On Oscillation Properties of Self-Adjoint Boundary Value Problems of Fourth Order

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Abstract—The connection between the number of internal zeros of nontrivial solutions to fourth-order self-adjoint boundary value problems and the inertia index of these problems is studied. We specify the types of problems for which such a connection can be established. In addition, we specify the types of problems for which a connection between the inertia index and the number of internal zeros of the derivatives of nontrivial solutions can be established. Examples demonstrating the effectiveness of the proposed new approach to an oscillatory problem are considered.

Keywords: boundary value problems for ordinary differential equations, spectral and oscillatory problems, inertia index, Kellogg kernels

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1. In this paper, we present some results and methods that make it possible to keep track of the relationship between the number of internal zeros of the solution to a self-adjoint boundary value problem of fourth order and the (negative) index of inertia of this problem. Classical results of this type for the Sturm–Liouville problem are well known (see, e.g., [1]), but, for a fourth-order problem, the situation is much more complicated.

Consider the boundary value problem

\begin{align}
(p y^{(iv)})'' - (q y')' + ry &= 0, \\
B y'' + C y' &= 0,
\end{align}

where $B$ and $C$ are block diagonal real matrices of order 4:

$$
B = \begin{cases} 
B_0 & 0 \\
0 & B_1
\end{cases}, \quad C = \begin{cases} 
C_0 & 0 \\
0 & C_1
\end{cases},
$$

$$
B_k, C_k \in \mathbb{R}^{2 \times 2}, \quad k = 0, 1,
$$

which satisfy the equation

$$
B^{-1}\text{im} C = \mathbb{R}^4 \ominus \ker C,
$$

where the symbol $B^{-1}$ denotes the total preimage. Here, we also used the auxiliary notation

$$
y^\ast \cong \begin{pmatrix} y(0) \\
y'(0) \\
y(1) \\
y'(1)
\end{pmatrix}, \quad y^\ast \cong \begin{pmatrix} y^[1](0) \\
y^[2](0) \\
y^[3](1) \\
y^[4](1)
\end{pmatrix},
$$

$$
y^[1] \cong p y'', \quad y^[3] \cong -(y^[1])' + q y'.
$$

In the study of problem (1), (2), three basic degrees of generality are possible. The first corresponds to the case of $p \in C^0[0, 1], q \in C^r[0, 1], r \in C[0, 1]$; moreover, the function $p$ is assumed to be positive. Here, Eq. (1) is understood directly and its solution $y$ is sought among the functions from the Sobolev space $W_2^4[0, 1]$. Fourth-order problems are still often considered in this formulation.

The second degree of generality corresponds to the case when $p, p^{-1} \in L_\infty[0, 1]$ and $q, r \in L[0, 1]$; moreover, the function $p$ is again assumed to be positive. Here, Eq. (1) is understood as formal notation for the equation $-(y^[3])' + ry = 0$, whose solutions are sought in the class of functions satisfying the conditions

$$
y, y', y^[2], y^[3] \in AC[0, 1] = W_1^1[0, 1]
$$

This treatment is also classical and can be found in [2, Section 15].

Finally, the third degree of generality occurs in the situation when the functions $p, p^{-1} \in L_\infty[0, 1]$ are posi-
tive, while the other coefficients are generalized functions (distributions): \( q \in W^{-2}_{2,1}[0,1] \) and \( r \in W^{-2}_{2,1}[0,1] \).

Here, the symbol \( W^{k}_{2,1}[0,1] \), where \( k = 1, 2 \), denotes the space of bounded linear functionals acting on the Sobolev space \( W^{k}_{2,1}[0,1] \). In this case, following [3], we introduce the space

\[
W^{2}_{2,1,1}[0,1] \ni \{ y \in W^{2}_{2,1}[0,1] : By \in \text{im}C \}
\]

and the solution of problem (1), (2) is understood as a function \( y \in W^{2}_{2,1,1}[0,1] \) satisfying the identity

\[
(\forall z \in W^{2}_{2,1,1}[0,1]) \quad \frac{1}{2} \int_{0}^{\infty} py''z''dx + \langle q, y'z' \rangle + \langle r, yz \rangle + \langle \xi, z \rangle_{R} = 0. \tag{7}
\]

Here, \( \xi \in \mathbb{R}^{4} \) is an arbitrary fixed vector with the property \( C\xi = -By' \), where the column \( y' \) is defined according to (5). Note that the terms in (7) are well defined, since \( y', y'' \in W^{2}_{2,1}[0,1] \) for all \( y, z \in W^{2}_{2,1,1}[0,1] \). It can be shown that this treatment coincides with the above-indicated classical definitions if \( p, q, \) and \( r \) are sufficiently smooth functions.

