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## Mariella Cecchi; Zuzana Došlá; Mauro Marini

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# SOME PROPERTIES OF THIRD ORDER <br> DIFFERENTIAL OPERATORS 

M. Cecchi, Firenze, Z. Došlá, Brno, M. Marinı, Firenze

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Abstract. Consider the third order differential operator $L$ given by

$$
L(\cdot) \equiv \frac{1}{a_{3}(t)} \frac{\mathrm{d}}{\mathrm{~d} t} \frac{1}{a_{2}(t)} \frac{\mathrm{d}}{\mathrm{~d} t} \frac{1}{a_{1}(t)} \frac{\mathrm{d}}{\mathrm{~d} t}(\cdot)
$$

and the related linear differential equation $L(x)(t)+x(t)=0$. We study the relations between $L$, its adjoint operator, the canonical representation of $L$, the operator obtained by a cyclic permutation of coefficients $a_{i}, i=1,2,3$, in $L$ and the relations between the corresponding equations.

We give the commutative diagrams for such equations and show some applications (oscillation, property A).

Keywords: Differential operators, linear differential equation of third order, canonical forms, adjoint equation, cyclic permutation, oscillatory solution, Kneser solution, property A

MSC 1991: 34C10, 34C20

## Introduction

Many authors have studied the oscillatory, nonoscillatory and asymptotic behaviour for the linear third order differential equation in the normal form, that is

$$
\begin{equation*}
x^{\prime \prime \prime}(t)+p(t) x^{\prime}(t)+q(t) x(t)=0 \tag{1}
\end{equation*}
$$

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where $p, q$ are continuous functions for $t \geqslant 0$. Among the numerous results dealing with this subject, we refer to the books [13, 16, 20] and the references contained therein. Other contributions may be obtained, as a particular case, from many papers concerning equations of higher order, see, e.g., [1, 17, 20] and/or nonlinear equations, e.g. [2, 5]. Some authors (see, e.g., $[12,14]$ ) consider the special forms

$$
\begin{align*}
& \left(r(t) x^{\prime}(t)\right)^{\prime \prime}+q(t) x(t)=0  \tag{2}\\
& \left(r(t) x^{\prime \prime}(t)\right)^{\prime}+q(t) x(t)=0 \tag{3}
\end{align*}
$$

Another often considered form of equation (1) arises when the second order differential equation

$$
\begin{equation*}
w^{\prime \prime}(t)+p(t) w(t)=0 \tag{4}
\end{equation*}
$$

is nonoscillatory. Let $h$ be a solution of (4), $h(t)>0$ for $t>T$; then, by standard computation, (1) may be written, for $t>T$, as

$$
\begin{equation*}
\left(h^{2}(t)\left(\frac{1}{h(t)} x^{\prime}(t)\right)^{\prime}\right)^{\prime}+q(t) h(t) x(t)=0 \tag{5}
\end{equation*}
$$

The natural question which arises is the relation among (2), (3), (5). To understand the various forms of the third order linear differential equation, we consider the equation

$$
\begin{equation*}
\left(\frac{1}{a_{2}(t)}\left(\frac{1}{a_{1}(t)} x^{\prime}(t)\right)^{\prime}\right)^{\prime}+a_{3}(t) x(t)=0 \tag{E}
\end{equation*}
$$

where
$(\mathrm{Hp}) \quad a_{i} \in C^{2}([0, \infty), \mathbb{R}), i=1,2,3, a_{1}(t)>0, a_{2}(t)>0, a_{3}(t) \neq 0$ on $[0, \infty)$.

We just remark that in the sequel the assumption " $a_{i} \in C^{2}([0, \infty), \mathbb{R})$ " is not always necessary, but it is considered for the sake of simplicity. For instance, if $a_{i} \in C^{3-i}([0, \infty), \mathbb{R}), i=1,2$, then (E) may be written in the normal form

$$
x^{\prime \prime \prime}(t)+r(t) x^{\prime \prime}(t)+p(t) x^{\prime}(t)+q(t) x(t)=0 .
$$

If $a_{i} \in C([0, \infty), \mathbb{R}), i=1,2,3$, then (E) may be interpreted as a first order linear differential system.

The assumption (Hp), as regards the sign of functions $a_{i}, i=1,2,3$, is motivated by the following fact: concerning the possible sign of coefficients of an equation of type ( E ), there are eight possible cases, but the significative ones with respect to the equation considered are only two, namely, for instance, $a_{1}>0, a_{2}>0, a_{3}>0$ or $a_{1}>0, a_{2}>0, a_{3}<0$.

Denote by $\left(\mathrm{E}_{+}\right)\left[\left(\mathrm{E}_{-}\right)\right]$the equation $(\mathrm{E})$ when $a_{3}(t)>0\left[a_{3}(t)<0\right]$ on $[0, \infty)$.
Equation (E) is related to the third order differential operator $L$ given by

$$
\begin{equation*}
L(x)(t)=\frac{1}{a_{3}(t)}\left(\frac{1}{a_{2}(t)}\left(\frac{1}{a_{1}(t)} x^{\prime}(t)\right)^{\prime}\right)^{\prime} \tag{6}
\end{equation*}
$$

Then (E) can be written as

$$
L(x)(t)+x(t)=0
$$

In a recent paper [7], other operators, which are obtained from $L$ by an ordered cyclic permutation of the functions $a_{i}$, have been considered. Further, for equations of higher order, in [7] some properties of such operators have been investigated and applied to the study of the oscillation as well as to the classification of solutions.

Together with $L$, other operators have often been considered in literature: in particular, the adjoint operator (see, e.g., [8]) and the operator in the canonical form ([22]). They are of particular interest in a large number of problems arising in the study of the qualitative behavior of solutions of (E), especially with regard to disconjugacy. Among numerous monographs dealing with this topic, we choose to refer to [13, 19].

The aim of this paper is to examine the relationships between the above quoted operators. Further, coming back to the related differential equations, we prove the commutation among them (Part I). In this way, we get some new criteria for (E) to have property A as well as the connection between (2) and (3).

Some definitions and notation are as follows:
A nontrivial solution of $(\mathrm{E})$ is said to be oscillatory if it has infinitely large zeros and nonoscillatory otherwise. Equation (E) is said to be oscillatory if there exists at least one oscillatory solution, and nonoscillatory if all its solutions are nonoscillatory.

Denote by $\mathcal{S}(\mathrm{E})$ the space of all solutions of $(\mathrm{E})$, by $\mathcal{N}(\mathrm{E})$ the set of all nonoscillatory solutions and by $\mathcal{O}(\mathrm{E})$ the set of all oscillatory solutions. By $x^{[i]}, i=1,2$, we denote the quasiderivatives of a solution $x$ of $(\mathrm{E})$. That is, if $x \in \mathcal{S}(\mathrm{E})$, then $x^{[1]}=\frac{1}{a_{1}} x^{\prime}, x^{[2]}=\frac{1}{a_{2}}\left(x^{[1]}\right)^{\prime}$.

Finally, let $u_{1}, u_{2}$ be two continuous functions different from zero on $[0, \infty)$ and define

$$
I\left(u_{i}\right)=\int_{0}^{\infty}\left|u_{i}(t)\right| \mathrm{d} t, \quad I\left(u_{i}, u_{j}\right)=\int_{0}^{\infty}\left|u_{i}(t)\right| \int_{0}^{t}\left|u_{j}(s)\right| \mathrm{d} s \mathrm{~d} t, \quad i, j=1,2 .
$$

The functions $u_{1}, u_{2}$ can be classified according to the convergence or divergence of the four integrals $I\left(u_{1}\right), I\left(u_{2}\right), I\left(u_{1}, u_{2}\right), I\left(u_{2}, u_{1}\right)$ into six possible cases:

$$
\begin{array}{ll}
\left(\mathrm{K}_{1}\right) & I\left(u_{1}\right)=I\left(u_{2}\right)=\infty \\
\left(\mathrm{K}_{2}\right) & I\left(u_{1}\right)<\infty, I\left(u_{2}\right)=I\left(u_{1}, u_{2}\right)=\infty, \\
\left(\mathrm{K}_{3}\right) & I\left(u_{2}\right)<\infty, I\left(u_{1}\right)=I\left(u_{2}, u_{1}\right)=\infty, \\
\left(\mathrm{K}_{4}\right) & I\left(u_{2}\right)<\infty, I\left(u_{2}, u_{1}\right)<\infty, I\left(u_{1}\right)=\infty,  \tag{7}\\
\left(\mathrm{K}_{5}\right) & I\left(u_{1}\right)<\infty, I\left(u_{1}, u_{2}\right)<\infty, I\left(u_{2}\right)=\infty, \\
\left(\mathrm{K}_{6}\right) & I\left(u_{1}\right)<\infty, I\left(u_{2}\right)<\infty
\end{array}
$$

## Part I. Relations between differential operators

Equation (E) is determined by the vector of functions

$$
\begin{equation*}
a=\left(a_{1}, a_{2}, a_{3}\right) . \tag{8}
\end{equation*}
$$

Moreover, ( E ) is related to the third order differential operator $L$ given in (6). We will say that $L$ is the operator associated with vector $a$ and $(\mathrm{E})$ is the equation associated with vector $a$. From now on we introduce the operator "cyclic permutation" of $L$, the "semi-adjoint operator" to $L$ and the "canonical form" of $L$.

