehat{y}$ satisfying (2.8), and $z \in \widetilde{\mathfrak{A}}_{m}$ if and only if there is an integrable function $g$ such that $z$ is the first coordinate of a vector function $\bar{z}$ satisfying (2.9).

The differential system (2.5) is identically normal in the sense that if $\bar{y}(x)$ is a solution of $\mathscr{L}[\hat{y}]=0$ with $\widehat{y}(x) \equiv 0$ on a subinterval $X$ of $[a, b]$, then $\bar{y}(x) \equiv 0$ on $X$. Indeed, if $\hat{y}$ is such a solution of (2.5), then $\widetilde{y}$ is a solution of $\widetilde{y}^{\prime}-\mathscr{A}_{2} \tilde{y}=0$ satisfying $\mathscr{A}_{4} \tilde{y}=0$ on $X$. This latter condition implies that $\widetilde{y}_{m}(x) \equiv 0$ on this subinterval, and the differential equation $\widetilde{y}^{\prime}-\mathscr{A}_{2} \widetilde{y}=0$ implies in turn that $\widetilde{y}_{j}(x) \equiv 0$ on $X$ for $j=m-1, \cdots, 1$. Similarly, system (2.6) is also identically normal. It follows from the identical normality of (2.5) that functions $y_{\alpha}$ in $\widetilde{\Re}_{n}$ are linearly independent solutions of $L[y]=0$ if and only if the corresponding vector functions $\widehat{y}_{\alpha}$ are linearly independent solutions of $\mathscr{P}[\hat{y}]=0$. Similarly, it follows from the identical normality of (2.6) that functions $z_{\alpha}$ in $\tilde{\mathfrak{A}}_{m}$ are linearly independent solutions of $L^{\text {匂 }}[z]=0$ if and only if the corresponding vector functions $\bar{z}_{\alpha}$ are linearly independent solutions of $\mathscr{L}^{\bar{\lambda}}[\bar{z}]=0$.
3. The adjoint operator. If $\mathscr{J}$ is the $(m+n) \times(m+n)$ matrix defined as in (2.7), then we may establish the following Lagrange identity by a simple inductive argument which does not use hypothesis (H).

Lemma 3.1. If $y \in \tilde{\mathfrak{A}}_{n}$ and $z \in \tilde{\mathfrak{A}}_{m}$, then

$$
\begin{equation*}
\bar{z} L[y]-\bar{L}^{\grave{ }}[z] y=\left(\bar{z}^{*} \mathscr{J} \bar{y}\right)^{\prime} \tag{3.1}
\end{equation*}
$$

Theorem 3.1. If $f \in \mathfrak{A}_{0}$, then there exists a $y$ in $\widetilde{\mathfrak{Q}}_{n}^{0}$ such that $L[y]=f$ if and only if $z$ in $\widetilde{\mathfrak{A}}_{m}$ and $L^{\natural}[z]=0$ implies that $(f, z)=0$.

Now if $y \in \tilde{\mathfrak{A}}_{n}^{0}, L[y]=f, z \in \tilde{\mathfrak{A}}_{m}$, and $L^{\star}[z]=0$, then, in view of Lemma 3.1,

$$
(f, z)=(L[y], z)-\left(y, L^{\star}[z]\right)=\left.\breve{z}^{*} \mathcal{J} \hat{y}\right|_{a} ^{b}=0 .
$$

On the other hand, suppose that $(f, z)=0$ whenever $z \in \tilde{\mathfrak{A}}_{m}$ and $L^{\star}[z]=0$, and let $y$ be the function in $\tilde{\mathfrak{N}}_{n}$ such that $L[y]=f$ and $\bar{y}(a)=0$. If $z_{j},(j=1, \cdots, m+n)$ are linearly independent solutions of $L^{\star}[z]=0$, then the $(m+n) \times(m+n)$ matrix $\bar{Z}(x)$ with column vectors $\breve{z}_{j}(x), \quad(j=1, \cdots, m+n)$, is nonsingular on $[a, b]$. From Lemma 3.1 we have the vector equation

$$
0=\left[\left(f, z_{j}\right)-\left(y, L^{\star}\left[z_{j}\right]\right)\right]=\left.\bar{Z}^{*} \mathscr{J} \hat{y}\right|_{a} ^{b}=\check{Z}^{*}(b) \mathscr{J} \hat{y}(b),
$$

and consequently $\hat{y}(b)=0$ also.
Theorem 3.2. If hypothesis (H) holds, then $\mathscr{D}_{0}{ }^{*}=\overline{\mathfrak{V}}_{m}$ and $f_{z}=$ $L^{\star}[z]$ on $\mathscr{D}_{0}{ }^{*}$.

That $\tilde{\mathfrak{U}}_{m} \subset \mathscr{D}_{0}{ }^{*}$ follows from Lemma 3.1. Now let $z_{0} \in \mathscr{D}_{0}{ }^{*}$ and suppose $f_{z_{0}}$ is a corresponding function in $\mathfrak{Y}_{0}$ such that $\left(L[y], z_{0}\right)=$ $\left(y, f_{z_{0}}\right)$ when $y \in \tilde{\mathfrak{Z}}_{n}^{0}$. Choose $w_{0}$ in $\tilde{\mathfrak{U}}_{m}$ such that $L^{\star}\left[w_{0}\right]=f_{z_{0}}$, and suppose that $z_{i} \in \widetilde{\mathfrak{Y}}_{m}$ are linearly independent solutions of $L^{\star}\left[z_{i}\right]=0$, with $\left(z_{i}, z_{j}\right)=\delta_{i j},(i, j=1, \cdots, m+n)$. If $w=w_{0}+\sum_{j=1}^{m+n}\left(z_{0}-w_{0}, z_{j}\right) z_{j}$, then $L^{\star}[w]=f_{z 0}$ and $\left(z_{0}-w, z\right)=0$ when $z \in \widetilde{\mathfrak{A}}_{m}$ and $L^{\star}[z]=0$. It follows that if $y \in \tilde{\mathfrak{A}}_{n}^{0}$, then

$$
\begin{equation*}
\left(L[y], z_{0}\right)=\left(y, f_{z_{0}}\right)=\left(y, L^{\star}[w]\right)=(L[y], w), \tag{3.2}
\end{equation*}
$$

so that $\left(L[y], z_{0}-w\right)=0$ when $y \in \widetilde{\mathfrak{M}}_{n}^{0}$. But it follows from Theorem 3.1 that there is a function $y$ in $\tilde{\mathscr{A}}_{a}^{0}$ such that $L[y]=z_{0}-w$. Consequently $\left(z_{0}-w, z_{0}-w\right)=0$ and $z_{0}=w \in \widetilde{\mathfrak{A}}_{m}$, so that $\mathscr{D}_{0}{ }^{*}=\widetilde{\mathscr{A}}_{m}$ and $f_{2_{0}}=L^{*}\left[z_{0}\right]$. This result extends Theorem 4.1 of Reid [7].

Now the operator $T_{0}^{*}$ adjoint to $T_{0}$ is defined to be the operator on $\mathscr{D}_{0}^{*}$ with value $f_{z}$ at $z$. In view of Theorem 3.2 we have $\mathscr{D}_{0}^{*}=$ $\tilde{\mathfrak{A}}_{m}$ and $T_{0}^{*}[z]=L^{\star}[z]$.
4. Extensions of the operator $T_{0}$. Let $\mathscr{D}$ be a linear subspace of $\tilde{\mathfrak{U}}_{n}$ containing $\tilde{\mathfrak{M}}_{n}^{0}$, and denote by $T$ the restriction of $L$ to $\mathscr{D}$. Denote by $\mathscr{D}^{*}$ the class of functions $z$ in $\mathfrak{U}_{0}$ which are essentially bounded and for which there exists an $f_{z}$ in $\mathfrak{X}_{0}$ such that ( $L[y], z$ ) = $\left(y, f_{z}\right)$ for all $y$ in $\mathscr{D}$. It follows from Theorem 3.2 that $\mathscr{D}^{*} \subset \mathfrak{N}_{m}$ and for each $z$ in $\mathscr{D}^{*}$ there is at most one $f_{z}$, namely $L^{\star}[z]$, such that $(L[y], z)=\left(y, f_{z}\right)$ for all $y$ in $\mathscr{D}$. The adjoint $T^{*}$ of $T$ is the
operator on $\mathscr{D}^{*}$ defined by the formula $T^{*}[z]=f_{z}$. The operator $T$ is said to be self-adjoint if and only if $\mathscr{D}=\mathscr{D}^{*}$ and $T=T^{*}$.

The following lemma will be helpful in describing $\mathscr{D}^{*}$. If $y_{j} \in \widetilde{\mathfrak{A}}_{n}$, $(j=1, \cdots, m+n)$, then $\bar{Y}$ will denote the matrix function defined by $\hat{Y}(x) \equiv\left[\widehat{y}_{j}(x)\right],(j=1, \cdots, m+n)$.

Lemma 4.1. If $\eta$ and $\zeta$ are $(m+n)$-vectors, then there exists a function $y \in \widetilde{\mathfrak{A}}_{n},\left(z \in \widetilde{\mathfrak{M}}_{m}\right)$, such that $\bar{y}(a)=\eta$ and $\widehat{y}(b)=\zeta,(\bar{z}(a)=\eta$ $d a n \bar{z}(b)=\zeta)$.

Since $\widetilde{\mathfrak{A}}_{n}$ is a vector space it is enough to show that there exist $m+n$ functions $y_{j}$ in $\widetilde{\mathfrak{A}}_{n}$ such that $\bar{y}_{j}(a)=0, \quad(j=1, \cdots, m+n)$ while $\bar{Y}(b)$ is nonsingular, and to show a corresponding result with $a$ and $b$ interchanged. To establish the existence of functions $y_{j}$ in $\widetilde{\mathfrak{A}}_{n}$ such that $\widehat{y}_{j}(a)=0, \quad(j=1, \cdots, m+n)$, and $\bar{Y}(b)$ is nonsingular, suppose to the contrary that for each collection of $m+n$ functions $y_{j}$ in $\widetilde{\mathfrak{M}}_{n}$ satisfying $\bar{y}_{j}(\alpha)=0,(j=1, \cdots, m+n)$, we have $\bar{Y}(b)$ singular. Let $z_{j}$ be $m+n$ linearly independent solutions of $L^{\text {漓 }}[z]=0$, and for $j=1, \cdots, m+n$ let $y_{j}$ be the function in $\widetilde{\mathfrak{A}}_{n}$ such that $L\left[y_{j}\right]=z_{j}$ and $\bar{y}_{j}(\alpha)=0$. Then there is a nonzero $(m+n)$-vector $\xi=\left(\xi_{j}\right)$ such that $\hat{Y}(b) \xi=0$. If $y(x)=\sum_{j=1}^{m+n} y_{j}(x) \xi_{j}$ and $z(x)=\sum_{j=1}^{m+n} z_{j}(x) \xi_{j}$, then $L[y]=z, L^{\hbar}[z]=0$ and $z(x) \not \equiv 0$, moreover, $y \in \widetilde{\mathfrak{A}}_{n}^{0}$. Hence it follows from Lemma 3.1 that

$$
0=(L[y], z)-\left(y, L^{\text {® }}[z]\right)=(z, z),
$$

which is impossible since $z(x) \not \equiv 0$. The numbers $a$ and $b$ may be interchanged and the preceding argument remains valid. The result for $\tilde{\mathfrak{A}}_{m}$ follows by interchanging the roles of $\tilde{\mathfrak{A}}_{n}$ and $\tilde{\mathfrak{A}}_{m}$, that is, by replacing $\left[p_{i j}\right]$ with $\left[p_{i j}\right]^{*}$.