In terms of operator theory, what was said above means that problem (1), (2) is associated with an operator \( T : W^{2}_{2,1,1}[0,1] \rightarrow W^{2}_{2,1,1}[0,1] \) that maps every function \( y \in W^{2}_{2,1,1}[0,1] \) to the bounded linear functional \( Ty \in W^{2}_{2,1,1}[0,1] \), acting on an arbitrary function \( z \in W^{2}_{2,1,1}[0,1] \) according to the rule

\[
\langle Ty, z \rangle := \frac{1}{2} \int_{0}^{\infty} py''z''dx + \langle q, y'z' \rangle + \langle r, yz \rangle + \langle \xi, z \rangle_{R}. \tag{8}
\]

Here, as before, \( C\xi = -By' \); moreover, by condition (4), the action of the operator \( T \) does not depend on the choice of the vector \( \xi \). The solutions of problem (1), (2) are exactly the elements of the kernel of the operator \( T \).

Note that the indicated general understanding of fourth-order boundary value problems goes back to the treatments of second-order boundary value problems proposed in [4, 5]. In what follows, problem (1), (2) is understood in the most general (third) sense. It follows from what was said above that all results are valid for narrower classical formulations.

Hereafter, we assume that the considered spaces are real and the operator \( T \) defined by (8) is symmetric, i.e., it satisfies the relation \( \langle Ty, z \rangle = \langle Tz, y \rangle \) for all \( y, z \in W^{2}_{2,1,1}[0,1] \). The operator \( T \) is called positive (or nonnegative) if, for any function \( y \in W^{2}_{2,1,1}[0,1] \) that does not vanish identically, we have \( \langle Ty, y \rangle > 0 \) (or \( \langle Ty, y \rangle \geq 0 \), respectively). The index of inertia of the operator \( T \) is defined as usual (see, e.g., [6]) as the maximum of the dimensions of subspaces \( \mathcal{M} \subset W^{2}_{2,1,1}[0,1] \) satisfying the condition

\[
(\forall y \in \mathcal{M}) \quad \langle Ty, y \rangle < 0.
\]

2. In this section, we consider an auxiliary problem corresponding to the case \( q \equiv 0 \) and \( r \leq 0 \). For a singular function \( r \in W^{-2}_{2,1}[0,1] \), the last inequality means that \( \langle r, y \rangle \leq 0 \) for any nonnegative function \( y \in W^{2}_{2,1}[0,1] \). In this case, it is easy to show that there exists a nondecreasing function \( R \in L_{\infty}[0,1] \) obeying the identity

\[
(\forall y \in W^{2}_{2,1}[0,1]) \quad \langle r, y \rangle = \int_{0}^{\infty} Ry' dx - \langle Ry \rangle_{R}'.
\]

We distinguish a class of boundary conditions (2) characterized by the following two assumptions.

**Assumption A.** There are no real numbers \( a_{0} > 0, a_{1} < 0, a_{2} > 0 \), and \( a_{3} > 0 \) for which

\[
B_{0} \begin{pmatrix} a_{0} \\ a_{1} \\ a_{2} \end{pmatrix} + C_{0} \begin{pmatrix} a_{3} \\ \end{pmatrix} = 0.
\]

**Assumption B.** There are no real numbers \( b_{0} > 0, b_{1} > 0, b_{2} > 0 \), and \( b_{3} < 0 \) for which

\[
B_{1} \begin{pmatrix} b_{0} \\ b_{1} \\ b_{2} \end{pmatrix} - C_{1} \begin{pmatrix} b_{3} \\ \end{pmatrix} = 0.
\]

**Theorem 1.** Suppose that \( q \equiv 0, r \leq 0, \) and the corresponding function \( R \in L_{\infty}[0,1] \) is not constant in neighborhoods of the points 0 and 1. Assume that boundary conditions (2) satisfy Assumptions A and B. Then the solution space of the problem has a dimension of at most 1. Moreover, for any nontrivial solution \( y \in W^{2}_{2,1,1}[0,1] \), there is no point \( x \in (0,1) \) satisfying the pair of equalities \( y(x) = y'(x) = 0 \) and the pair of equalities \( y'(x) = (py'')(x) = 0 \).