## I-1 Cyclic permutation.

Let us consider the map $\mathcal{C}: C^{2}\left([0, \infty), \mathbb{R}^{3}\right) \rightarrow C^{2}\left([0, \infty), \mathbb{R}^{3}\right)$ given by

$$
\mathcal{C}: a=\left(a_{1}, a_{2}, a_{3}\right) \mapsto \mathcal{C}(a)=\left(a_{2}, a_{3}, a_{1}\right)
$$

and let $L^{\mathcal{C}}, L^{\mathcal{C C}}$ be the operators associated with the vectors $\mathcal{C}(a), \mathcal{C}(\mathcal{C}(a))$ respectively, that is the operators given by

$$
\begin{aligned}
L^{\mathcal{C}}(y)(t) & =\frac{1}{a_{1}(t)}\left(\frac{1}{a_{3}(t)}\left(\frac{1}{a_{2}(t)} y^{\prime}(t)\right)^{\prime}\right)^{\prime} \\
L^{\mathcal{C C}}(z)(t) & =\frac{1}{a_{2}(t)}\left(\frac{1}{a_{1}(t)}\left(\frac{1}{a_{3}(t)} z^{\prime}(t)\right)^{\prime}\right)^{\prime}
\end{aligned}
$$

We will call these operators cyclic permutation operators of $L$. Let $\left(\mathrm{E}^{\mathcal{C}}\right)$ be the equation associated with the vector $\mathcal{C}(a)$, that is the equation associated with the first cyclic permutation of the vector $a$ given in (8):

$$
\begin{equation*}
L^{\mathcal{C}}(y)(t)+y(t)=\frac{1}{a_{1}(t)}\left(\frac{1}{a_{3}(t)}\left(\frac{1}{a_{2}(t)} y^{\prime}(t)\right)^{\prime}\right)^{\prime}+y(t)=0 \tag{C}
\end{equation*}
$$

where $y=x^{[1]}=\frac{1}{a_{1}} x^{\prime}$, with $x$ a solution of (E).
The equation associated with the vector $\mathcal{C}(\mathcal{C}(a))$, i.e. the equation associated with the second cyclic permutation of (8), is

$$
\begin{equation*}
L^{\mathcal{C} C}(z)(t)+z(t)=\frac{1}{a_{2}(t)}\left(\frac{1}{a_{1}(t)}\left(\frac{1}{a_{3}(t)} z^{\prime}(t)\right)^{\prime}\right)^{\prime}+z(t)=0 \tag{CC}
\end{equation*}
$$

where $z=y^{[1]}=\frac{1}{a_{2}} y^{\prime}$, with $y$ a solution of $\left(\mathrm{E}^{\mathcal{C}}\right)$. This phenomenon is studied in detail in [7] for the general $n$.

## I-2 Semi-adjoint operator.

Consider now the map $\mathcal{A}: C^{2}\left([0, \infty), \mathbb{R}^{3}\right) \rightarrow C^{2}\left([0, \infty), \mathbb{R}^{3}\right)$ given by

$$
\mathcal{A}: a=\left(a_{1}, a_{2}, a_{3}\right) \mapsto \mathcal{A}(a)=\left(a_{2}, a_{1},-a_{3}\right)
$$

and let $L^{\mathcal{A}}$ be the operator associated with the vector $\mathcal{A}(a)$, that is the operator given by

$$
L^{\mathcal{A}}(u)(t)=-\frac{1}{a_{3}(t)}\left(\frac{1}{a_{1}(t)}\left(\frac{1}{a_{2}(t)} u^{\prime}(t)\right)^{\prime}\right)^{\prime}
$$

The equation associated with vector $\mathcal{A}(a)$ is $L^{\mathcal{A}}(u)(t)+u(t)=0$, i.e.,

$$
\begin{equation*}
\frac{1}{a_{3}(t)}\left(\frac{1}{a_{1}(t)}\left(\frac{1}{a_{2}(t)} u^{\prime}(t)\right)^{\prime}\right)^{\prime}-u(t)=0 \tag{A}
\end{equation*}
$$

where $u=x^{[1]} y-x y^{[1]}, x, y$ being two independent solutions of (E). A similar definition holds for the operator $L^{\mathcal{A A}}$, that is the operator associated with $\mathcal{A}(\mathcal{A}(a))$, and for the equation ( $\mathrm{E}^{\mathcal{A} \mathcal{A}}$ ).

We will call the operator $L^{\mathcal{A}}$ the semi-adjoint operator to $L$ since it is related to the operator $L^{*}$ adjoint to the operator $L$ defined as (see, e.g., [8])

$$
L^{*} \equiv-\frac{\mathrm{d}}{\mathrm{~d} t} \frac{1}{a_{1}} \frac{\mathrm{~d}}{\mathrm{~d} t} \frac{1}{a_{2}} \frac{\mathrm{~d}}{\mathrm{~d} t} \frac{.}{a_{3}}
$$

Indeed, the adjoint equation of $(\mathrm{E})$ is the equation $L^{*}(v)(t)+v(t)=0$ which can be written as $\left(\mathrm{E}^{\mathcal{A}}\right)$ via the $\operatorname{map} u=\frac{1}{a_{3}} v$.

## I-3 Canonical form.

Let $u_{i} \in C^{3-i}([0, \infty), \mathbb{R}), i=1,2,3$, with a fixed sign on $[0, \infty)$ and consider the third order differential operator $D$ given by

$$
D(x)(t)=\frac{1}{u_{3}(t)}\left(\frac{1}{u_{2}(t)}\left(\frac{1}{u_{1}(t)} x^{\prime}(t)\right)^{\prime}\right)^{\prime}
$$

The operator $D$ is disconjugate on $[0, \infty)$, that is every solution of the equation $D(x)(t)=0$ has at most two zeros on $[0, \infty)$.

Hence by Theorem 1 in [22], $D$ can be written in a certain canonical form, that is

$$
D \equiv \frac{1}{u_{3} b_{3}} \frac{\mathrm{~d}}{\mathrm{~d} t} \frac{1}{b_{2}} \frac{\mathrm{~d}}{\mathrm{~d} t} \frac{1}{b_{1}} \frac{\mathrm{~d}}{\mathrm{~d} t} \frac{-}{b_{0}}
$$

where the functions $b_{i}, i=0,1,2,3$, have a fixed sign, they are continuous, such that

$$
\int^{\infty}\left|b_{1}\right|=\infty, \quad \int^{\infty}\left|b_{2}\right|=\infty
$$

and determined up to positive multiplicative constants with product 1.
The explicit forms of functions $b_{i}$ depend on the convergence or divergence of the four integrals $I\left(u_{1}\right), I\left(u_{2}\right), I\left(u_{1}, u_{2}\right), I\left(u_{2}, u_{1}\right)$, that is on the six cases $\left(\mathrm{K}_{i}\right)$, $i=1, \ldots, 6$, given in (7). They may be calculated using the proof of Lemmas 1 and 2 in [22] and, for reader's convenience, they are given in Table 1.