Denote by $\mathscr{B}$ the subspace of $2(m+n)$-dimensional complex space consisting of the end values ( $\hat{y}(a), \widetilde{y}(a), \widehat{y}(b), \tilde{y}(b))$ for functions $y$ in $\mathscr{D}$. Similarly, $\mathscr{B}^{*}$ will denote the subspace of end values $(\check{z}(a), \widetilde{z}(a), \check{z}(b), \widetilde{z}(b))$ for functions $z$ in $\mathscr{D}^{*}$. If $k<2 m+2 n$ and the dimension of $\mathscr{B}$ is $2 m+2 n-k$, then let $P$ and $Q$ be $(m+n) \times(2 m+2 n-k)$ matrices such that the columns of $\left[-P^{*} Q^{*}\right]^{*}$ form a basis for $\mathscr{B}$. If $k>0$ also, then let $M$ and $N$ be $k \times(m+n)$ matrices such that the $k \times 2(m+n)$ matrix [ $M N$ ] has rank $k$ and $M P-N Q=0$. Then in view of Lemma 4.1 we have that $\mathscr{D}$ is characterized as the class of functions $y$ in $\widetilde{\mathcal{A}}_{n}$ with the property that

$$
\begin{equation*}
s(\widehat{y}) \equiv M \widehat{y}(a)+N \widehat{y}(b)=0 \tag{4.1}
\end{equation*}
$$

If $k=0$, then by Lemma 4.1 we have $\mathscr{D}=\widetilde{\mathfrak{A}}_{n}$.

Theorem 4.1. $\operatorname{Dim} \mathscr{B}+\operatorname{dim} \mathscr{B}^{*}=2 m+2 n$; if $\operatorname{dim} \mathscr{B}>0$ and $P, Q$ are $(m+n) \times(2 m+2 n-k)$ matrices such that the column vectors of $\left[-P^{*} Q^{*}\right]^{*}$ from a basis for $\mathscr{B}$, then $\mathscr{D}^{*}$ is the class of functions $z$ in $\tilde{\mathfrak{A}}_{m}$ for which

$$
\begin{equation*}
P^{*} \mathscr{J}^{*} \bar{z}(a)+Q^{*} \mathscr{J}^{*} \bar{z}(b)=0 \tag{4.2}
\end{equation*}
$$

First note that if $\operatorname{dim} \mathscr{B}=0$, then $\mathscr{D}^{*}=\widetilde{\mathfrak{A}}_{m}$ by Theorem 3.2, and thus by Lemma 4.1 we have $\operatorname{dim} \mathscr{B}^{*}=2 m+2 n$. Now suppose that $\operatorname{dim} \mathscr{B}>0, z \in \tilde{\mathfrak{A}}_{m}$, and (4.2) holds. Then for $y$ in $\mathscr{D}$ and $\xi$ a $(2 m+2 n-k)$-vector chosen so that $\bar{y}(a)=-P \xi$ and $\bar{y}(b)=Q \xi$ it follows from Lemma 3.1 that

$$
(L[y], z)-\left(y, L^{\text {亠 }}[z]\right)=\left.\bar{z}^{*} \mathscr{J} \bar{y}\right|_{a} ^{b}=\left\{P^{*} \mathscr{J}^{*} \breve{z}(a)+Q^{*} \mathscr{J}^{*} \breve{z}(b)\right\}^{*} \xi=0
$$

and hence $z \in \mathscr{D}^{*}$. On the other hand, if $z \in \mathscr{D}^{*}$ then it follows from Theorem 3.2 that $z \in \tilde{\mathfrak{A}}_{m}$, since $\tilde{\mathfrak{M}}_{n}^{0} \subset \mathscr{D}$. Then (4.2) follows from Lemma 3.1, Lemma 4.1 and the choice of $P$ and $Q$. Therefore, in view of Lemma 4.1, it follows that $\operatorname{dim} \mathscr{B}+\operatorname{dim} \mathscr{B}^{*}=2 m+2 n$.

Corollary I. If $\operatorname{dim} \mathscr{B}>0$, and $R$ and $S$ are $(2 m+2 n-k) \times$ $(m+n)$ matrices, then $\mathscr{D}^{*}$ is the collection of functions $z$ in $\widetilde{\mathfrak{U}}_{m}$ for which

$$
\begin{equation*}
R \bar{z}(a)+S \breve{z}(b)=0 \tag{4.3}
\end{equation*}
$$

if and only if the $(2 m+2 n-k) \times 2(m+n)$ matrix $[R S]$ has rank $2 m+2 n-k$ and $M \mathscr{J}^{*} R^{*}-N \cdot \mathscr{J}^{*} S^{*}=0$.

Corollary II. The adjoint of $T^{*}$ is $T$.
The index of compatibility for a system $L[y]=0, y \in \mathscr{D}$ is defined to be $\operatorname{dim}\{y: y \in \mathscr{D}$ and $L[y]=0\}$. The next two theorems are consequences of the equivalence of the equations $L[y]=f$ and $L^{\text {吝 }}[z]=g$ to the systems (2.8) and (2.9), respectively, and corresponding theorems on first order systems. Analogous theorems for $n$th order linear differential equations are given in [2, Chapter 11], and those results may be extended to first order systems.

Theorem 4.2. If $\operatorname{dim} \mathscr{B}^{*}=k$ and the index of compatibility of the system $L[y]=0, y \in \mathscr{D}$ is $r$, then $\rho=k+r-m-n$ is the index of compatibility for the system $L^{\natural}[z]=0, z \in \mathscr{D}^{*}$.

THEOREM 4.3. If $f \in \mathfrak{N}_{0}$, then there exists a function $y$ in $\mathscr{D}$ such that $L[y]=f$ if and only if $(f, z)=0$ for all $z$ in $\mathscr{D}^{*}$ satisfying $L^{\star}[z]=0$.

The next two theorems are analogues of Theorems 6.1 and 6.2 of Reid [7]. The second of the two gives necessary and sufficient conditions for the operator $T$ to be self-adjoint when $\left[p_{i j}(x)\right]$ is Hermitian. If $y_{j} \in \widetilde{\mathfrak{N}}_{n}$ and $\bar{Y}=\left[\widehat{y}_{j}\right], \quad(j=1, \cdots, m+n)$, then the symbols $s(\bar{Y})$ and $s^{-}(\hat{Y})$ are used for the $k \times(m+n)$ matrices $M \hat{Y}(a)+N \hat{Y}(b)$ and $M \hat{Y}(a)-N \hat{Y}(b)$, respectively. Similarly, if $z_{j} \in \widetilde{\mathscr{Z}}_{m}$ and $\breve{Z}=\left[\bar{z}_{j}\right]$, $(j=1, \cdots, m+n)$, then $t(\breve{Z})$ and $t^{-}(\breve{Z})$ denote $R \check{Z}(a)+S \check{Z}(b)$ and $R \check{Z}(a)-S \breve{Z}(b)$, respectively.

TheORem 4.4. Suppose that $2(m+n)>\operatorname{dim} \mathscr{B}>0, y_{j}$ and $z_{j}$, $(j=1, \cdots, m+n)$, are linearly independent solutions of $L[y]=0$ and $L^{\star}[z]=0$, respectively, and let $\Delta=\left(\breve{Z}^{*} \mathscr{J} \bar{Y}\right)^{-1}$. Then $\Delta$ is constant on $[a, b]$ and $\mathscr{D}^{*}$ is the collection of functions $z$ in $\widetilde{\mathfrak{A}}_{m}$ satisfying (4.3) if and only if the $(2 m+2 n-k) \times 2(m+n)$ matrix $[R S]$ has rank $2 m+2 n-k$ and

$$
\begin{equation*}
s(\hat{Y}) \Delta\left\{t^{-}(\breve{Z})\right\}^{*}+s^{-}(\hat{Y}) \Delta\{t(\breve{Z})\}^{*}=0 \tag{4.4}
\end{equation*}
$$

Theorem 4.5. Suppose that $m=n, \quad\left[p_{i j}(x)\right], \quad(i, j=0, \cdots, n$; $x \in[a, b])$, is Hermitian and $\operatorname{dim} \mathscr{B}=2 n$. Let $y_{j},(j=1, \cdots, 2 n)$, be linearly independent solutions of $L[y]=0$, and let $\Delta=\left(\hat{Y}^{*} \mathscr{J} \bar{Y}\right)^{-1}$. Then $\Delta$ is constant on $[a, b]$, and $T$ is self-adjoint if and only if the $2 n \times 2 n$ matrix $s^{-}(\hat{Y}) \Delta\{s(\bar{Y})\}^{*}$ is Hermitian.
5. Generalized Green's functions. The subspaces $\mathscr{D}, \mathscr{D}^{*}$ of $\tilde{\mathfrak{A}}_{n}$ and $\tilde{\mathfrak{A}}_{m}$, respectively, and the subspaces $\mathscr{B}, \mathscr{B}^{*}$ of $2(m+n)$ dimensional complex space are as defined in §4. If $0<\operatorname{dim} \mathscr{B}<2 m+$ $2 n$, then the matrices $M, N, P$, and $Q$ are as specified in $\S 4$.

If $f \in \mathfrak{N}_{0}$ then we are concerned with solutions of the quasidifferential system

$$
\begin{equation*}
L[y]=f, \quad y \in \mathscr{O} . \tag{5.1}
\end{equation*}
$$

Of prime importance is the homogeneous system

$$
\begin{equation*}
L[y]=0, \quad y \in \mathscr{D} \tag{5.2}
\end{equation*}
$$

and its adjoint system

$$
\begin{equation*}
L^{\text {漓 }}[z]=0, \quad z \in \mathscr{D}^{*} . \tag{5.3}
\end{equation*}
$$

By definition a generalized Green's function for the system (5.2) is an essentially bounded and measurable function $g$ on $\square \equiv$ $\{(x, t): a \leqq x \leqq b, a \leqq t \leqq b\}$ with the property that if $f$ is $a$ function in $\mathfrak{N}_{0}$ for which (5.1) has a solution, then a particular solution $y$
of (5.1) is given by

$$
\begin{equation*}
y(x)=\int_{a}^{b} g(x, t) f(t) d t \tag{5.4}
\end{equation*}
$$

Reid [5] has shown the existence of a generalized Green's matrix for a compatible first order system with two-point boundary conditions, where the number of independent boundary conditions is equal to the number of differential equations. If $\operatorname{dim} \mathscr{B}=m+n$, then Reid's results could be used to obtain a generalized Green's function for (5.2). In this section the existence and some properties of a generalized Green's function will be shown when $\operatorname{dim} \mathscr{B}$ is not necessarily equal to $m+n$. The technique used here may be modified to extend Reid's results to the case where the number of independent boundary conditions is different from the number of differential equations.