3. One of the basic concepts in the general theory of oscillation is the Kellogg kernel (see, e.g., [7, Chapter IV, Section 3, Theorem 1] or [8]). In the case of positive definite problems of fourth order with separated boundary conditions, it is known [9] that the Green’s function \( K(x,s) := \langle \delta_{x}, T^{-1} \delta_{s} \rangle \) of the corresponding operator is a Kellogg kernel if and only if it is positive on the open square \( (0,1) \times (0,1) \). Using this fact and applying general results concerning integral equations with Kellogg kernels [7, Chapter IV, Section 3, Theorems 1, 2], we obtain the following assertion.

**Theorem 2.** Let the operator \( T \) be defined by (8) and \( H : W^{2}_{2,1,1}[0,1] \rightarrow W^{-2}_{2,1,1}[0,1] \) be the operator of multiplication by a positive generalized function \( h \in W^{-2}_{2,1}[0,1] \).
Suppose that the operator \( T + H : W^{2}_{2,B,C}[0, 1] \rightarrow W^{2}_{2,B,C}[0, 1] \) is positive and satisfies the estimate \( \langle \delta_s, (T + H)^{-1} \delta_t \rangle > 0 \) for all points \( x, s \in (0, 1) \). Then the dimension of the kernel of \( T \) is at most 1 and the number of sign changes for any nontrivial function from this kernel equals the index of inertia of \( T \).

In the general situation, checking whether the Green’s function of a fourth-order differential operator is positive represents a nontrivial task. In many cases, however, it suffices to answer this question when \( q \equiv 0 \) and \( r \equiv 0 \).

**Theorem 3.** Suppose that \( q \equiv 0 \) and \( r \equiv 0 \), boundary conditions (2) obey Assumptions A and B, and the corresponding operator \( T : W^{2}_{2,B,C}[0, 1] \rightarrow W^{2}_{2,B,C}[0, 1] \) is positive. Then the Green’s function \( K(x, s) \equiv \langle \delta_s, T^{-1} \delta_t \rangle \) is positive on the square \((0, 1) \times (0, 1)\).

4. If the function \( h \in W^{2}_{2}[0, 1] \) is not positive, Theorem 2 generally does not hold. For this reason, in a number of particular problems, it is useful to combine Theorem 1 with more general results (than Theorem 2) concerning operators that do not lower the number of sign changes for functions. These results are stated in the present section.

Consider the following situation. Let \( H : W^{2}_{2,B,C}[0, 1] \rightarrow W^{2}_{2,B,C}[0, 1] \) be a nonnegative completely continuous operator for which the corresponding operator \( T + H \) is positive. The sequence of the eigenvalues \(-1 < \lambda_0 \leq \lambda_1 \leq \lambda_2 \leq \ldots\) (listed with multiplicity) of the linear pencil \( \lambda \mapsto T - \lambda H \) is associated with a sequence \( \{f_k\}_{k=0}^{\infty} \) of its eigenfunctions. Let \( I : \mathcal{S} \rightarrow W^{2}_{2,B,C}[0, 1] \) be a bounded and boundedly invertible operator, where \( \mathcal{S} \) is a Hilbert space continuously embedded in \( C[0, 1] \). The following result holds.

**Theorem 4.** Suppose that the operator \( I^{-1}(T + H)^{-1} \), \( HI : \mathcal{S} \rightarrow \mathcal{S} \), does not increase the number of sign changes for any function and the following assumptions are valid:

1. All nonnegative eigenvalues of the pencil \( \lambda \mapsto T - \lambda H \) are simple.

2. For any \( m \in \mathbb{N} \) with property \( \lambda_m \geq 0 \), inside the linear span of the set \( \{f_k\}_{k=0}^m \), there exists a neighborhood \( \mathcal{U} \) of the function \( f_m \) such that, for any function \( y \in \mathcal{U} \), the number of sign changes for the function \( I^{-1}y \) does not exceed its counterpart for \( I^{-1}f_m \).

A generalized function \( h \in W^{2}_{2}[0, 1] \) is said to be positive if \( \langle h, y \rangle > 0 \) for any nonnegative test function \( y \in W^{2}_{2}[0, 1] \) that does not vanish identically.