Since the operator $L$ given in (6) is disconjugate, it can be written in the canonical form, that is

$$
L(\cdot) \equiv \frac{1}{a_{3} b_{3}} \frac{\mathrm{~d}}{\mathrm{~d} t} \frac{1}{b_{2}} \frac{\mathrm{~d}}{\mathrm{~d} t} \frac{1}{b_{1}} \frac{\mathrm{~d}}{\mathrm{~d} t} \frac{-}{b_{0}}
$$

Moreover, (E) may be written as

$$
\begin{equation*}
\left(\frac{1}{b_{2}(t)}\left(\frac{1}{b_{1}(t)}\left(\frac{x(t)}{b_{0}(t)}\right)^{\prime}\right)^{\prime}\right)^{\prime}+a_{3}(t) b_{3}(t) x(t)=0 \tag{1}
\end{equation*}
$$

Let us consider the map $\mathcal{K}: C^{2}\left([0, \infty), \mathbb{R}^{3}\right) \rightarrow C^{2}\left([0, \infty), \mathbb{R}^{3}\right)$ given by

$$
\mathcal{K}_{i}: a=\left(a_{1}, a_{2}, a_{3}\right) \mapsto \mathcal{K}_{i}(a)=\left(b_{1}, b_{2}, a_{3} b_{0} b_{3}\right)
$$

where $b_{j}$ are the coefficients of the canonical form of $L$ in case $\left(\mathrm{K}_{i}\right), i=2,3,4,5,6$ (see Table 1 with $u_{1}=a_{1}, u_{2}=a_{2}$ ).

Now let us consider the differential operator $L^{\mathcal{K}_{i}}$ associated with the vector $\mathcal{K}_{i}(a)$ given by

$$
L^{\mathcal{K}_{i}}(v)(t)=\left(\frac{1}{a_{3}(t) b_{0}(t) b_{3}(t)} \frac{\mathrm{d}}{\mathrm{~d} t} \frac{1}{b_{2}(t)} \frac{\mathrm{d}}{\mathrm{~d} t} \frac{1}{b_{1}(t)} \frac{\mathrm{d}}{\mathrm{~d} t} v(t)\right) .
$$

The equation associated with the vector $\mathcal{K}_{i}(a)$ is $L^{\mathcal{K}_{i}}(v)(t)+v(t)=0$, i.e.,
$\left(\mathrm{E}^{\mathcal{K}_{i}}\right) \quad\left(\frac{1}{b_{2}(t)}\left(\frac{1}{b_{1}(t)} v^{\prime}(t)\right)^{\prime}\right)^{\prime}+a_{3}(t) b_{0}(t) b_{3}(t) v(t)=0$,
which gives equation $\left(\mathrm{E}_{1}\right)$ via the map $v=\frac{1}{b_{0}} x$. If we do not need to specify the case ( $\mathrm{K}_{i}$ ), then we will use the notation ( $\mathrm{E}^{\mathcal{K}}$ ) instead of ( $\mathrm{E}^{\mathcal{K}_{i}}$ ).

In I-4 we focus our attention on cases $\left(\mathrm{K}_{1}\right)-\left(\mathrm{K}_{5}\right)$; the case $\left(\mathrm{K}_{6}\right)$ is discussed in Part II, Remark 7.

In accordance with the above notation, the identity " $\equiv$ " between two equations means that the associated vectors are the same. Obviously, if ( $\mathrm{K}_{1}$ ) holds, then $\left(\mathrm{E}^{\mathcal{K}_{1}}\right) \equiv(\mathrm{E})$. If $\left(\mathrm{K}_{2}\right)$ holds, then the canonical form of $L$ yields the equation
$\left(\mathrm{E}^{\mathcal{K}_{2}}\right) \quad\left(\frac{1}{a_{2}(t) \int_{t}^{\infty} a_{1}}\left(\frac{\left(\int_{t}^{\infty} a_{1}\right)^{2}}{a_{1}(t)} v^{\prime}\right)^{\prime}\right)^{\prime}+a_{3}(t)\left(\int_{t}^{\infty} a_{1}\right) v=0$,
where $v=\left(\int_{t}^{\infty} a_{1}\right)^{-1} x$ with $x$ a solution of $(\mathrm{E})$. If $\left(\mathrm{K}_{3}\right)$ holds, then $\left(\mathrm{E}^{\mathcal{K}}\right)$ has the form

$$
\left(\mathrm{E}^{\mathcal{K}_{3}}\right) \quad\left(\frac{\left(\int_{t}^{\infty} a_{2}\right)^{2}}{a_{2}(t)}\left(\frac{x^{\prime}}{a_{1}(t) \int_{t}^{\infty} a_{2}}\right)^{\prime}\right)^{\prime}+a_{3}(t)\left(\int_{t}^{\infty} a_{2}\right) x=0
$$

Similarly we obtain $\left(\mathrm{E}^{\mathcal{K}_{4}}\right)$ and $\left(\mathrm{E}^{\mathcal{K}_{5}}\right)$.

Remark 1. Coefficient $b_{0}$ in the canonical forms is either the constant function $b_{0}(t) \equiv 1$ (in cases $\left.\left(\mathrm{K}_{1}\right),\left(\mathrm{K}_{3}\right)\right)$ or a positive decreasing function in all other cases. Indeed, in cases $\left(\mathrm{K}_{i}\right), i=2,4,5$ we have $b_{0}(t)=\int_{t}^{\infty} f(s) \mathrm{d} s$, where $f>0, \int^{\infty} f<\infty$ and thus $b_{0}^{\prime}(t)=-f(t)<0$. Similarly, by routine computation we get the conclusion in case ( $\mathrm{K}_{6}$ ).

## I-4 Relations between the corresponding equations.

Let us examine all equations under consideration:
(E)

$$
\begin{array}{ll}
(\mathrm{E}) & \frac{1}{a_{3}(t)}\left(\frac{1}{a_{2}(t)}\left(\frac{1}{a_{1}(t)} x^{\prime}(t)\right)^{\prime}\right)^{\prime}+x(t)=0, \\
\left(\mathrm{E}^{\mathcal{C}}\right) & \frac{1}{a_{1}(t)}\left(\frac{1}{a_{3}(t)}\left(\frac{1}{a_{2}(t)} y^{\prime}(t)\right)^{\prime}\right)^{\prime}+y(t)=0, \\
\left(\mathrm{E}^{\mathcal{A}}\right) & \frac{1}{a_{3}(t)}\left(\frac{1}{a_{1}(t)}\left(\frac{1}{a_{2}(t)} u^{\prime}(t)\right)^{\prime}\right)^{\prime}-u(t)=0, \\
\left(\mathrm{E}^{\mathcal{K}_{i}}\right) & \left(\frac{1}{b_{2}(t)}\left(\frac{1}{b_{1}(t)} v^{\prime}(t)\right)^{\prime}\right)^{\prime}+b_{0}(t) b_{3}(t) a_{3}(t) v(t)=0
\end{array}
$$

where $b_{j}$ are the coefficients of the canonical form in case $\left(\mathrm{K}_{i}\right)$ (see Table 1 with $u_{1}=a_{1}, u_{2}=a_{2}$ ). Our aim here is to describe some relations between these differential equations. In the sequel we will say that (E) yields, via the map $\mathcal{C},\left(\mathrm{E}^{\mathcal{C}}\right)$ (in the sense of associated vectors) and make use of the notation

$$
(\mathrm{E}) \xrightarrow{\mathcal{C}}\left(\mathrm{E}^{\mathcal{C}}\right) .
$$

Similarly for the maps $\mathcal{A}, \mathcal{K}_{i}$ and their composition. Hence

$$
(\mathrm{E}) \xrightarrow{\mathcal{A}}\left(\mathrm{E}^{\mathcal{A}}\right), \quad(\mathrm{E}) \xrightarrow{\mathcal{K}}\left(\mathrm{E}^{\mathcal{K}}\right) .
$$

In what follows we denote by $L^{\mathcal{A C}}\left[\left(\mathrm{E}^{\mathcal{A C}}\right)\right]$ the operator [equation] associated with the vector $\mathcal{C}(\mathcal{A}(a))$ and similarly for the other composition of the maps $\mathcal{C}, \mathcal{A}, \mathcal{K}_{i}$.

Concerning the explicit forms of the coefficients $b_{i}$, defining the canonical forms in cases under consideration, they may be obtained, by standard calculation, from Table 1. For instance, the explicit form of coefficients $b_{i}$ which define the operator $L^{\mathcal{C}} \mathcal{K}_{i}$ associated with the vector $\mathcal{K}_{i}(\mathcal{C}(a))$ are given in Table 2.

Theorems 1-2 below present other results related to these equations and Theorem 3 below gives some relations between the corresponding spaces of solutions.