For a $\nu$ th order linear differential operator $\sum_{j=0}^{y} q_{j}(x) y^{(j)}$ with $q_{j} \in C_{j},(j=0,1, \cdots, \nu)$, and $q_{\nu}(x) \neq 0$, the generalized Green's function has been treated by Greub and Rheinboldt [4] and Wyler [10]; a more comprehensive treatment of an algebraic theory of operator solutions of boundary problems, which includes this case as a special instance, is given in Wyler [11].

Lemma 5.1. If $y_{j},(j=1, \cdots, m+n)$, are linearly independent solutions of $L[y]=0$, then there exist $m+n$ linearly independent solutions $z_{j}$ of $L^{\star}[z]=0$ such that

$$
\begin{equation*}
\bar{Z}^{*} \mathscr{J} \hat{Y}=E_{m+n} \tag{5.5}
\end{equation*}
$$

This result follows from Lemma 3.1 and the existence and uniqueness theorems for the equations $\mathscr{L}[\hat{y}]=0$ and $\mathscr{L}^{\psi}[\tilde{z}]=0$.

If $y_{j} \in \widetilde{\mathfrak{A}}_{n}$ and $z_{j} \in \tilde{\mathfrak{A}}_{m}, \quad(j=1, \cdots, m+n)$, then define matrix functions $\hat{Y}, \tilde{Y}, \check{Z}$, and $\widetilde{Z}$ as follows: $\hat{Y}(x)=\left[\hat{y}_{j}(x)\right], \tilde{Y}(x)=\left[\widetilde{y}_{j}(x)\right]$, $\check{Z}(x)=\left[\check{z}_{j}(x)\right]$, and $\widetilde{Z}(x)=\left[\widetilde{z}_{j}(x)\right],(j=1, \cdots, m+n)$.

Corollary. If $y_{j}$ and $z_{j},(j=1, \cdots, m+n)$, are as in Lemma 5.1, then

$$
\begin{array}{ll}
\hat{Y}(x) \check{Z}^{*}(x) \equiv 0_{n m}, & \hat{Y}(x) \widetilde{Z}^{*}(x) \equiv E_{n} \\
\widetilde{Y}(x) \check{Z}^{*}(x) \equiv-E_{m}, & \widetilde{Y}(x) \widetilde{Z}^{*}(x) \equiv 0_{m n} \tag{5.6}
\end{array}
$$

Theorem 5.1. If $\tau \in[a, b], \xi_{j}$ is $a$ constant, $y_{j}$ and $z_{j},(j=1, \cdots$, $m+n$ ), are as in Lemma 5.1, then the solution $y$ of $L[y]=f$ satisfying $\widehat{y}(\tau)=\sum_{j=1}^{m+n} \widehat{y}_{j}(\tau) \xi_{j}$ is given by the first component of the vector

$$
\begin{equation*}
\widehat{y}(x)=\sum_{j=1}^{m+n} \widehat{y}_{j}(x) \xi_{j}+\int_{\tau}^{x} \sum_{j=1}^{m+n} \bar{y}_{j}(x) \bar{z}_{j}(t) f(t) d t . \tag{5.7}
\end{equation*}
$$

Indeed, if $\xi=\left(\xi_{j}\right),(j=1, \cdots, m+n)$, and we set $\widehat{y}(x)=\hat{Y}(x) u(x)$, for $u$ an $(m+n)$-vector function, then $\hat{y}$ is a solution of $\mathscr{S}[\hat{y}]=f e^{\left(m+n, x^{1}\right)}$, $\bar{y}(\tau)=\bar{Y}(\tau) \xi$ if and only if

$$
\mathscr{J} \bar{Y}(x) u^{\prime}(x)=e^{(m+n, 1)} f(x), \quad u(\tau)=\xi .
$$

Hence $u^{\prime}(x)=\breve{Z}^{*}(x) e^{(m+n, 1)} f(x)$ and

$$
u(x)=\xi+\int_{\tau}^{x} \breve{Z}^{*}(s) e^{(m+n, 1)} f(s) d s
$$

from which the theorem follows.
Now suppose that $y_{j},(j=1, \cdots, m+n)$, are linearly independent solutions of $L[y]=0$ and that $z_{j},(j=1, \cdots, m+n)$, are chosen as in Lemma 5.1. If $\operatorname{dim} \mathscr{B}=2 m+2 n-k, k>0$, then $s(\hat{Y})$ and $s^{-}(\hat{Y})$ are $k \times(m+n)$ matrices defined as $s(\hat{Y})=M \bar{Y}(a)+N \bar{Y}(b)$ and $s^{-}(\hat{Y})=M \hat{Y}(a)-N \hat{Y}(b)$. If $r$ is the index of compatibility for (5.2), then $s(\hat{Y})$ has rank $m+n-r$. If $r>0$, then let $S$ be an $(m+n) \times r$ matrix with the property that $S^{*} S=E_{r}$ and $s(\bar{Y}) S=0$. If $r>m+n-k$, then $T$ will represent a $k \times(k-m-n+r)$ matrix such that $T^{*} T=$ $E_{k-m-n+r}$ and $T^{*} s(\hat{Y})=0$. It follows that the $(k+r) \times(k+r)$ matrix

$$
\left[\begin{array}{cc}
s(\hat{Y}) & T  \tag{5.8}\\
S^{*} & 0
\end{array}\right]
$$

is nonsingular, and its inverse is of the form

$$
\left[\begin{array}{ll}
D & S  \tag{5.9}\\
T^{*} & 0
\end{array}\right]
$$

The $(m+n) \times k$ matrix $D$ is the generalized reciprocal of $s(\hat{Y})$ in the sense of E. H. Moore, (see [9, Section 14]). If $r=0$, then the matrix $S$ does not appear, if $r=m+n-k$, then $T$ does not appear.

Now if $\operatorname{dim} \mathscr{B}<2(m+n)$, let $G(x, t)$ be the $(m+n) \times(m+n)$ matrix defined by

$$
\begin{aligned}
& G(x, t)=\frac{1}{2} \hat{Y}(x)\left[\frac{|x-t|}{x-t} E_{m+n}+D s^{-}(\hat{Y})\right] \breve{Z}^{*}(t), \quad x \neq t ; \\
& G(x, x)=\frac{1}{2} \hat{Y}(x) D s^{-}(\widehat{Y}) \breve{Z}^{*}(x), \quad x \in[a, b] .
\end{aligned}
$$

If $\operatorname{dim} \mathscr{B}=2(m+n)$, let $G(x, t)$ be defined by

$$
\begin{aligned}
& G(x, t)=\frac{1}{2} \frac{|x-t|}{x-t} \widehat{Y}(x) \breve{Z}^{*}(t), \quad x \neq t \\
& G(x, x)=0, \quad x \in[a, b]
\end{aligned}
$$

Let $g_{0}$ be the function with domain $\square$ whose value at $(x, t)$ is the element in the first row and first column of $G(x, t)$, that is

$$
\begin{array}{ll}
g_{0}(x, t)=g_{0,1}(x, t)+g_{0,2}(x, t) & \text { if } \operatorname{dim} \mathscr{B}<2(m+n), \\
g_{0}(x, t)=g_{0,1}(x, t) & \text { if } \operatorname{dim} \mathscr{B}=2(m+n),
\end{array}
$$

where

$$
\begin{aligned}
& g_{0,1}(x, t)=\frac{1}{2} \operatorname{sgn}(x-t) \sum_{i=1}^{m+n} y_{i}(x) \bar{z}_{i}(t), \\
& g_{0,2}(x, t)=\frac{1}{2} \sum_{i, j=0}^{m+n} y_{i}(x) \mathscr{K}_{i j} \bar{z}_{j}(t),
\end{aligned}
$$

provided $\left[\mathscr{K}_{i j}\right]$ is the matrix $D s^{-}(\hat{Y})$ and $\operatorname{sgn} u=|u| / u$ for $u \neq 0$, $\operatorname{sgn} 0=0$.

Theorem 5.2. The function $g_{0}$ defined above is a generalized Green's function for (5.2).

If $\operatorname{dim} \mathscr{B}=2(m+n)$, then this result follows directly from Theorem 5.1. Now suppose that $\operatorname{dim} \mathscr{B}<2(m+n)$, and $f$ is an integrable function for which (5.1) has a solution. If $y$ is a solution of $L[y]=f$, then for a suitable vector $\xi$ one has

$$
\widehat{y}(x)=\frac{1}{2}\left[\bar{Y}(x) \xi+\int_{a}^{x} \widehat{Y}(x) \breve{Z}^{*}(t) e^{(m+n, 1)} f(t) d t-\int_{x}^{b} \bar{Y}(x) \breve{Z}^{*}(t) e^{(m+n, 1)} f(t) d t\right]
$$

Thus, since (5.9) is the inverse of (5.8), it follows that $y$ is a solution of (5.1) if and only if

$$
T^{*} s^{-}(\widehat{Y}) \int_{a}^{b} \breve{Z}^{*}(t) e^{(m+n, 1)} f(t) d t=0
$$

and for some $r$-vector $\eta$ we have

$$
\xi=D s^{-}(\widehat{Y}) \int_{a}^{b} \breve{Z}^{*}(t) e^{(m+n, 1)} f(t) d t+S \eta
$$

Therefore,

$$
\begin{aligned}
\bar{y}(x)=\frac{1}{2}\left[\bar{Y}(x) S \eta+\bar{Y}(x) D s^{-}(\bar{Y})\right. & \int_{a}^{b} \breve{Z}^{*}(t) e^{(m+n, 1)} f(t) d t \\
& \left.+\int_{a}^{b} \bar{Y}(x) \frac{|x-t|}{x-t} \breve{Z}^{*}(t) e^{(m+n, 1)} f(t) d t\right]
\end{aligned}
$$

from which the theorem follows since $\eta$ may be chosen to be zero.
The symbol $g_{0}^{(i, j)}$ will be used to signify the partial derivative $\partial^{i+j} g_{0} / \partial t^{j} \partial x^{i}$. Generalized partial derivatives $g_{0}^{\langle\alpha, \beta\rangle}$ will now be defined for $g_{0}$. If $\alpha<n$ and $\beta<m$, then $g_{0}^{\langle\alpha, \beta\rangle}(x, t)=g_{0}^{\langle\alpha, \beta\rangle}(x, t)$. If $\alpha<n$, then $g_{0}^{\langle\alpha, m+j\rangle},(j=0, \cdots, n-1)$, is defined as follows:

$$
g_{0}^{\langle\alpha, m\rangle}(x, t)=\sum_{i=0}^{m} \bar{p}_{i n}(t) g_{0}^{(\alpha, i)}(x, t) ;
$$

if $g^{\langle\alpha, m-1+j\rangle}$ is a.c. in its second argument, then

$$
\begin{aligned}
g_{0}^{\langle\alpha, m+j\rangle}(x, t)=\sum_{i=0}^{m} \bar{p}_{i n-j}(t) g_{0}^{(\alpha, i\rangle}(x, t)-\partial / \partial t & g_{0}^{\langle\alpha, m-1+j\rangle}(x, t), \\
& (j=1, \cdots, n-1) .
\end{aligned}
$$