(3) For \( m \in \mathbb{N} \) with property \( \lambda_m = 0 \), there exists a subspace \( \mathcal{M} \subset W^{2}_{2,B,C}[0, 1] \) of dimension \( m + 1 \) having a trivial intersection with the kernel of \( H \) such that, for any function \( y \in \mathcal{M} \), the corresponding function \( I^{-1}y \) has at most \( m \) sign changes.

Then, for any nontrivial function \( y \in \ker T \), the number of sign changes for the corresponding function \( I^{-1}y \in \mathcal{S} \) equals the index of inertia of \( T \).

In addition to the trivial case, when the operator \( I : W^{2}_{2,B,C}[0, 1] \rightarrow W^{2}_{2,B,C}[0, 1] \) is the identity mapping, the following two situations are typical:

1. The operator \( I \) performs the change of variable \( [Iu](x) = u(\tau(x)) \), where the function \( \tau \in W^{2}_{2}[0, 1] \) has an everywhere positive derivative. The role of \( \mathcal{S} \) is played by a subspace of the Sobolev space \( W^{2}_{2}[0, 1] \).

2. The space \( W^{2}_{2,B,C}[0, 1] \) does not contain nontrivial constant functions, and \( I^{-1} \) is the differentiation operator. The role of \( \mathcal{S} \) is played by a subspace of the Sobolev space \( W^{2}_{2}[0, 1] \).

The first of these situations occurs, for example, in the standard procedure of eliminating the second term of the left-hand side of Eq. (1) (see, e.g., [6, 10–12]). Specifically, suppose that there is a function \( \tau \in W^{2}_{2}[0, 1] \) with an everywhere positive derivative such that \( \tau(0) = 0, \tau(1) = 1 \), and

\[
(\forall y \in W^{2}_{2}[0, 1])
\int_{0}^{1} p\tau'' y' dx + \langle q, \tau' y \rangle = \gamma_0 \cdot y(0) + \gamma_1 \cdot y(1),
\]

where \( \gamma_0, \gamma_1 \in \mathbb{R} \) are constants. The operator \( I : \mathcal{S} \rightarrow W^{2}_{2,B,C}[0, 1] \) is defined according to the rule \( [Iu](x) = u(\tau(x)) \). Straightforward calculations show that the operator \( I*TI : \mathcal{S} \rightarrow \mathcal{S}^* \) is then associated with the boundary value problem

\[
(\tilde{p}y'')'' + \tilde{r}y = 0, \quad \tilde{B}y' + \tilde{C}y = 0,
\]

where \( \tilde{p} \circ \tau = p \circ (\tau)^3, \tilde{r} \circ \tau = r \circ (\tau)^{-1} \), and

\[
\tilde{B} = B \cdot \text{diag}[1, \tau'(0), 1, \tau'(1)] - C \cdot \text{diag}[0, \gamma_0, 0, \gamma_1],
\]

\[
\tilde{C} = C \cdot \text{diag}[1, (\tau'(0))^{-1}, 1, (\tau'(1))^{-1}].
\]

In this case, the space \( \mathcal{S} \) is specified as \( \mathcal{S} = W^{2}_{2,B,C}[0, 1] \subseteq W^{2}_{2}[0, 1] \).

As applied to the second of the above-indicated situations, the following result is useful (cf. [12]):

**Theorem 5.** Assume that the condition \( y(0) = 0 \) holds for all functions \( y \in W^{2}_{2,B,C}[0, 1] \). Let \( I^{-1} : W^{2}_{2,B,C}[0, 1] \rightarrow \mathcal{S} \)
be the differentiation operator and the action of the operator $H: W^2_{2,β,γ} [0, 1] → W^2_{2,β,γ} [0, 1]$ be defined by the rule

$$\langle Hy, z \rangle = \langle h, y'z' \rangle + \alpha_0 \cdot (y'z')(0) + \alpha_1 \cdot (y'z')(1),$$

where $\alpha_0, \alpha_1 ≥ 0$ and $h \in W^2_{2,β,γ} [0, 1]$ is a nonnegative generalized function. Suppose that $J^*(T + H)\mathbb{S} → \mathbb{S}$ is a positive definite Sturm–Liouville operator. Then the operator $I^{-1}(T + H)^{-1} H: \mathbb{S} → \mathbb{S}$ does not increase the number of sign changes for any function.