Theorem 1. The maps $\mathcal{A}, \mathcal{C}$ and $\mathcal{K}_{i}, i=2,3,4,5$ make the following diagrams commutative


Diagram 1

$(\mathrm{E}) \stackrel{\mathcal{A}}{\longleftrightarrow}\left(\mathrm{E}^{\mathcal{A}}\right)$


Diagram 2


Diagram 3

Proof. Equation (E) yields, via the map $\mathcal{K}_{i} \circ \mathcal{A}$, the equation

$$
\left(\mathrm{E}^{\mathcal{K} \mathcal{A}}\right) \quad\left(\frac{1}{b_{1}(t)}\left(\frac{1}{b_{2}(t)} u^{\prime}(t)\right)^{\prime}\right)^{\prime}-a_{3}(t) b_{0}(t) b_{3}(t) u(t)=0
$$

Similarly, the equation (E) yields, via the $\operatorname{map} \mathcal{C} \circ \mathcal{A}$, the equation $\left(\mathrm{E}^{\mathcal{C A}}\right) \quad\left(\frac{1}{a_{2}(t)}\left(\frac{1}{a_{3}(t)} y^{\prime}(t)\right)^{\prime}\right)^{\prime}-a_{1}(t) y(t)=0$.

Diagram 1: We have $(\mathrm{E}) \equiv\left(\mathrm{E}^{\mathcal{C A C A}}\right),\left(\mathrm{E}^{\mathcal{A}}\right) \equiv\left(\mathrm{E}^{\mathcal{A} \mathcal{A C A C}}\right),\left(\mathrm{E}^{\mathcal{A}}\right) \equiv\left(\mathrm{E}^{\mathcal{C A C}}\right)$.
Diagram 2: Using Table 1 we get $\left(\mathrm{E}^{\mathcal{K}_{3} \mathcal{A}}\right) \equiv\left(\mathrm{E}^{\mathcal{A} \mathcal{K}_{2}}\right)$ and $\left(\mathrm{E}^{\mathcal{K}_{2} \mathcal{A}}\right) \equiv\left(\mathrm{E}^{\mathcal{A} \mathcal{K}_{3}}\right)$.
By the similar symmetry between the cases $\left(\mathrm{K}_{4}\right)$ and $\left(\mathrm{K}_{5}\right)$ (see Table 1) we get Diagram 3.

Theorem 2. 1. Let $I\left(a_{2}\right)<\infty, I\left(a_{2}, a_{1}\right)=I\left(a_{2}, a_{3}\right)=\infty$. Then the maps $\mathcal{A}, \mathcal{C}$ and $\mathcal{K}_{i}, i=2,3$ make the following diagram commutative

2. Let $I\left(a_{3}\right)<\infty, I\left(a_{3}, a_{1}\right)=I\left(a_{3}, a_{2}\right)=\infty$. Then the maps $\mathcal{A}, \mathcal{C}$ and $\mathcal{K}_{i}, i=2,3$ make the following diagram commutative

3. Let $I\left(a_{1}\right)<\infty, I\left(a_{1}, a_{3}\right)=I\left(a_{1}, a_{2}\right)=\infty$. Then the maps $\mathcal{A}, \mathcal{C}$ and $\mathcal{K}_{i}, i=2,3$ make the following diagram commutative


Proof. To prove the whole diagram is commutative, it is sufficient to prove that two square diagrams commute. The first concerns the composed transformations $\mathcal{K}_{2} \circ \mathcal{C}$ and $\mathcal{C} \circ \mathcal{K}_{3}$, the second $\mathcal{C} \circ \mathcal{A}$ and $\mathcal{A} \circ \mathcal{C}$.

Let $I\left(a_{2}\right)<\infty, I\left(a_{2}, a_{1}\right)=I\left(a_{2}, a_{3}\right)=\infty$. We calculate the canonical form of (E) [of $\left(\mathrm{E}^{\mathcal{C}}\right)$ ] according to Table 1 , where the case $\left(\mathrm{K}_{3}\right)$ [ $\left(\mathrm{K}_{2}\right)$, respectively] is considered. By simple computation we get $\left(\mathrm{E}^{\mathcal{C} \mathcal{K}_{2}}\right) \equiv\left(\mathrm{E}^{\mathcal{K}_{3} \mathcal{C}}\right)$. The second part of the diagram follows from Diagram 1 of Theorem 1, where ( E ) is replaced by ( $\mathrm{E}^{\mathcal{K}_{3}}$ ).

The other two cases follow from the cyclic permutation of $(E)$ and $\left(E^{\mathcal{C}}\right)$, respectively.

Transformations $\mathcal{C}$ and $\mathcal{K}_{i}$ induce isomorphisms $T_{\mathcal{C}}: \mathcal{S}(\mathrm{E}) \rightarrow \mathcal{S}\left(\mathrm{E}^{\mathcal{C}}\right)$, defined by

$$
T_{\mathcal{C}}: x \in \mathcal{S}(\mathrm{E}) \mapsto x^{[1]} \in \mathcal{S}\left(\mathrm{E}^{\mathcal{C}}\right)
$$

and $T_{\mathcal{K}}: \mathcal{S}(\mathrm{E}) \rightarrow \mathcal{S}\left(\mathrm{E}^{\mathcal{K}}\right)$, defined by

$$
\begin{equation*}
T_{\mathcal{K}}: x \in \mathcal{S}(\mathrm{E}) \mapsto \frac{x}{b_{0}} \in \mathcal{S}\left(\mathrm{E}^{\mathcal{K}}\right) \tag{8}
\end{equation*}
$$

These isomorphisms preserve oscillation, as the following theorem shows.

Theorem 3. The maps $T_{\mathcal{C}}$ and $T_{\mathcal{K}}$ are isomorphisms which maintain the oscillatory properties of the solutions. That is,

$$
x \in \mathcal{O}(\mathrm{E}) \Longleftrightarrow T_{\mathcal{C}}(x) \in \mathcal{O}\left(\mathrm{E}^{\mathcal{C}}\right), \quad x \in \mathcal{O}(\mathrm{E}) \Longleftrightarrow T_{\mathcal{K}}(x) \in \mathcal{O}\left(\mathrm{E}^{\mathcal{K}}\right)
$$

Equation $\left(\mathrm{E}^{\mathcal{A}}\right)$ satisfies

$$
\mathcal{O}(\mathrm{E}) \neq \emptyset \Longleftrightarrow \mathcal{O}\left(\mathrm{E}^{\mathcal{A}}\right) \neq \emptyset
$$

Proof. The statement on $T_{\mathcal{C}}$ is proved in [7], Theorem 1 (Correspondence Principle). We will prove the statement for $T_{\mathcal{K}}$. According to (8), if $x \in \mathcal{S}(\mathrm{E})$, then $y=x / b_{0} \in \mathcal{S}\left(\mathrm{E}^{\mathcal{K}}\right)$. It is obvious that if $x_{1} \neq x_{2}$ then $y_{1} \neq y_{2}$ and if $y \in \mathcal{S}\left(\mathrm{E}^{\mathcal{K}}\right)$, then the function $x=b_{0} y$ is from $\mathcal{S}(\mathrm{E})$, i.e. $T_{\mathcal{K}}$ is isomophism. Further, $x$ is oscillatory iff $y$ is oscillatory. The statement for the adjoint equation is proved in [10].

Remark 2. (i) The map $\mathcal{A}$ does not induce an isomorphism between spaces of solutions because these spaces have different structure with respect to the oscillation (see, e.g., $[6,10]$ ).
(ii) The existence of an oscillatory solution of any equation in the diagrams derived in Theorems 1 and 2, ensures by Theorem 3 that all other equations in these diagrams are oscillatory, too. Similarly this holds for nonoscillatory (disconjugacy) criteria for (E). ((E) is nonoscillatory if and only if it is eventually disconjugate, i.e., there exists $T \in \mathbb{R}$ such that every nontrivial solution has at most two zeros on $[T, \infty)$, see [10].)
(iii) From the fact that $b_{0}(t) \rightarrow 0$ as $t \rightarrow \infty$ in cases $\left(\mathrm{K}_{2}\right),\left(\mathrm{K}_{4-6}\right)$ of canonical forms and from (8) it follows that in all cases $\left(\mathrm{K}_{i}\right)(i=2, \ldots, 6), T_{\mathcal{K}}^{-1}$ preserves solutions tending to zero and $T_{\mathcal{K}}$ preserves unbounded solutions.