If $\beta<m$, then $g_{0}^{\langle n+i, \beta\rangle},(i=0, \cdots, m-1)$, is defined as follows:

$$
g_{0}^{\langle n, \beta\rangle}(x, t)=\sum_{j=0}^{n} p_{m j}(x) g_{0}^{(j, \beta)}(x, t) ;
$$

if $g^{\langle n-1+i, \beta\rangle}$ is a.c. in its first argument, then

$$
\begin{aligned}
g_{0}^{\langle n+i, \beta\rangle}(x, t)=\sum_{j=1}^{n} p_{m-i j}(x) g_{0}^{(j, \beta)}(x, t)-\partial / \partial x & g_{0}^{\langle n-1+i, \beta\rangle}(x, t) \\
& (i=1, \cdots, m-1) .
\end{aligned}
$$

Theorem 5.3. If $\alpha+\beta \leqq m+n-2$, and $g_{0}$ is the function of Theorem 5.2, then $g_{0}^{\langle\alpha, \beta\rangle}$ exists and is continuous on $\square$.

This result clearly holds for $g_{0,2}$, hence one need only consider specifically $g_{0,1}$. Let $\alpha+\beta \leqq m+n-2$, and suppose first that $\alpha<n$. If $\beta<m$, then the theorem follows from the fact that $\hat{Y}(x) \check{Z}^{*}(x) \equiv 0$. If $\beta=m-1+j, \quad(j=1, \cdots, n-\alpha-1)$, then use the identity $\hat{Y}(x) \widetilde{Z}^{*}(x) \equiv E_{m}$. On the other hand, if $\beta<m$ and $\alpha=n-1+i$, $(i=1, \cdots, m-\beta-1)$, then use the identity $\tilde{Y}(x) \check{Z}^{*}(x) \equiv-E_{m}$.

Theorem 5.4. The generalized Green's function for the system (5.2) is not unique. If $u_{1}, \cdots, u_{r}$ form a basis for the solutions of (5.2), $v_{1}, \cdots, v_{\rho}$ form a basis for the solutions of (5.3), and $g_{0}$ is one generalized Green's function for (5.2) then a function $g$ on $\square$ is also a generalized Green's function for (5.2) if and only if there exist essentially bounded and measurable functions $\psi_{1}, \cdots, \psi_{r}, \varphi_{1}, \cdots, \varphi_{p}$ such that if $(x, t) \in \square$, then

$$
\begin{equation*}
g(x, t)=g_{0}(x, t)+\sum_{i=1}^{r} u_{i}(x) \psi_{i}(t)+\sum_{j=1}^{\rho} \varphi_{j}(x) \bar{v}_{j}(t) \tag{5.10}
\end{equation*}
$$

If $g$ is a function on $\square$ satisfying (5.10), then in view of Theorem
4.3 it follows that $g$ is a generalized Green's function for (5.2).

To establish the converse we may assume without loss of generality that $\left(u_{i}, u_{j}\right)=\delta_{i j},(i, j=1, \cdots, r)$, and $\left(v_{\alpha}, v_{\beta}\right)=\delta_{\alpha \beta},(\alpha, \beta=1, \cdots, \rho)$. If $w \in \mathfrak{H}_{0}$ and $f(x)=w(x)-\sum_{j=1}^{\rho}\left(w, v_{j}\right) v_{j}(x)$, then $\left(f, v_{\alpha}\right)=0, \quad(\alpha=$ $1, \cdots, \rho)$. Thus for this choice of $f$ it follows from Theorem 4.3 that (5.1) has a solution. Suppose that $g$ is a second generalized Green's function for (5.2) and let $d(x, t)=g(x, t)-g_{0}(x, t)$. Then there are constants $\xi_{1}, \cdots, \xi_{r}$ such that

$$
\int_{a}^{b} d(x, t) f(t) d t=\sum_{i=1}^{r} u_{i}(x) \xi_{i},
$$

and if $\Phi(x, t)=d(x, t)-\sum_{j=1}^{p} \bar{v}_{j}(t) \int_{a}^{b} d(x, s) v_{j}(s) d s$, then

$$
\begin{equation*}
\int_{a}^{b} \Phi(x, t) f(t) d t=\sum_{i=1}^{r} u_{i}(x) \xi_{i} \tag{5.11}
\end{equation*}
$$

Multiplying (5.11) by $\bar{u}_{i}(x)$, and integrating with respect to $x$, we have

$$
\int_{a}^{b} \int_{a}^{b} \bar{u}_{i}(x) \Phi(x, t) f(t) d t d x=\xi_{i}, \quad(i=1, \cdots, r)
$$

and consequently

$$
\int_{a}^{b}\left[\Phi(x, t)-\sum_{i=1}^{r} u_{i}(x) \int_{a}^{b} \bar{u}_{i}(s) \Phi(s, t) d s\right] w(t) d t=0
$$

But $w$ is an arbitrary integrable function, and hence

$$
\Phi(x, t)-\sum_{i=1}^{r} u_{i}(x) \int_{a}^{b} \bar{u}_{i}(s) \Phi(s, t) d s=0 \quad \text { on } \square
$$

and

$$
d(x, t)=\sum_{i=1}^{r} u_{i}(x) \int_{a}^{b} \bar{u}_{i}(s) \Phi(s, t) d s+\sum_{j=1}^{\infty} \bar{v}_{j}(t) \int_{a}^{b} d(x, s) v_{j}(s) d s .
$$

Hence (5.10) holds with $\psi_{i}$ and $\varphi_{j}$ defined by $\psi_{i}(t)=\int_{a}^{b} u_{i}(s) \Phi(s, t) d s$. and $\varphi_{j}(x)=\int_{a}^{b} d(x, s) v_{j}(s) d s, \quad(i=1, \cdots, r ; j=1, \cdots, \rho)$, and clearly these functions are essentially bounded and measurable.

We now show that a generalized Green's function $g$ for (5.2) has the property that the function $h$ defined by $h(x, t)=\bar{g}(t, x)$ is a generalized Green's function for the adjoint system (5.3). Preliminary to this result we shall prove the following theorem.

Theorem 5.5. Suppose that $u_{1}, \cdots, u_{r}$ form a basis for the solutions of (5.2), $v_{1}, \cdots, v_{\rho}$ from a basis for the solutions of (5.3), and $\Theta=\left\{\theta_{1}, \cdots, \theta_{r}\right\}, \Omega=\left\{\omega_{1}, \cdots, \omega_{\rho}\right\}$ are sets of integrable functions
with the property that the matrices $\left[\left(u_{i}, \theta_{j}\right)\right],(i, j=1, \cdots, r)$, and $\left[\left(v_{\alpha}, \omega_{\beta}\right)\right],(\alpha, \beta=1, \cdots, \rho)$, are nonsingular. Then there exists a unique generalized Green's function $g_{L}(, ; \theta, \Omega)$ for (5.2) satisfying the conditions

$$
\begin{array}{ll}
\int_{a}^{b} g_{L}(x, t ; \Theta, \Omega) \omega_{\alpha}(t) d t=0, & (\alpha=1, \cdots, \rho),  \tag{5.12}\\
\int_{a}^{b} \bar{\theta}_{2}(x) g_{L}(x, t ; \Theta, \Omega) d x=0, & (i=1, \cdots, r) .
\end{array}
$$

Without any loss of generality we can assume that $\left[\left(u_{i}, \theta_{j}\right)\right]=E_{r}$ and $\left[\left(v_{\alpha}, \omega_{\beta}\right)\right]=E_{\rho}$. Let $g_{0}$ be the generalized Green's function for (5.2) described in Theorem 5.2. We now determine functions $\psi_{1}, \cdots, \psi_{r}$ and functions $\varphi_{1}, \cdots, \varphi_{\rho}$ such that the generalized Green's function given by (5.10) satisfies conditions (5.12). Such a generalized Green's function $g$ will satisfy the conditions (5.12) if and only if the functions $\psi_{i},(i=1, \cdots, r)$, and $\rho_{\alpha},(\alpha=1, \cdots, \rho)$, satisfy the equations

$$
\begin{align*}
\psi_{i}(x)+\int_{a}^{b} \sum_{\beta=1}^{\rho} \bar{\theta}_{i}(s) \varphi_{\beta}(s) \bar{v}_{\beta}(x) d s+\int_{a}^{b} \bar{\theta}_{i}(s) g_{0}(s, x) d s & =0, \\
& (i=1, \cdots, r),  \tag{5.13}\\
\varphi_{\alpha}(x)+\int_{a}^{b} \sum_{j=1}^{r} u_{j}(x) \psi_{\nu}(s) \omega_{\alpha}(s) d s+\int_{a}^{b} g_{0}(x, s) \omega_{\alpha}(s) d s & =0, \\
(\alpha & =1, \cdots, \rho) .
\end{align*}
$$

A particular set of solutions for equations (5.13) is

$$
\begin{align*}
\varphi_{\alpha}(x)= & -\int_{a}^{b} g_{0}(x, s) \omega_{\alpha}(s) d s, & (\alpha=1, \cdots, \rho), \\
\psi_{i}(x)= & \int_{a}^{b} \int_{a}^{b} \sum_{\beta=1}^{p} \bar{\theta}_{i}(t) g_{0}(t, s) \omega_{\beta}(s) \bar{v}_{\beta}(x) d s d t &  \tag{5.14}\\
& -\int_{a}^{b} \bar{\theta}_{i}(t) g_{0}(t, x) d t, & (i=1, \cdots, r) .
\end{align*}
$$

Moreover, if $\psi_{i}$ and $\varphi_{\alpha},(i=1, \cdots, r ; \alpha=1, \cdots, \rho)$, is any collection of solutions of (5.13), then after substituting the value of $\psi_{i}(x)$ given by the first equation into the second equation of (5.13) it can be shown by straightforward computation that the value of

$$
\sum_{i=1}^{r} u_{i}(x) \psi_{i}(t)+\sum_{\alpha=1}^{p} \varphi_{n}(x) \bar{v}_{\alpha}(t)
$$

is independent of the particular $\psi_{i}$ and $\varphi_{\alpha}$. Hence there is a unique generalized Green's function for (5.2) satisfying (5.12).