As another, more complicated example, we consider the spectral problem

$$y^{(4)} − (q y')' − \lambda y = 0,$$

$$\begin{bmatrix}
1000 \\
0100 \\
0000 \\
0000
\end{bmatrix} y^\wedge = \begin{bmatrix}
000\tilde{a}\lambda + b \\
0000 \\
0000 \\
0000
\end{bmatrix} y^\gamma = 0,
$$

which was studied in [14], assuming that the function $q ∈ C^2[0, 1]$ is nonnegative, $ω \in \left[\frac{π}{2}, π\right]$, and $σ = bc − ad > 0$. Applying the change-of-variable transformation from Section 4, we obtain the problem

$$\begin{bmatrix}
1000 \\
0100 \\
0000 \\
0000
\end{bmatrix} y^\wedge = \begin{bmatrix}
000\sigma \lambda \cdot (\tilde{a}\lambda + b) + \eta \cdot (c\lambda + d) \\
0000 \\
0000 \\
0000
\end{bmatrix} y^\gamma = 0$$

with coefficients $\theta > 0$ and $\eta ≥ 0$. In this case, Theorem 1 automatically guarantees that any positive eigenvalue of the problem is simple. Now we set

$$\hat{h} = \begin{cases}
\frac{\lambda}{\delta_1} & \text{if } ω = π \\
\lambda \delta - (\cot \omega) \cdot \delta_1 & \text{if } ω ≠ π,
\end{cases}$$

where $\delta_1$ is the delta function supported at 1. If $c\lambda + d = 0$ or $(\tilde{a}\lambda + b) \cdot (c\lambda + d)^{-1} ≤ -\eta \delta^{-1}$, Theorems 2 and 3 immediately imply that the number of zeros of an eigenfunction on the interval $(0, 1)$ equals the index of inertia of the problem.

If $(\tilde{a}\lambda + b) \cdot (c\lambda + d)^{-1} < -\eta \delta^{-1}$, it is possible that the resulting operator $\hat{H} + \hat{h}$ is not positive, so the indicated procedure cannot be applied. Accordingly, in this case, we should set

$$\langle \hat{H}y, z \rangle = \begin{cases}
\langle y', z' \rangle + \eta + \theta \cdot \frac{\tilde{a}\lambda + b}{c\lambda + d} & \text{if } ω = π \\
\langle y', z' \rangle & \text{if } ω ≠ π.
\end{cases}$$
For \( \omega = \pi \), this operator \( \hat{H} \) does not increase the number of sign changes for any function, so we can apply Theorem 4. For \( \omega \neq \pi \), the operator \( \hat{I} : \hat{S} \rightarrow W^2_{2, \pi, \omega}([0, 1]) \) should be specified as an integration operator, and then Theorems 5 and 4 can be applied. In this case, the index of inertia of the problem is equal to the number of internal zeros of the eigenfunction derivative \( y' \in W^2_{2}[0,1] \). The number of zeros of the eigenfunction itself is either equal to the index of inertia or is fewer by 1, depending on the sign of \( y(1)y'(1) \).

It should be noted that Theorem 2.2 in [14], which presents results concerning the considered example, involves inaccuracies. Namely, according to this theorem, in the case \( a/c \geq 0 \), the number of positive eigenvalues for which the number of zeros of corresponding eigenfunctions is different from the index of inertia cannot exceed 2. Actually, this is not true. For example, the spectral problem

\[
\begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 \\ 0 & 0 & 0 & 0 \end{pmatrix} y^\beta + \begin{pmatrix} 0 & 0 & 0 & 0 \\ 0 & 0 & -1 & 0 \\ 0 & 0 & 0 & -\lambda + 10^6 \\ 0 & 0 & 0 & 0 \end{pmatrix} y^\gamma = 0
\]

has ten such eigenvalues to the left of the point \( \lambda = 10^6 \). This is verified directly by analyzing the behavior of the solution

\[
y(x, \lambda) = (\sinh \sqrt{\lambda} - \sin \sqrt{\lambda})(\sinh \sqrt{\lambda}x - \sin \sqrt{\lambda}x) - (\cosh \sqrt{\lambda} + \cos \sqrt{\lambda})(\cosh \sqrt{\lambda}x - \cos \sqrt{\lambda}x)
\]

to the boundary value problem

\[
y^{(4)} - \lambda y = 0, \quad y(0) = y'(0) = y^{''}(1) = 0,
\]

depending on the parameter \( \lambda \in (0, 10^6) \).

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