## Part II. Applications: Property A

In [14, Th. $5.10 ; 18$, Th. 1.2] it has been observed that if $q>0, p \leqslant 0$ and equation (1) has an oscillatory solution then every nonoscillatory solution $x$ satisfies

$$
\operatorname{sgn} x \operatorname{sgn} x^{\prime}<0, \quad \operatorname{sgn} x \operatorname{sgn} x^{\prime \prime}>0 \quad \text { on }[0, \infty)
$$

Such a nonoscillatory solution is called a Kneser solution and the property that every solution of equation (1) is either an oscillatory solution or a Kneser solution tending to zero for $t \rightarrow \infty$ is called property $A$, see [16].

For this reason we consider equation ( $\mathrm{E}_{+}$) and study the problem how the transformations $T_{\mathcal{K}}^{-1}, T_{\mathcal{C}}$ preserve property A of ( $\mathrm{E}_{+}$).

The set $\mathcal{N}\left(\mathrm{E}_{+}\right)$of all nontrivial nonoscillatory solutions of ( $\mathrm{E}_{+}$) can be divided into the following four classes:

$$
\begin{aligned}
& N_{0}=\left\{x \in \mathcal{N}\left(\mathrm{E}_{+}\right), \exists t_{x}: x(t) x^{[1]}(t)<0, x(t) x^{[2]}(t)>0 \quad \text { for } t \geqslant t_{x}\right\} \\
& N_{1}=\left\{x \in \mathcal{N}\left(\mathrm{E}_{+}\right), \exists t_{x}: x(t) x^{[1]}(t)>0, x(t) x^{[2]}(t)<0 \quad \text { for } t \geqslant t_{x}\right\} \\
& N_{2}=\left\{x \in \mathcal{N}\left(\mathrm{E}_{+}\right), \exists t_{x}: x(t) x^{[1]}(t)>0, x(t) x^{[2]}(t)>0 \text { for } t \geqslant t_{x}\right\} \\
& N_{3}=\left\{x \in \mathcal{N}\left(\mathrm{E}_{+}\right), \exists t_{x}: x(t) x^{[1]}(t)<0, x(t) x^{[2]}(t)<0 \quad \text { for } t \geqslant t_{x}\right\}
\end{aligned}
$$

Analogously to equation (1), solutions in the class $N_{0}$ are called Kneser solutions (see, e.g., [16]).

Definition. Equation ( $\mathrm{E}_{+}$) is said to have property $A$ if each of its solutions either is oscillatory or is a Kneser solution tending to zero for $t \rightarrow \infty$.

Remark 3. When equations $\left(E_{+}\right),\left(E_{+}^{\mathcal{C}}\right)$ and $\left(E_{+}^{\mathcal{C C}}\right)$ have property $A$, then any Kneser solution $x$ of any of these equations satisfies $\lim _{t \rightarrow \infty} x^{[i]}=0, i=0,1,2$, where $x^{[0]}=x$. This property coincides with property A defined for the system of differential equations, see [15].

Lemma 1. If $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ has property $A$ then $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ and $\left(\mathrm{E}_{+}\right)$are oscillatory.
Proof. We have $\mathcal{N}\left(\mathrm{E}_{+}^{\mathcal{K}}\right)=N_{0} \cup N_{2} ; N_{2} \neq \emptyset$ if and only if $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ is nonoscillatory, see Lemmas 2.1 and 2.2 in [16]. Hence if $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ has property $A$, then $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ is oscillatory and so, by Theorem $3,\left(\mathrm{E}_{+}\right)$is oscillatory, too.

Lemma 2. If $I\left(a_{1}\right)=\infty$ or $I\left(a_{1}, a_{2}\right)=\infty$ then $N_{3}=\emptyset$.
Proof. Let $N_{3} \neq \emptyset$. Suppose that

$$
x(t)<0, \quad x^{[1]}(t)>0, \quad x^{[2]}(t)>0 \quad \text { for large } t .
$$

Thus $x^{[3]}(t)>0$ and we have

1. $x<0$ and $x^{[1]}>0$ implies that $x<0$ is increasing, so $x$ is bounded.
2. $x^{[1]}>0$ and $x^{[2]}>0$ implies that $x^{[1]}>0$ is increasing, so

$$
x(t)-x(T)>x^{[1]}(T) \int_{T}^{t} a_{1}(s) \mathrm{d} s \quad \text { for } t>T
$$

If $I\left(a_{1}\right)=\infty$, then this is a contradiction with the fact that $x$ is bounded.
3. $x^{[2]}>0$ and $x^{[3]}>0$ implies that $x^{[2]}>0$ is increasing. Thus we get, integrating twice from $T$ to $t, t>T$,

$$
\begin{gathered}
x^{[2]}(t)>x^{[2]}(T)>0 \\
x^{[1]}(t)>x^{[1]}(T)+x^{[2]}(T) \int_{T}^{t} a_{2}(s) \mathrm{d} s \\
x(t)-x(T)>x^{[1]}(T)+x^{[2]}(T) \int_{T}^{t} a_{1}(s) \int_{T}^{s} a_{2}(s) \mathrm{d} s \mathrm{~d} t \rightarrow \infty,
\end{gathered}
$$

which contradicts the boundedness of $x$.

Proposition 1. (i) If $I\left(a_{i}\right)=\infty, i=1,2,3$, then ( $\mathrm{E}_{+}$) has property $A$.
(ii) The maps induced by cyclic permutations preserve Kneser solutions, i.e., Kneser solutions of $\left(\mathrm{E}_{+}\right)$are transformed to Kneser solutions of $\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ and so on.
(iii) Let $I\left(a_{1}\right)=I\left(a_{2}\right)=\infty$. Then

$$
\begin{equation*}
\int_{0}^{\infty} a_{3}(t) \int_{0}^{t} a_{2}(s) \int_{0}^{s} a_{1}(u) \mathrm{d} u \mathrm{~d} s \mathrm{~d} t=\infty \tag{9}
\end{equation*}
$$

if and only if every Kneser solution of ( $\mathrm{E}_{+}$) tends to zero for $t \rightarrow \infty$.
(iv) If

$$
\begin{equation*}
\int_{0}^{\infty} a_{3}(t) \int_{0}^{t} a_{1}(s) \int_{0}^{s} a_{2}(u) \mathrm{d} u \mathrm{~d} s \mathrm{~d} t<\infty \tag{10}
\end{equation*}
$$

then ( $\mathrm{E}_{+}$) is nonoscillatory.
Proof. (i) See, e.g., $[15,16]$.
(ii) See [7, Collorary 3].
(iii) See [21, Theorem 7].
(iv) See $[6$, Theorem 5].

Lemma 3. Let $\left(\mathrm{E}_{+}\right)$have the canonical form $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$. In cases $\left(\mathrm{K}_{i}\right), i=2,3,4$, the map $T_{\mathcal{K}}^{-1}: \mathcal{S}\left(\mathrm{E}_{+}^{\mathcal{K}}\right) \rightarrow \mathcal{S}\left(\mathrm{E}_{+}\right)$preserves property $A$, i.e., if $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ has property $A$ then ( $\mathrm{E}_{+}$) has property $A$.

In the other cases $\left(\mathrm{K}_{i}\right), i=5,6$, if $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ has property $A$ then $\mathcal{N}\left(\mathrm{E}_{+}\right)=N_{0} \cup N_{3}$ and every nonoscillatory solution of ( $\mathrm{E}_{+}$) tends to zero for $t \rightarrow \infty$.

Proof. Let $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ have property A. Then every nonoscillatory solution $u$ of $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ is a Kneser solution tending to zero. Then $x=b_{0} u$ is a solution of $\left(\mathrm{E}_{+}\right)$and using Remark 1 we get

$$
x^{[1]}=\frac{1}{a_{1}}\left(b_{0} u\right)^{\prime}=\frac{1}{a_{1}}\left(b_{0}^{\prime} u+b_{0} u^{\prime}\right)<0
$$

hence $x \in N_{0} \cup N_{3}$. In addition, using Remark 2 (iii), $x$ tends to zero. If case ( $\mathrm{K}_{i}$ ), $i=2,3,4$, occurs, then using Lemma 2 we get $N_{3}=\emptyset$.