The conditions of Theorem 5.5 are clearly satisfied by the sets $\theta_{i}=u_{i},(i=1, \cdots, r)$, and $\omega_{\alpha}=v_{\alpha},(\alpha=1, \cdots, \rho)$; in particular, for linear homogeneous differential operators whose coefficients satisfy
suitable differentiability conditions, the treatment of Greub and Rheinboldt [4] is limited to this specification.

It is to be remarked that, in view of the definition of $g_{0}$, if $\psi_{i}$ and $\varphi_{\alpha_{2}}(i=1, \cdots, r ; \alpha=1, \cdots, \rho)$, is any collection of solutions of (5.13), then $\varphi_{\alpha} \in \tilde{\mathfrak{A}}_{n},(a=1, \cdots, \rho)$, and $\bar{\psi}_{i} \in \tilde{\mathfrak{N}}_{m},(i=1, \cdots, r)$.

Correspondingly, there exists a unique generalized Green's function $g_{L^{\sharp}}(, ; \Omega, \Theta)$ for the system (5.3) which satisfies the conditions

$$
\begin{align*}
\int_{a}^{b} \bar{\omega}_{\alpha}(x) g_{L^{\star}}(x, t ; \Omega, \Theta) d x=0, & (\alpha=1, \cdots, \rho), \\
\int_{a}^{b} g_{L^{\sharp}}(x, t ; \Omega, \Theta) \theta_{i}(t) d t=0, & (i=1, \cdots, r) . \tag{5.15}
\end{align*}
$$

For brevity, denote by $b_{a}$ and $b_{\theta}$ the functions defined on $\square$ by the formulas

$$
b_{\Omega}(x, t)=\sum_{j=1}^{\rho} \omega_{j}(x) \bar{v}_{j}(t), \quad b_{\theta}(x, t)=\sum_{i=1}^{r} \theta_{i}(x) \bar{u}_{i}(t) .
$$

Theorem 5.6. If $g_{L}(, ; \Theta, \Omega)$ is the unique generalized Green's function satisfying (5.12), then the following conditions (5.16)-(5.20) are satisfied:
(5.16) $g_{L}^{\langle j, 0\rangle}(, ; \Theta, \Omega),(j=0, \cdots, m+n-2)$, exists and is continuous on $\square$ while $g_{L}^{\langle m+n-1,0\rangle}(x, t ; \Theta, \Omega)$ and $\partial / \partial x g_{L}^{\langle m+n-1,0\rangle}(x, t ; \theta, \Omega)$ exist on the individual domains $a \leqq t<x, a<x<b$ and $a \leqq x<b, x<t \leqq b$;
(5.17) if $t \in[a, b]$, then the function whose value at $x \neq t$ is $g_{L}^{\langle m+n-1,0\rangle}(x, t ; \Theta, \Omega)$ has a right and a left limit at $t$, denoted by $g_{L}^{\langle m+n-1,0\rangle}\left(t^{+}, t ; \Theta, \Omega\right)$ and $g_{L}^{\langle m+n-1,0\rangle}\left(t^{-}, t ; \Theta, \Omega\right)$, respectively, and

$$
g_{L}^{\langle m+n-1,0\rangle}\left(t^{-}, t ; \Theta, \Omega\right)-g_{L}^{\langle m+n-1,0\rangle}\left(t^{+}, t ; \Theta, \Omega\right)=1 ;
$$

(5.18) if $t \in[a, b]$, then $L\left[g_{L}(, t ; \Theta, \Omega)\right]=b \Omega(, t)$ on $[a, t)$ and $(t, b]$;
(5.19) if $t \in(a, b)$, then the function whose value at $x$ is $g_{L}(x, t ; \theta, \Omega)$ satisfies the boundary conditions which characterize the set $\mathscr{D}$;

$$
\begin{equation*}
\int_{a}^{b} \bar{\theta}_{i}(x) g_{L}(x, t ; \Theta, \Omega) d x=0, \quad(i=1, \cdots, r ; t \in[a, b]) \tag{5.20}
\end{equation*}
$$

Conditions (5.16)-(5.18) may be verified directly using the properties of $g_{0}$ and the remark following the proof of Theorem 5.5. Condition (5.20) is merely one of the conditions in (5.12). If $\mathscr{D}=\widetilde{\mathfrak{A}}_{n}$, then (5.19) is trivially satisfied. Otherwise, let $w$ be any integrable function, and define $f$ by

$$
f(x)=w(x)-\sum_{\alpha=1}^{\rho} \omega_{\alpha}(x)\left(w, v_{\alpha}\right)=w(x)-\int_{\alpha}^{b} b_{2}(x, t) w(t) d t
$$

In view of the assumption that $\left[\left(v_{\alpha}, \omega_{\beta}\right)\right]=E_{\rho}$, it follows that $\left(f, v_{\alpha}\right)=0$, ( $\alpha=1, \cdots, \rho$ ), and therefore the function $u$ defined by

$$
u(x)=\int_{a}^{b} g_{L}(x, t ; \Theta, \Omega) f(t) d t
$$

is a solution of (5.1). But it follows from (5.12) that

$$
\int_{a}^{b} g_{L}(x, t, \Theta, \Omega) f(t) d t=\int_{a}^{b} g_{L}(x, t ; \Theta, \Omega) w(t) d t
$$

Therefore,

$$
\begin{aligned}
0 & =M \hat{u}(a)+N \hat{u}(b) \\
& =\int_{a}^{b}\left(M \widehat{g}_{L}(a, t ; \Theta, \Omega)+N \widehat{g}_{L}(b, t ; \Theta, \Omega)\right) w(t) d t,
\end{aligned}
$$

from which (5.19) follows in view of the arbitrariness of the function $w$.
Corollary. If $w \in \mathfrak{A}_{0}$ and $y$ is defined by

$$
y(x)=\int_{a}^{b} g_{L}(x, t ; \Theta, \Omega) w(t) d t
$$

then

$$
\begin{aligned}
& L[y]=w-\int_{a}^{b} b_{\Omega}(, t) w(t) d t \\
& y \in \mathscr{D}, \quad\left(y, \theta_{i}\right)=0,
\end{aligned} \quad(i=1, \cdots, r) .
$$

It should be noted that the unique generalized Green's function $g_{L^{\star}}(, ; \Omega, \Theta)$ for (5.3) which satisfies (5.15) also satisfies conditions analogous to (5.16)-(5.20).

Theorem 5.7. If $x, t \in[a, b]$, then $g_{L^{\star}}(x, t ; \Omega, \Theta)=\bar{g}_{L}(t, x ; \Theta, \Omega)$.
Let $w$ and $h$ be arbitrary integrable functions and define $y$ and $z$ by

$$
\begin{aligned}
& y(x)=\int_{a}^{b} g_{L}(x, t ; \theta, \Omega) w(t) d t \\
& z(x)=\int_{a}^{b} g_{L^{\star}}(x, t ; \Omega, \Theta) h(t) d t
\end{aligned}
$$

respectively. Then it follows from the corollary to Theorem 5.6 and its analogue that $y \in \mathscr{D}, z \in \mathscr{D}^{*}$, and therefore

$$
(L[y], z)-\left(y, L^{\grave{幺}}[z]\right)=0 .
$$

But it also follows from the corollary to Theorem 5.6 that $L[y]=$
$w-\int_{a}^{b} b_{\Omega}(, t) w(t) d t, L^{\text {৯ }}[z]=h-\int_{a}^{b} b_{\theta}(, t) h(t) d t$, and therefore in view of (5.12), (5.15), and the definition of $b_{\Omega}$ and $b_{\theta}$, we have

$$
\int_{a}^{b} \int_{a}^{b} \bar{h}(x)\left[\bar{g}_{L^{\sharp}}(t, x ; \Omega, \Theta)-g_{L}(x, t ; \Theta, \Omega] w(t) d t d x=0,\right.
$$

from which the theorem follows since $w$ and $h$ are arbitrary integrable functions.

Corollary I. The function $g_{L}(, ; \Theta, \Omega)$ is characterized by conditions (5.16)-(5.20), and the function $g_{L^{\sharp}}(, ; \theta, \Omega)$ is characterized by analogous conditions.

As a consequence of Theorems 5.4 and 5.7 one has the following result:

Corollary II. If $g$ is a generalized Green's function for (5.2), then the function $h$ defined by $h(x, t)=\bar{g}(t, x)$ is a generalized Green's function for (5.3).
6. A canonical form for boundary conditions. Let $\left[f_{i j}\right]$ and $\left[g_{i j}\right],(i=0, \cdots, m \geqq 1 ; j=0, \cdots, n \geqq 1)$, be $(m+1) \times(n+1)$ integrable matrix functions. Suppose that the matrix function $\left[f_{i j}\right]$, $(i=0, \cdots, m ; j=0, \cdots, n)$, satisfies hypothesis $(H)$, and $g_{m j}(x) \equiv$ $g_{i n}(x) \equiv 0,(i=0, \cdots, m ; j=0, \cdots, n)$.

For a complex number $\lambda$ let $p_{i j}(; \lambda)$ be the function defined on $[a, b]$ by

$$
p_{i j}(x ; \lambda)=f_{i j}(x)+\lambda g_{i j}(x), \quad(i=0, \cdots, m ; j=0, \cdots, n)
$$

It follows that for each number $\lambda$ hypothesis (H) holds for the matrix function $\left[p_{i j}(; \lambda)\right]$. For suitable $y$ in $\mathfrak{U}_{n}$ let $\widetilde{y}_{1}(; \lambda), \cdots, \widetilde{y}_{m}(; \lambda)$ be defined on $[a, b]$ as follows:

$$
\begin{align*}
& \widetilde{y}_{m}(x ; \lambda)=\sum_{j=0}^{n} p_{m j}(x ; \lambda) y^{(j)}(x)=\sum_{j=0}^{n} f_{m j}(x) y^{(j)}(x) \\
& \text { if } \widetilde{y}_{i+1}(; \lambda) \in \mathfrak{H}_{1}, \text { then } \widetilde{y}_{i}(x ; \lambda)=\sum_{j=0}^{n} p_{i j}(x ; \lambda) y^{(j)}(x)-\widetilde{y}_{i+1}^{\prime}(x ; \lambda),  \tag{6.1}\\
& \quad(i=m-1, \cdots, 1) .
\end{align*}
$$

The class of functions $y$ in $\mathfrak{A}_{n}$ for which $\widetilde{y}_{1}(, \lambda), \cdots, \widetilde{y}_{m}(; \lambda)$ are a.c. will be denoted by $\tilde{\mathfrak{A}}_{n}(\lambda)$, and $L[; \lambda]$ will be the operator with domain $\widetilde{\mathfrak{A}}_{n}(\lambda)$, and defined by

$$
\begin{equation*}
L[y ; \lambda]=\sum_{j=0}^{n} p_{0 j}(; \lambda) y^{(j)}-\widetilde{y}_{1}^{\prime}(; \lambda) . \tag{6.2}
\end{equation*}
$$