Lemma 4. Let $I\left(a_{1}\right)=\infty$. Then $T_{\mathcal{C}}: \mathcal{S}\left(\mathrm{E}_{+}\right) \rightarrow \mathcal{S}\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ preserves property $A$, i.e., if $\left(\mathrm{E}_{+}\right)$has property $A$ then ( $\mathrm{E}_{+}^{\mathcal{C}}$ ) has property $A$.

Similarly, if $I\left(a_{2}\right)=\infty$ and $\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ has property $A$, then $\left(\mathrm{E}_{+}^{C \mathcal{C}}\right)$ has property $A$; if $I\left(a_{3}\right)=\infty$ and $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ has property $A$, then $\left(\mathrm{E}_{+}\right)$has property $A$.

Proof. Let $I\left(a_{1}\right)=\infty$. Suppose that ( $\mathrm{E}_{+}$) has property A. Then by Theorem 3 and Proposition 1(ii) we get that every solution of $\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ is either oscillatory or a

Kneser solution. Suppose that $y$ is a Kneser solution of $\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ tending to nonzero constant $c>0$. Since $y=x^{[1]}>0$, with $x$ a solution of ( $\mathrm{E}_{+}$) and $x<0$, we have for large $t$ and $\varepsilon>0$ sufficiently small

$$
0>x(t)=x(0)+\int_{0}^{t} a_{1}(s) x^{[1]}(s) \mathrm{d} s>x(0)+(c-\varepsilon) \int_{0}^{t} a_{1}(s) \mathrm{d} s \rightarrow \infty \text { for } t \rightarrow \infty
$$

which is a contradiction. Similarly we get the other two conclusions.

Lemma 5. Let ( $\mathrm{E}_{+}$) have the canonical form $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$. Suppose

$$
\begin{equation*}
\int^{\infty} a_{3}(t) b_{0}(t) b_{3}(t) \mathrm{d} t=\infty \tag{11}
\end{equation*}
$$

Then $\left(\mathrm{E}_{+}^{\mathcal{K}}\right)$ has property $A$.
Proof. Since $\int^{\infty} b_{i}=\infty, i=1,2$, and (11) holds, the conclusion follows from Proposition 1(i).

Theorem 4. Let one of the following three conditions hold:
(i) $I\left(a_{1}\right)=I\left(a_{3}\right)=I\left(a_{2}, a_{3}\right)=I\left(a_{2}, a_{1}\right)=\infty$
(ii) $I\left(a_{1}\right)=I\left(a_{2}\right)=I\left(a_{3}, a_{1}\right)=I\left(a_{3}, a_{2}\right)=\infty$
(iii) $\quad I\left(a_{2}\right)=I\left(a_{3}\right)=I\left(a_{1}, a_{2}\right)=I\left(a_{1}, a_{3}\right)=\infty$.

Then $\left(\mathrm{E}_{+}\right),\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ and $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ are oscillatory and have property $A$.
Proof. First, let us remark that if $I\left(a_{1}\right)=I\left(a_{2}\right)=I\left(a_{3}\right)=\infty$, the conclusion holds in view of Proposition 1(i) and Lemma 1.

Hence, suppose that $I\left(a_{2}\right)<\infty$ and (i) holds. Then ( $\mathrm{E}_{+}$) may be written in the canonical form $\left(\mathrm{E}_{+}^{\mathcal{K}_{3}}\right)$. Since $I\left(a_{2}, a_{3}\right)=\infty$, the condition (11) is satisfied and we get from Lemma 5 that ( $\mathrm{E}_{+}^{\mathcal{K}_{3}}$ ) has property A. Hence, by Lemma 3 , ( $\mathrm{E}_{+}$) has property A and, in addition, by Lemma 1, it has an oscillatory solution. By Theorem 3, ( $\mathrm{E}_{+}^{\mathcal{C}}$ ) and $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ are oscillatory, too. Since $I\left(a_{2}, a_{1}\right)=\infty$, we have $I\left(a_{1}\right)=\infty$ and in view of Lemma 4 the equation ( $E_{+}^{\mathcal{C}}$ ) has property A.

It remains to prove that also $\left(E_{+}^{\mathcal{C C}}\right)$ has property $A$. From the fact that $\left(E_{+}^{\mathcal{C}}\right)$ has property A it follows, in view of Proposition 1(ii), that every nontrivial solution of $\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ is either oscillatory or a Kneser solution. Since $I\left(a_{2}, a_{1}\right)=\infty$, we have

$$
\int_{0}^{\infty} a_{2}(t) \int_{0}^{t} a_{1}(s) \int_{0}^{s} a_{3}(u) \mathrm{d} u \mathrm{~d} s \mathrm{~d} t=\infty
$$

Using Proposition 1(iii) for $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$, we get that every Kneser solution of $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ tends to zero.

Similarly, if $I\left(a_{3}\right)<\infty$ and (ii) [or $I\left(a_{1}\right)<\infty$ and (iii)] hold, we apply Lemma 5 for $\left(E_{+}^{\mathcal{C K}}\right.$ ) [for $\left(\mathrm{E}_{+}^{\mathcal{C C}}{ }_{3}\right)$ ]. Then we use Lemma 3 and get the conclusion for ( $\mathrm{E}_{+}^{\mathcal{C}}$ ) [or $\left.\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)\right]$. Using Lemma 4 we get the conclusion for $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ [for ( $\left.\left.\mathrm{E}_{+}\right)\right]$and finally, using Proposition 1(iii), for $\left(\mathrm{E}_{+}\right)\left[\left(\mathrm{E}_{+}^{\mathcal{C}}\right)\right]$, respectively.

Theorem 5. Let one of the following three conditions hold:
(i) $I\left(a_{1}\right)=\infty, I\left(a_{2}, a_{1}\right)<\infty$ and

$$
\int_{0}^{\infty} a_{3}(t)\left(\int_{t}^{\infty} a_{2}(s) \mathrm{d} s\right)\left(\int_{t}^{\infty} a_{1}(s) \int_{s}^{\infty} a_{2}(\tau) d \tau \mathrm{~d} s\right) \mathrm{d} t=\infty
$$

(ii) $I\left(a_{2}\right)=\infty, \quad I\left(a_{3}, a_{2}\right)<\infty$ and

$$
\int_{0}^{\infty} a_{1}(t)\left(\int_{t}^{\infty} a_{3}(s) \mathrm{d} s\right)\left(\int_{t}^{\infty} a_{2}(s) \int_{s}^{\infty} a_{3}(\tau) d \tau \mathrm{~d} s\right) \mathrm{d} t=\infty
$$

(iii) $I\left(a_{3}\right)=\infty, \quad I\left(a_{1}, a_{3}\right)<\infty$ and

$$
\int_{0}^{\infty} a_{2}(t)\left(\int_{t}^{\infty} a_{1}(s) \mathrm{d} s\right)\left(\int_{t}^{\infty} a_{3}(s) \int_{s}^{\infty} a_{1}(\tau) d \tau \mathrm{~d} s\right) \mathrm{d} t=\infty
$$

Then $\left(\mathrm{E}_{+}\right),\left(\mathrm{E}_{+}^{\mathcal{C}}\right),\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ are oscillatory and have property $A$.
Proof. The idea of the proof is similar to that of Theorem 4.
Suppose (i). Then ( $\mathrm{E}_{+}$) may be written in the canonical form ( $\mathrm{E}_{+}^{\mathcal{K}_{4}}$ ). By Lemma 5 , $\left(\mathrm{E}_{+}^{\mathcal{K}_{4}}\right)$ has property A and so, by Lemma 3 , also ( $\mathrm{E}_{+}$) has property A. In addition, by Lemma $1,\left(\mathrm{E}_{+}\right)$is oscillatory. Using Lemma 4 we get the conclusion for ( $\left.\mathrm{E}_{+}^{\mathcal{C}}\right)$.

Consider ( $\mathrm{E}_{+}^{\mathcal{C}}$ ). Since ( $\mathrm{E}_{+}^{\mathcal{C}}$ ) has property A, it follows, in view of Proposition 1 (ii) and Theorem 3, that every nontrivial solution of $\left(E_{+}^{\mathcal{C C}}\right)$ is either oscillatory or a Kneser solution. The assumptions of the theorem imply that

$$
\begin{gathered}
\infty=\int_{0}^{\infty} a_{3}(t) \int_{t}^{\infty} a_{1}(s) \int_{s}^{\infty} a_{2}(u) \mathrm{d} u \mathrm{~d} s \mathrm{~d} t \\
=\int_{0}^{\infty} a_{3}(t) \int_{t}^{\infty} a_{2}(s) \int_{0}^{s} a_{1}(u) \mathrm{d} u \mathrm{~d} s \mathrm{~d} t=\int_{0}^{\infty} a_{2}(t) \int_{0}^{t} a_{1}(s) \int_{0}^{s} a_{3}(u) \mathrm{d} u \mathrm{~d} s \mathrm{~d} t
\end{gathered}
$$

so, using Proposition 1(iii) as in the proof of Theorem 4, every Kneser solution of $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ tends to zero. Then $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ is oscillatory and has property $A$. In the two remaining cases we proceed similarly.

Remark 4. Theorems 4 and 5 generalize for $n=3$ Corollary 5 of [7] where it is proved that if $I\left(a_{i}\right)=\infty, i=1,2,3$, then $\left(\mathrm{E}_{+}\right),\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ and $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ have property A.

Remark 5. Concerning oscillatory properties, we have proved in [6, Theorems 8 and 10] a more general statement: $\left(\mathrm{E}_{+}\right),\left(\mathrm{E}_{+}^{\mathcal{C}}\right)$ and $\left(\mathrm{E}_{+}^{\mathcal{C C}}\right)$ are oscillatory provided any of the following condition is satisfied:

$$
\begin{align*}
\text { (i) } \quad I\left(a_{1}\right)=I\left(a_{3}\right)=I\left(a_{2}, a_{3}\right)=\infty, \\
\text { (ii) } \quad I\left(a_{1}\right)=I\left(a_{2}\right)=I\left(a_{3}, a_{1}\right)=\infty,  \tag{12}\\
\text { (iii) } \quad I\left(a_{2}\right)=I\left(a_{3}\right)=I\left(a_{1}, a_{2}\right)=\infty .
\end{align*}
$$

Using (ii) and Proposition 1(iii) we can also prove the existence of an unbounded oscillatory solution for equation ( $\mathrm{E}_{+}$).

Corollary 1. If

$$
\begin{equation*}
I\left(a_{1}\right)=I\left(a_{2}\right)=I\left(a_{3}, a_{1}\right)=\infty \tag{13}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{0}^{\infty} a_{3}(t) \int_{0}^{t} a_{2}(s) \int_{0}^{s} a_{1}(u) \mathrm{d} u \mathrm{~d} s \mathrm{~d} t<\infty \tag{14}
\end{equation*}
$$

then there exists an unbounded oscillatory solution of $\left(\mathrm{E}_{+}\right)$.
Proof. If (14) holds, then, by Proposition 1(iii), there exists a Kneser solution tending to a nonzero constant. By Remark 5, (13) ensures that there exists an oscillatory solution. Assume that every oscillatory solution is bounded. Then ( $\mathrm{E}_{+}$) has a weakly oscillatory solution $y$, i.e., $y \neq 0, y^{\prime}$ is oscillatory, which contradicts [7, Corollary 2].

Remark 6. As already proved (Lemma 1), if $I\left(a_{1}\right)=I\left(a_{2}\right)=\infty$ and ( $\mathrm{E}_{+}$) has property A then this equation possesses an oscillatory solution. From (10), (14) and Theorem 3 we obtain the well-known fact that the equation $x^{\prime \prime \prime}+q(t) x=0$ is oscillatory if and only if it has property A (see, e.g., [16, Lemma 2.8]). But in general, the converse is not true. The following conclusion holds.

Corollary 2. Let $I\left(a_{1}\right)=I\left(a_{2}\right)=\infty$. If $\left(\mathrm{E}_{+}\right)$is oscillatory and has not property $A$, then ( $\mathrm{E}_{-}$) is nonoscillatory.

Proof. Since ( $E_{+}$) is oscillatory and has not property A, there exists a Kneser solution of ( $\mathrm{E}_{+}$) tending to a nonzero constant. By Proposition 1(iii) we have (14).

This condition ensures, in view of Proposition 1(iv), that the equation

$$
\left(\frac{1}{a_{1}(t)}\left(\frac{1}{a_{2}(t)} x^{\prime}(t)\right)^{\prime}\right)^{\prime}+a_{3}(t) x(t)=0
$$

is nonoscillatory. By Theorem 3 its adjoint equation, which is ( $\mathrm{E}_{-}$), is nonoscillatory, too.

Corollaries 1 and 2 are illustrated by the following example.
Example. Let $\varepsilon \in(0,1)$ and $T>1$. Consider the equation

$$
\begin{equation*}
\left(t \ln t\left(\frac{y^{\prime}(t)}{\ln t}\right)^{\prime}\right)^{\prime}+\frac{y(t)}{t^{2}(\ln t)^{1+\varepsilon}}=0, \quad t \in[T, \infty) \tag{15}
\end{equation*}
$$

Then

$$
\begin{gathered}
\int_{T}^{t} a_{1}(s) \mathrm{d} s=t(\ln t-1)+c_{1} \rightarrow \infty, \quad \int_{T}^{t} a_{2}(s) \mathrm{d} s=\ln \ln t+c_{2} \rightarrow \infty \text { for } t \rightarrow \infty, \\
\int_{T}^{t} a_{2}(s) \int_{T}^{s} a_{1}(s) \mathrm{d} s \mathrm{~d} t=t-\int_{T}^{t} \frac{\mathrm{~d} s}{\ln s}-T-\int_{T}^{t} \frac{c \mathrm{~d} s}{s \ln s} \rightarrow \infty \quad \text { for } t \rightarrow \infty, \\
\int_{T}^{\infty} a_{3}(t) \int_{T}^{t} a_{2}(s) \int_{T}^{s} a_{1}(u) \mathrm{d} u \mathrm{~d} s \mathrm{~d} t \leqslant \int_{T}^{\infty} \frac{\mathrm{d} t}{t(\ln t)^{1+\varepsilon}}=\int_{\ln T}^{\infty} \frac{\mathrm{d} u}{u^{1+\varepsilon}}<\infty, \\
\int_{T}^{\infty} a_{3}(t) \int_{T}^{t} a_{1}(s) \mathrm{d} s \mathrm{~d} t \geqslant \int_{T}^{\infty} \frac{\ln t-1}{t(\ln t)^{1+\varepsilon}} \mathrm{d} t=\int_{T}^{\infty} \frac{\mathrm{d} t}{t(\ln t)^{\varepsilon}} \mathrm{d} t=\int_{\ln T}^{\infty} \frac{\mathrm{d} s}{s^{\varepsilon}} \mathrm{d} s=\infty,
\end{gathered}
$$

i.e., (14) and (13) hold. Hence, by Corollary 1, equation (15) has an unbounded oscillatory solution and a Kneser solution tending to nonzero constant.

By Corollary 2 the equation

$$
\left(t \ln t\left(\frac{y^{\prime}(t)}{\ln t}\right)^{\prime}\right)^{\prime}-\frac{y(t)}{t^{2}(\ln t)^{1+\varepsilon}}=0
$$

is nonoscillatory.
Remark 7. Theorems 4 and 5 concern the cases $\left(K_{3}\right)$ and ( $K_{4}$ ). In case ( $K_{2}$ ), i.e., if $I\left(a_{1}\right)<\infty$ and $I\left(a_{1}, a_{2}\right)=\infty$, we obtain, applying Lemmas 1,3 and 5 to equation ( $\mathrm{E}_{+}^{\mathcal{K}_{2}}$ ), the same conclusion as that given in Theorem 4 (iii).