The vector function $\left(\widetilde{y}_{i}(; \lambda)\right)$, $(i=1, \cdots, m)$, will be represented by $\widetilde{y}(; \lambda)$, and $\widehat{y}(; \lambda)$ will signify the $(n+m)$-vector function ( $y, \cdots$, $\left.y^{(n-1)}, \widetilde{y}_{1}(; \lambda), \cdots, \widetilde{y}_{m}(; \lambda)\right)$. For a complex number $\nu$ let $p_{j i}^{\prime}(; \nu)$ be the function on $[a, b]$ defined by

$$
p_{j \nu}^{\stackrel{\rightharpoonup}{j}}(x ; \nu)=\bar{f}_{i j}(x)+\nu \bar{\nu}_{i j}(x), \quad(i=0, \cdots, m ; j=0, \cdots, n) .
$$

For suitable $z$ in $\mathfrak{U}_{m}$ define $\widetilde{z}_{1}(; \nu), \cdots, \widetilde{z}_{n}(; \nu)$ by

$$
\begin{align*}
& \widetilde{z}_{n}(x ; \nu)=\sum_{\imath=0}^{m} p_{n i}^{\ngtr}(x ; \nu) z^{(i)}(x)=\sum_{i=1}^{m} \bar{f}_{i n}(x) z^{(i)}(x) ; \\
& \text { if } \widetilde{z}_{j+1}(; \nu) \in \mathfrak{H}_{1}, \text { then } \widetilde{z}_{j}(x ; \nu)=\sum_{i=1}^{m} p_{j i}^{\stackrel{\rightharpoonup}{\grave{i}}}(x ; \nu) z^{(i)}(x)-\widetilde{z}_{j+1}^{\prime}(x ; \nu) ;  \tag{6.3}\\
& \\
& \quad(j=n-1, \cdots, 1) .
\end{align*}
$$

The class of functions $z$ in $\mathfrak{A}_{m}$ for which $\widetilde{z}_{1}(; \nu), \cdots, \widetilde{z}_{n}(; \nu)$ are a.c. will be denoted by $\widetilde{\mathfrak{A}}_{m}(\nu)$ and $L^{\star}[; \nu]$ will be operator with domain $\widetilde{\mathfrak{A}}_{m}(\nu)$, and defined by

$$
\begin{equation*}
L^{\star}[z ; \nu]=\sum_{i=1}^{m} p_{0 \dot{\sim}}^{\dot{\sim}}(; \nu) z^{(i)}-\widetilde{z}_{1}^{\prime}(; \nu) . \tag{6.4}
\end{equation*}
$$

The vector function $\left(\widetilde{z}_{j}(; \nu)\right),(j=1, \cdots, n)$, will be represented by $\widetilde{z}(; \nu)$, and $\breve{z}(; \nu)$ will denote the vector function $\left(z, \cdots, z^{(m-1)}, \widetilde{z}_{1}(, \nu), \cdots\right.$, $\left.\widetilde{z}_{n}(; \nu)\right)$. Let $A_{10}, A_{11}, A_{20}$, and $A_{21}$ be $k \times n$ matrices, and let $B_{1}$ and $B_{2}$ be $k \times m$ matrices, $(1 \leqq k \leqq 2 m+2 n-1)$, such that for each number $\lambda$ the $k \times 2(m+n)$ matrix

$$
\left[A_{1}(\lambda)-B_{1} A_{2}(\lambda) B_{2}\right]
$$

has rank $k$, where $A_{1}(\lambda)=A_{10}+\lambda A_{11}$ and $A_{2}(\lambda)=A_{20}+\lambda A_{21}$. Let $\mathscr{D}(\lambda)$ be the collection of functions $y$ in $\tilde{\mathfrak{A}}_{n}(\lambda)$ for which

$$
\begin{equation*}
A_{1}(\lambda) \hat{y}(a)-B_{1} \widetilde{y}(a ; \lambda)+A_{2}(\lambda) \hat{y}(b)+B_{2} \widetilde{y}(b ; \lambda)=0 . \tag{6.5}
\end{equation*}
$$

This section is concerned with the particular Euler type quasi-differential system

$$
\begin{equation*}
L[y ; \lambda]=0, \quad y \in \mathscr{D}(\lambda) \tag{6.6}
\end{equation*}
$$

It follows from Theorem 3.2 that the system adjoint to (6.6) is

$$
\begin{equation*}
L^{\star}[z ; \bar{\lambda}]=0, \quad z \in \mathscr{D}^{*}(\bar{\lambda}), \tag{6.7}
\end{equation*}
$$

where $\mathscr{D} *(\bar{\lambda}) \subset \tilde{\mathfrak{A}}_{m}(\bar{\lambda})$. The following assumption is made about $\mathscr{D}^{*}(\bar{\lambda})$ :
Hypothesis $\left(\mathrm{H}_{1}\right)$. There exist $(2 m+2 n-k) \times m$ matrices $A_{3}(\nu)=$ $A_{30}+\nu A_{31}, A_{4}(\nu)=A_{40}+\nu A_{41}$ and $(2 m+2 n-k) \times n$ matrices $B_{3}, B_{4}$ such that for arbitrary $\lambda$ the set $\mathscr{D}^{*}(\bar{\lambda})$ is the collection of function $z$ in $\widetilde{\mathfrak{A}}_{m}(\bar{\lambda})$ for which

$$
\begin{equation*}
A_{3}(\bar{\lambda}) \check{z}(a)-B_{3} \tilde{z}(a ; \bar{\lambda})+A_{4}(\bar{\lambda}) \check{z}(b)+B_{4} \tilde{z}(b ; \bar{\lambda})=0 . \tag{6.8}
\end{equation*}
$$

It shoud be noted that the assumption used by Zimmerberg to obtain Theorem 2.1 of [13] does not imply that hypothesis $\left(\mathrm{H}_{1}\right)$ holds. For if $m=n=1$ and $k=2 n$, then let the matrices $A_{10}, A_{11}, B_{1}, A_{20}$, $A_{21}, B_{2}$ be defined as

$$
\begin{array}{lll}
A_{10}^{*}=\left[\begin{array}{ll}
1 & 1
\end{array}\right], & \mathrm{A}_{11}^{*}=\left[\begin{array}{ll}
0 & 1
\end{array}\right], & B_{1}^{*}=\left[\begin{array}{ll}
2 & 1
\end{array}\right] \\
A_{20}^{*}=\left[\begin{array}{lll}
1 & 0
\end{array}\right], & A_{21}^{*}=\left[\begin{array}{lll}
0 & 1
\end{array}\right], & B_{2}^{*}=\left[\begin{array}{lll}
0 & 1
\end{array}\right]
\end{array}
$$

Then the hypothesis of Theorem 2.1 of [13] is satisfied, but hypothesis $\left(H_{1}\right)$ does not hold.

If hypothesis $\left(H_{1}\right)$ holds then for each complex number $\nu$ the $(2 m+2 n-k) \times 2(m+n)$ matrix

$$
\begin{equation*}
\left[A_{3}(\nu) B_{3} A_{4}(\nu) B_{4}\right] \tag{6.9}
\end{equation*}
$$

has rank $2 m+2 n-k$. Moreover, by a proof quite analogous to that used by Reid to obtain (11.11') of [6] one may establish the following result.

Lemma 6.1. If hypothesis $\left(H_{1}\right)$ holds, then $\mathscr{D}(\lambda)$ is the collection of functions $y$ in $\tilde{\mathfrak{A}}_{n}(\lambda)$ for which there is a $(2 m+2 n-k)$-vector $e_{0}$ such that

$$
\begin{array}{ll}
\widehat{y}(a)=B_{3}^{*} e_{0}, & \widetilde{y}(a ; \lambda)=A_{3}^{*}(\bar{\lambda}) e_{0} \\
\widehat{y}(b)=B_{4}^{*} e_{0}, & \widetilde{y}(b ; \lambda)=-A_{4}^{*}(\bar{\lambda}) e_{0} \tag{6.10}
\end{array}
$$

and $\mathscr{D}^{*}(\bar{\lambda})$ is the collection of functions $z$ in $\widetilde{\mathfrak{A}}_{m}(\bar{\lambda})$ for which there is a k-vector $e_{1}$ such that

$$
\begin{array}{ll}
\check{z}(\alpha)=B_{1}^{*} e_{1}, & \widetilde{z}(\alpha ; \bar{\lambda})=A_{1}^{*}(\lambda) e_{1}, \\
\check{z}(b)=B_{2}^{*} e_{1}, & \widetilde{z}(b ; \bar{\lambda})=-A_{2}^{*}(\lambda) e_{1}, \tag{6.11}
\end{array}
$$

where $A_{i}^{*}(\nu)=\left(A_{i}(\nu)\right)^{*},(i=1,2,3,4)$.
Now let $K_{10}=A_{10} B_{3}^{*}+A_{20} B_{4}^{*}, K_{11}=A_{11} B_{3}^{*}+A_{21} B_{4}^{*}, K_{1}(\lambda)=K_{10}+$ $\lambda K_{11}, K_{20}=A_{30} B_{1}^{*}+A_{40} B_{2}^{*}, K_{21}=A_{31} B_{1}^{*}+A_{41} B_{2}^{*}$, and $K_{2}(\lambda)=K_{20}+\lambda K_{21}$. Then the next result follows from Lemma 6.1 and Lemma 3.1.