In the two remaining cases $\left(\mathrm{K}_{5}\right)$ and $\left(\mathrm{K}_{6}\right)$ we get, applying Lemmas 1,3 and 5 to the equation $\left(\mathrm{E}_{+}^{\mathcal{K}_{5}}\right)$ or $\left(\mathrm{E}_{+}^{\mathcal{K}_{6}}\right)$, the following conclusion:

Let one of the following two conditions hold:
(i) $I\left(a_{2}\right)=\infty, I\left(a_{1}, a_{2}\right)<\infty$ and

$$
\int_{0}^{\infty} a_{3}(t) \int_{t}^{\infty} a_{2}(s) \mathrm{d} s \int_{t}^{\infty} a_{1}(s) \int_{s}^{\infty} a_{2}(\tau) d \tau \mathrm{~d} s \mathrm{~d} t=\infty,
$$

(ii) $I\left(a_{1}\right)<\infty, I\left(a_{2}\right)<\infty$ and

$$
\int_{0}^{\infty} a_{3}(t) \int_{t}^{\infty} a_{1}(s) \mathrm{d} s \int_{t}^{\infty} \nu(s) \mathrm{d} s \int_{t}^{\infty} a_{2}(s) \int_{s}^{\infty} a_{1}(\tau) d \tau \mathrm{~d} s \mathrm{~d} t=\infty,
$$

where $\nu(t)=\left(\frac{1}{J_{t}^{\infty} a_{1}}\right)^{\prime} \int_{t}^{\infty} a_{2} \int_{s}^{\infty} a_{1}$.
Then ( $\mathrm{E}_{+}$), ( $\mathrm{E}_{+}^{\mathcal{C}}$ ) and ( $\mathrm{E}_{+}^{\mathcal{C C}}$ ) are oscillatory and, by Theorem 3, also the adjoint equations $\left(E_{+}^{\mathcal{A}}\right),\left(E_{+}^{\mathcal{C A}}\right)$ and $\left(E_{+}^{C \mathcal{C A}}\right)$ are oscillatory. In addition, every nonoscillatory solution of $\left(\mathrm{E}_{+}\right)$is either a Kneser solution or of class $N_{3}$ and tends to zero for $t \rightarrow \infty$.

Analogous statements can be obtained by application of Lemmas 1, 3 and 5 to equations ( $E_{+}^{\mathcal{C}}$ ) and ( $E_{+}^{\mathcal{C L}}$ ).

Coming back to the equations (2), (3), if we apply Theorem 4 (iii) to equation (2), Theorem 4 (ii) to equation $\left(2^{\mathcal{C}}\right)$ and Theorem 3 to the adjoint equations to (2) and (3), we obtain the following theorem.

Corollary 3. Let $q \neq 0, r>0$ and let one of the following conditions hold:
(i) $\int^{\infty} t / r(t) \mathrm{d} t=\infty$ and $I\left(\frac{1}{r},|q|\right)=\infty$,
(ii) $\int^{\infty} t|q(t)| \mathrm{d} t=\infty$ and $I\left(|q|, \frac{1}{r}\right)=\infty$.

Then equations (2), (3) are both oscillatory. In addition, if $q>0$, then these equations have property $A$.

The results presented in Parts I and II can be used for proving integral criteria for ( $\mathrm{E}_{-}$) to have property B as well as for the nonlinear equation $L(x)+f(x)=0$ to have property A or property B. This will be shown elsewhere.

| Case | $b_{0}(t)$ | $b_{1}(t)$ | $b_{2}(t)$ | $b_{3}(t)$ |
| :--- | :--- | :--- | :--- | :--- |
| $\left(\mathrm{K}_{1}\right)$ | 1 | $u_{1}$ | $u_{2}$ | 1 |
| $\left(\mathrm{~K}_{2}\right)$ | $\int_{t}^{\infty} u_{1}$ | $\left(\frac{1}{J_{t}^{\infty} u_{1}}\right)^{\prime}$ | $u_{2} \int_{t}^{\infty} u_{1}$ | 1 |
| $\left(\mathrm{~K}_{3}\right)$ | 1 | $u_{1} \int_{t}^{\infty} u_{2}$ | $\left(\frac{1}{J_{t}^{\infty} u_{2}}\right)^{\prime}$ | $\int_{t}^{\infty} u_{2}$ |
| $\left(\mathrm{~K}_{4}\right)$ | $\int_{t}^{\infty} u_{1} \int_{s}^{\infty} u_{2}$ | $\left(\frac{1}{J_{t}^{\infty} u_{1} \int_{s}^{\infty} u_{2}}\right)^{\prime}$ | $\left(\frac{1}{J_{t}^{\infty} u_{2}}\right)^{\prime} \int_{t}^{\infty} u_{1} \int_{s}^{\infty} u_{2}$ | $\int_{t}^{\infty} u_{2}$ |
| $\left(\mathrm{~K}_{5}\right)$ | $\int_{t}^{\infty} u_{1}$ | $\left(\frac{1}{J_{t}^{\infty} u_{1}}\right)^{\prime} \int_{t}^{\infty} u_{2} \int_{s}^{\infty} u_{1}$ | $\left(\frac{1}{J_{t}^{\infty} u_{2} \int_{s}^{\infty} u_{1}}\right)^{\prime}$ | $\int_{t}^{\infty} u_{2} \int_{s}^{\infty} u_{1}$ |
| $\left(\mathrm{~K}_{6}\right)$ | $\int_{t}^{\infty} u_{1} \int_{t}^{\infty} \nu$ | $\left(\frac{1}{J_{t}^{\infty} \nu}\right)^{\prime}$ | $\int_{t}^{\infty} \nu\left(\frac{1}{\left.J_{t}^{\infty} u_{2} \int_{s}^{\infty} u_{1}\right)}\right)^{\prime}$ | $\int_{t}^{\infty} u_{2} \int_{s}^{\infty} u_{1}$ |

Table 1: Canonical forms for the operator $D$ where $\nu(t)=\left(\frac{1}{f_{t}^{\infty} u_{1}}\right)^{\prime} \int_{t}^{\infty} u_{2} \int_{s}^{\infty} u_{1}$.

| Case | $b_{0}(t)$ | $b_{1}(t)$ | $b_{2}(t)$ | $b_{3}(t)$ |
| :--- | :--- | :--- | :--- | :--- |
| $\left(\mathrm{K}_{1}\right)$ | 1 | $a_{2}$ | $a_{3}$ | 1 |
| $\left(\mathrm{~K}_{2}\right)$ | $\int_{t}^{\infty} a_{2}$ | $\left(\frac{1}{J_{t}^{\infty} a_{2}}\right)^{\prime}$ | $a_{3} \int_{t}^{\infty} a_{2}$ | 1 |
| $\left(\mathrm{~K}_{3}\right)$ | 1 | $a_{2} \int_{t}^{\infty} a_{3}$ | $\left(\frac{1}{J_{t}^{\infty} a_{3}}\right)^{\prime}$ | $\int_{t}^{\infty} a_{3}$ |
| $\left(\mathrm{~K}_{4}\right)$ | $\int_{t}^{\infty} a_{2} \int_{s}^{\infty} a_{3}$ | $\left(\frac{1}{J_{t}^{\infty} a_{2} \int_{s}^{\infty} a_{3}}\right)^{\prime} \cdot$ | $\left(\frac{1}{J_{t}^{\infty} a_{3}}\right)^{\prime} \int_{t}^{\infty} a_{2} \int_{s}^{\infty} a_{3}$ | $\int_{t}^{\infty} a_{3}$ |
| $\left(\mathrm{~K}_{5}\right)$ | $\int_{t}^{\infty} a_{2}$ | $\left(\frac{1}{J_{t}^{\infty} a_{2}}\right)^{\prime} \int_{t}^{\infty} a_{3} \int_{s}^{\infty} a_{2}$ | $\left(\frac{1}{J_{t}^{\infty} a_{3} \int_{t}^{\infty} a_{2}}\right)^{\prime}$ | $\int_{t}^{\infty} a_{3} \int_{s}^{\infty} a_{2}$ |
| $\left(\mathrm{~K}_{6}\right)$ | $\int_{t}^{\infty} a_{2} \int_{t}^{\infty} \mu$ | $\left(\frac{1}{J_{t}^{\infty} \mu}\right)^{\prime}$ | $\int_{t}^{\infty} \mu\left(\frac{1}{\left.J_{t}^{\infty} a_{3} J_{s}^{\infty} a_{2}\right)}\right)^{\prime}$ | $\int_{t}^{\infty} a_{3} \int_{s}^{\infty} a_{2}$ |

Table 2: Canonical forms for the operator $L^{\mathcal{C}}$ where $\mu(t)=\left(\frac{1}{J_{t}^{\infty} a_{2}}\right)^{\prime} \int_{t}^{\infty} a_{3} \int_{s}^{\infty} a_{2}$.

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Authors' addresses: M. Cecchi, M. Marini, Depart. of Electr. Eng., University of Florence, Via S. Marta 3, 50139 Firenze, Italy; Z. Došlá, Depart. of Mathematics, Masaryk University, Janáčkovo nám. 2a, 66295 Brno, Czech Republic.