Lemma 6.2. If hypothesis $\left(H_{1}\right)$ holds, then $K_{2}^{*}(\bar{\lambda})=K_{1}(\lambda)$.
Lemma 6.3. Suppose that hypothesis $\left(H_{1}\right)$ holds, the $k \times 2$ m matrix [ $B_{1} B_{2}$ ] has rank $k-p$, and the $(2 m+2 n-k) \times 2 n$ matrix $\left[B_{3} B_{4}\right]$ has rank $2 m+2 n-k-q$. Then there exist $p \times n$ matrices $\psi_{1}, \psi_{2}$ and $q \times m$ matrices $\psi_{3}, \psi_{4}$ such that the $p \times 2 n$ matrix $\left[\psi_{1} \psi_{2}\right]$ has rank $p$, the $q \times 2 m$ matrix $\left[\psi_{3} \psi_{4}\right.$ ] has rank $q$, and

$$
\begin{array}{ll}
\psi_{1} \hat{y}(a)+\psi_{2} \hat{y}(b)=0, & \text { for } y \in \mathscr{D}(\lambda), \\
\psi_{3} \check{z}(a)+\psi_{4} \check{z}(b)=0, & \text { for } z \in \mathscr{D}^{*}(\bar{\lambda}) . \tag{6.13}
\end{array}
$$

Suppose that $R$ is a $p \times k$ matrix of rank $p$ such that $R\left[B_{1} B_{2}\right]=0$, and define $\psi_{1}$ and $\psi_{2}$ as $\psi_{1}=R A_{10}, \psi_{2}=R A_{20}$. In view of Lemma 6.2 and the fact that for arbitrary complex $\lambda$ the $k \times 2(m+n)$ matrix $\left[A_{1}(\lambda) B_{1} A_{2}(\lambda) B_{2}\right]$ has rank $k$ it follows that there exists a $p \times p$ matrix $V$ such that

$$
\left[R A_{1}(\lambda) R A_{2}(\lambda)\right]=\left(E_{p}+\lambda V\right) R\left[A_{10} A_{20}\right]
$$

Hence $E_{p}+\lambda V$ is nonsingular and the equation (6.12) is equivalent to

$$
R A_{1}(\lambda) \widehat{y}(a)+R A_{2}(\lambda) \widehat{y}(b)=0
$$

If $R_{0}$ is a $q \times(2 m+2 n-k)$ matrix of rank $q$ such that $R_{0}\left[B_{3} B_{4}\right]=0$, and $\psi_{3}, \psi_{4}$ are defined as $\psi_{3}=R_{0} A_{30}, \psi_{4}=R_{0} A_{40}$, then equation (6.13) may be verified in a similar fashion. The conclusion concerning the ranks of $\left[\psi_{1} \psi_{2}\right]$ and $\left[\psi_{3} \psi_{4}\right]$ is clear.

From Lemma 6.2 it then follows that $\left[B_{1} B_{2}\right]\left[\psi_{3} \psi_{4}\right]^{*}=0$ and $\left[B_{3} B_{4}\right]\left[\psi_{1} \psi_{2}\right]^{*}=0$, so that $q \leqq 2 m-(k-p)$ and $p \leqq 2 n-[2 m+2 n-$ $k-q]=k+q-2 m$, from which one has the following result.

Lemma 6.4. If hypothesis $\left(H_{1}\right)$ holds, then the columns of $\left[\psi_{3} \psi_{4}\right]^{*}$ form a basis for the null space of $\left[B_{1} B_{2}\right]$ and the columns of $\left[\psi_{1} \psi_{2}\right]^{*}$ form a basis for the null space of $\left[B_{3} B_{4}\right]$.

The following theorem gives a simultaneous canonical representation of the boundary conditions for (6.6) and (6.7) in terms of parameter matrices $\psi_{i}, Q_{i}, G_{i},(i=1,2,3,4)$, and is the central result of this section.

Theorem 6.1. Suppose that hypothesis $\left(H_{1}\right)$ holds. Then there exist $m \times n$ matrices $Q_{i}$ and $G_{i},(i=1,2,3,4)$, such that $y \in \mathscr{D}(\lambda)$ if and only if there exists a $q$-vector $\eta_{1}$ such that

$$
\begin{align*}
\psi_{1} \hat{y}(a)+\psi_{2} \hat{y}(b) & =0, \\
\left(Q_{1}-\lambda G_{1}\right) \hat{y}(a)+\left(Q_{2}-\lambda G_{2}\right) \hat{y}(b)+\psi_{3}^{*} \eta_{1}-\widetilde{y}(a ; \lambda) & =0,  \tag{6.14}\\
\left(Q_{3}-\lambda G_{3}\right) \hat{y}(a)+\left(Q_{4}-\lambda G_{4}\right) \hat{y}(b)+\psi_{4}^{*} \eta_{1}+\widetilde{y}(b ; \lambda) & =0 .
\end{align*}
$$

Moreover, $z \in \mathscr{D}^{*}(\bar{\lambda})$ if and only if there exists a p-vector $\eta_{2}$ such that

$$
\begin{align*}
\dot{\psi}_{3} \check{z}(a)+\psi_{4} \check{z}(b) & =0, \\
\left(Q_{1}^{*}-\bar{\lambda} G_{1}^{*}\right) \check{z}(a)+\left(Q_{3}^{*}-\bar{\lambda} G_{3}^{*}\right) \check{z}(b)+\dot{\psi}_{1}^{*} \eta_{2}-\widetilde{z}(a ; \bar{\lambda}) & =0,  \tag{6.15}\\
\left(Q_{2}^{*}-\bar{\lambda} G_{2}^{*}\right) \check{z}(a)+\left(Q_{4}^{*}-\bar{\lambda} G_{4}^{*}\right) \check{z}(b)+\psi_{2}^{*} \eta_{2}+\widetilde{z}(b ; \bar{\lambda}) & =0 .
\end{align*}
$$

Suppose that the matrices $K_{10}$ and $K_{11}$ have ranks $q_{0}$ and $q_{1}$, respectively. Let $D_{10}$ and $D_{11}$ be $(2 m+2 n-k) \times\left(2 m+2 n-k-q_{0}\right)$ and $(2 m+2 n-k) \times\left(2 m+2 n-k-q_{1}\right)$ matrices, respectively, whose individual column vectors form orthonormal bases for the null spaces of $K_{10}$ and $K_{11}$, that is, $K_{10} D_{10}=0$ and $K_{11} D_{11}=0$. As $K_{20}=K_{10}^{*}$ and $K_{21}=K_{11}^{*}$ by Lemma 6.2, there exist matrices $D_{20}$ and $D_{21}$ of respective orders $k \times\left(k-q_{0}\right)$ and $k \times\left(k-q_{1}\right)$ whose individual column vectors form orthonormal bases for the null spaces of $K_{20}$ and $K_{21}$. Then

$$
\left[\begin{array}{cc}
K_{10} & D_{20}  \tag{6.16}\\
D_{20}^{*} & 0
\end{array}\right], \quad\left[\begin{array}{cc}
K_{11} & D_{21} \\
D_{11}^{*} & 0
\end{array}\right], \quad\left[\begin{array}{cc}
K_{20} & D_{10} \\
D_{20}^{*} & 0
\end{array}\right], \quad\left[\begin{array}{cc}
K_{21} & D_{11} \\
D_{21}^{*} & 0
\end{array}\right]
$$

are nonsingular and have inverses of the form

$$
\left[\begin{array}{cc}
H_{10} & D_{10}  \tag{6.17}\\
D_{20}^{*} & 0
\end{array}\right], \quad\left[\begin{array}{cc}
H_{11} & D_{11} \\
D_{21}^{*} & 0
\end{array}\right], \quad\left[\begin{array}{cc}
H_{10}^{*} & D_{20} \\
D_{10}^{*} & 0
\end{array}\right], \quad\left[\begin{array}{cc}
H_{11}^{*} & D_{21} \\
D_{11}^{*} & 0
\end{array}\right],
$$

respectively. The matrices $H_{10}, H_{11}, H_{10}^{*}$, and $H_{11}^{*}$ are generalized reciprocals of the respective matrices $K_{10}, K_{11}, K_{20}=K_{10}^{*}$, and $K_{21}=K_{11}^{*}$. Let $Q_{i}$ and $G_{i}, \quad(i=1,2,3,4)$, be defined as $Q_{1}=A_{30}^{*} H_{10} A_{10}, Q_{2}=$ $A_{30}^{*} H_{10} A_{20}, Q_{3}=A_{40}^{*} H_{10} A_{10}, Q_{4}=A_{40}^{*} H_{10} A_{20}, G_{1}=-A_{31}^{*} H_{11} A_{11}, G_{2}=-A_{31}^{*} H_{11} A_{21}$, $G_{3}=-A_{41}^{*} H_{11} A_{11}$, and $G_{4}=-A_{41}^{*} H_{11} A_{21}$.

Now if $y \in \mathscr{D}(\lambda)$ then in view of Lemma 6.3 we need only verify the last two equations of (6.14). Suppose that $e_{0}$ is determined by (6.10). Then it follows from (6.10) and the fact that the matrices (6.17) are the inverses of the matrices (6.16) that

$$
\begin{align*}
& e_{0}=H_{10} A_{10} \hat{y}(a)+H_{10} A_{20} \hat{y}(b)+D_{10} D_{10}^{*} e_{0}, \\
& e_{0}=H_{11} A_{11} \hat{y}(a)+H_{11} A_{21} \hat{y}(b)+D_{11} D_{11}^{*} e_{0} \tag{6.18}
\end{align*}
$$

Now it follows from (6.10) and (6.18) that

$$
\begin{align*}
\left(Q_{1}-\lambda G_{1}\right) \hat{y}(a) & +\left(Q_{2}-\lambda G_{2}\right) \hat{y}(b)+\left(A_{30}^{*} D_{10} D_{10}^{*}+\lambda A_{31}^{*} D_{11} D_{11}^{*}\right) e_{0} \\
& -\widetilde{y}(a ; \lambda)=0 \\
\left(Q_{3}-\lambda G_{3}\right) \hat{y}(a) & +\left(Q_{4}-\lambda G_{4}\right) \hat{y}(b)+\left(A_{40}^{*} D_{10} D_{10}^{*}+\lambda A_{11}^{*} D_{11} D_{11}^{*}\right) e_{0}  \tag{6.19}\\
& +\widetilde{y}(b ; \lambda)=0
\end{align*}
$$

But $B_{1}\left(A_{30}^{*} D_{10} D_{10}^{*}+\lambda A_{31}^{*} D_{11} D_{11}^{*}\right)+B_{2}\left(A_{40}^{*} D_{10} D_{10}^{*}+\lambda A_{41}^{*} D_{11} D_{11}^{*}\right)=K_{20}^{*} D_{10} D_{10}^{*}+$ $\lambda K_{21}^{*} D_{11} D_{11}^{*}=0$, and consequently the two equations of (6.19) may be written in the form of the last two equations of (6.14) involving the parameter vector $\eta_{1}$.

On the other hand, suppose that $y \in \widetilde{\mathfrak{A}}_{n}(\lambda)$ and (6.14) holds. Now the first equation of (6.14) implies that there is a $(2 m+2 n-k)$-vector $e_{0}$ such that $\hat{y}(a)=B_{3}^{*} e_{0}$ and $\hat{y}(b)=B_{4}^{*} e_{0}$. Hence it follows from (6.16) and (6.17) that (6.18) holds for this value of $e_{0}$. Solving the equations
(6.18) for $H_{10} A_{10} \hat{y}(a)+H_{10} A_{20} \hat{y}(b)$ and $H_{11} A_{11} \hat{y}(a)+H_{11} A_{21} \hat{y}(b)$, multiplying the first equation on the left by $A_{30}^{*}$ and $A_{40}^{*}$, and the second equation on the left by $\lambda A_{31}^{*}$ and $\lambda A_{41}^{*}$, respectively, and adding it can be shown that the last two equations of (6.14) may be written as

$$
\begin{align*}
& A_{30}^{*}\left(e_{0}-D_{10} D_{10}^{*} e_{0}\right)+\lambda A_{31}^{*}\left(e_{0}-D_{11} D_{11}^{*} e_{0}\right)+\psi_{3}^{*} \eta_{1}-\widetilde{y}(a ; \lambda)=0, \\
& A_{40}^{*}\left(e_{0}-D_{10} D_{10}^{*} e_{0}\right)+\lambda A_{11}^{*}\left(e_{0}-D_{11} D_{11}^{*} e_{0}\right)+\psi_{4}^{*} \eta_{1}+\widetilde{y}(b ; \lambda)=0 . \tag{6.20}
\end{align*}
$$

In view of Lemma 6.2, the definition of the matrices $D_{10}, D_{11}$, and the choice of the vector $e_{0}$, one sees after multiplying the first equation of (6.20) by $B_{1}$, the second equation by $B_{2}$, and adding the two equations, that $y$ satisfies the boundary conditions of (6.6). The conclusion concerning $D^{*}(\bar{\lambda})$ may be established in a similar manner.

The next theorem is an application of Theorem 6.1, where it is to be noticed that if $m=n$ and $\left[f_{i j}(x)\right],\left[g_{i j}(x)\right]$ are Hermitian, then $\widetilde{\mathfrak{N}}_{n}(\lambda)=\widetilde{\mathfrak{A}}_{n}(\lambda)$; in particular, if $z \in \widetilde{\mathfrak{A}}_{n}(\lambda)$, then $\bar{z}(; \lambda)=\bar{z}(; \lambda)$.

Theorem 6.2. Suppose that $m=n,\left[f_{i j}(x)\right]$ and $\left[g_{i j}(x)\right]$ are Hermitian on $[a, b], k=2 n$, and $\mathscr{D}^{*}(\bar{\lambda})=\mathscr{D}(\bar{\lambda})$. Then the system (6.6) is equivalent to the Euler-Lagrange equations and transversality conditions for minimizing the functional

$$
\hat{y}^{*}(\alpha)\left[Q_{1} \hat{y}(\alpha)+Q_{2} \hat{y}(b)\right]+\hat{y}^{*}(b)\left[Q_{2}^{*} \hat{y}(\alpha)+Q_{4} \hat{y}(b)\right]+\int_{a}^{b} \sum_{\alpha=0}^{n} \bar{y}^{(\alpha)} f_{\alpha \beta} y^{(\beta)}
$$

subject to the restraints

$$
\begin{aligned}
& \psi_{1} \hat{y}(a)+\psi_{2} \hat{y}(b)=0, \\
& \hat{y}^{*}(\alpha)\left[G_{1} \hat{y}(a)+G_{2} \hat{y}(b)\right]+\hat{y}^{*}(b)\left[G_{2}^{*} \hat{y}(a)+G_{4} \hat{y}(b)\right]+\int_{a \alpha, \beta=0}^{b} \sum^{n-1} \bar{y}^{(\alpha)} g_{\alpha \beta} y^{(\beta)} \\
&=\mathrm{const} .
\end{aligned}
$$

If $m=n$, the problem is restricted to the field of real numbers, $g_{i j}(x) \equiv f_{i j}(x) \equiv 0$ for $i \neq j$, and if $f_{i i}, g_{i i} \in \mathfrak{C}_{i},(i, j=0, \cdots, n)$, then the results of this section are the same as obtained by Zimmerberg [12], provided that the formula (2.4) of that paper is corrected by replacing $f_{i}, f_{i+1}, \cdots, f_{n-1}$ by $f_{i}-\lambda g_{i}, f_{i+1}-\lambda g_{i+1}, \cdots, f_{n-1}-\lambda g_{n-1}$, respectively. If, moreover, $g_{i i}(x) \equiv 0$ for $i \geqq 1$, then these are the same results as obtained by Reid [6, Section 11].
7. An application. In this section the results of Section 6 and a theorem of Reid [7] will be used to show that the boundary conditions for'a rather large class of linear $\nu$ th order differential operators may be written in the form given by Theorem 6.1.

Reid [7] has considered $\nu$ th order linear differential operators $L$ of the form

$$
\begin{equation*}
L[y]=\sum_{j=0}^{\nu} q_{j}(x) y^{(j)}, \quad \nu \geqq 1, \tag{7.1}
\end{equation*}
$$

with integrable coefficients. Functions $\Lambda_{i}(y ; p),(i=0,1,2, \cdots)$, were defined as

$$
\begin{aligned}
\Lambda_{0}(y ; p) & \equiv p(x) y, \quad \Lambda_{2 r}(y ; p) \equiv\left(p(x) y^{(r)}\right)^{(r)}, \\
\Lambda_{2 r-1}(y ; p) & \equiv \frac{1}{2}\left[\left(p(x) y^{(r-1)}\right)^{(r)}+\left(p(x) y^{(r)}\right)^{(r-1}\right], \quad(r=1,2, \cdots),
\end{aligned}
$$

with the understanding that $p \in \mathfrak{X}_{r}$ in the definition of $\Lambda_{2 r}$ and $\Lambda_{2 r-1}$. The primary result of that paper, and the one of most interest here, is Theorem 3.2, to the effect that if the polynomials $1, x, \cdots, x^{n} / n!$, where $n=\nu / 2$ or $n=(\nu+1) / 2$ according as $\nu$ is even or odd, belong to the domain of the adjoint operator $T_{0}^{*}$, then there exist functions $\pi_{j},(j=0, \cdots, \nu)$, with $\pi_{0} \in \mathfrak{U}_{0}, \pi_{2 \alpha-1} \in \mathfrak{U}_{\alpha}, \pi_{2 \alpha} \in \mathfrak{U}_{\alpha}$ such that $L[y]$ is given by

$$
\begin{equation*}
L[y]=\sum_{j=0}^{\nu} \Lambda_{j}\left(y ; \pi_{j}\right), \tag{7.2}
\end{equation*}
$$

while $\mathfrak{A}_{\nu}$ is contained in the domain of the adjoint operator $T_{0}^{*}$ and

$$
\begin{equation*}
T_{0}^{*}(z)=L^{\sharp}[z] \equiv \sum_{j=0}^{\nu} \Lambda_{j}\left(z ;(-1)^{j} \bar{\pi}_{j}\right) \quad \text { for } \quad z \in \mathfrak{A}_{\nu} . \tag{7.3}
\end{equation*}
$$

In view of the differentiability properties of $\pi_{j},(j=1, \cdots, \nu)$, it follows that (7.2) and (7.3) are of the form (6.2) and (6.4), respectively, which in turn reduce to (2.2) and (2.4), respectively, provided that $m=n, g_{i j}(x) \equiv 0$ when $i \geqq 1$ or $j \geqq 1$, and for $i, j=0, \cdots, n$ one defines $f_{i j}(x)$ as follows: $f_{i i}(x)=(-1)^{i} \pi_{2 i}(x) ; f_{i i-1}(x)=(-1)^{i}(1 / 2) \pi_{2 i-1}(x)$, $(i=1, \cdots, n) ; f_{i i+1}(x)=(-1)^{i}(1 / 2) \pi_{2 i+1}(x),(i=0, \cdots, n-1) ; f_{i j}(x) \equiv 0$, ( $j<i-1$ and $j>i+1$ ).

In particular, if $\nu=2 n$ and $\pi_{2 n}(x) \not \equiv 0$, then the vector $\hat{y}(x)$ consists of $y(x)$ and its first $n-1$ derivatives. Similarly, $\check{z}(x)$ consists of $z(x)$ and its first $n-1$ derivatives. The coordinates $\widetilde{y}_{i}(x)$ of the $n$-vector $\widetilde{y}(x)$ are defined by (2.1), and may be expressed in terms of $y(x)$ and its first $2 n-j$ derivatives, $(j=1, \cdots, n-1)$, and similarly for the coordinates of $\widetilde{z}(x)$, defined by (2.3). Consequently, $L[y]$ and $L^{凶}[z]$ are defined for $y, z \in \mathfrak{A}_{\nu}$.

If $\nu=2 n-1$, and $\pi_{\nu}(x) \not \equiv 0$, then $L$ is an operator of odd order and we modify the above defined matrix $\left[f_{i j}(x)\right]$ in the following way: delete the last row, replace $f_{n-1 n}(x)$ with $(-1)^{n-1} \pi_{2 n-1}(x)$, and replace $f_{n-1 n-1}(x)$ with $(-1)^{n-1}\left(\pi_{2 n-2}(x)+(1 / 2) \pi_{2 n-1}^{\prime}(x)\right)$. This change from an $(n+1) \times(n+1)$ matrix $\left[f_{i j}(x)\right]$ to the $n \times(n+1)$ matrix $\left[f_{i j}^{0}\right]$ changes neither the value of $L[y]$ nor the value of $L^{\text {㐫 }}[z]$. Now if $\pi_{2 n-1} \in \mathfrak{A}_{n}$,
then $\pi_{2 n-1}^{\prime} \in \mathfrak{U}_{n-1}$ so that $\widetilde{y}_{j}(x)$ may still be differentiated out and written in terms of $y$ and its first $2 n-j$ derivatives, $(j=1, \cdots, n-2)$, and similarly $\widetilde{z}_{i}(x),(i=1, \cdots, n-1)$, may be written in terms of $z(x)$ and its first $2 n-i$ derivatives. Consequently we still have that $L$ and $L^{\star}$ have the common domain $\mathfrak{N}_{\nu}$.

If now it is assumed that there is an $\varepsilon>0$ such that $\left|q_{\nu}(x)\right| \geqq \varepsilon$ almost everywhere, then it follows from Theorem 3.2, or Theorem 4.1 of [7], that the domain of the adjoint operator $T_{0}^{*}$ is $\mathfrak{A}_{\nu}$. Moreover, in view of the formulas which give the canonical variables $\widetilde{y}_{j}(x)$ and $\widetilde{z}_{i}(x)$ in terms of $y(x), \cdots, y^{(n-1)}(x)$ and $z(x), \cdots, z^{(m-1)}(x)$, respectively, we see that there exist nonsingular linear transformations $T$ and $T_{1}$ which transform the vector functions ( $y, y^{\prime}, \cdots, y^{(\nu-1)}$ ) and ( $z, z^{\prime}, \cdots$, $z^{(\nu-1)}$ ) into the vector functions ( $y, y^{\prime}, \cdots, y^{(n-1)}, \widetilde{y}_{1}, \cdots, \widetilde{y}_{m}$ ) and ( $z, z^{\prime}, \cdots$, $\left.\boldsymbol{z}^{(m-1)}, \widetilde{z}_{1}, \cdots, \widetilde{z}_{n}\right)$, respectively. Therefore, in view of Theorem 3.2 of Reid [7] and Theorem 6.1, it follows that boundary conditions for a $\nu$ th order differential operator of the type described above which involve linearly $y$ and its first $\nu-1$ derivatives at two points may be written as (6.14), and the adjoint boundary conditions may be written as (6.15).

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