# ON SMALL VIBRATIONS OF A DAMPED STIELTJES STRING 

Olga Boyko and Vyacheslav Pivovarchik<br>Communicated by Alexander Gomilko


#### Abstract

Inverse problem of recovering masses, coefficients of damping and lengths of the intervals between the masses using two spectra of boundary value problems and the total length of the Stieltjes string (an elastic thread bearing point masses) is considered. For the case of point-wise damping at the first counting from the right end mass the problem of recovering the masses, the damping coefficient and the lengths of the subintervals by one spectrum and the total length of the string is solved.


Keywords: damping, Dirichlet boundary condition, point mass, Hermite-Biehler polynomial, continued fraction, eigenvalues.

Mathematics Subject Classification: 35Q99, 39A99.

## 1. INTRODUCTION

Investigation of spectral problems for damped systems was started in [8]. As far as we know first results on inverse problems for damped strings were obtain in [2] and $[15,16]$. In these papers the left end of the string was supposed to be free and the right end damped or in other words the right end of the string could move with viscous friction in the direction orthogonal to the equilibrium position of the string. In these papers very wide classes of strings were considered. In [2] the class of so-called regular strings was used, i.e. strings of finite mass and length, while in $[15,16]$ the class of so-called S-strings, i.e. strings of finite lengths and finite first momentum of mass distribution. Conditions necessary and sufficient for a sequence of complex numbers to be the spectrum of a damped string were given in [2] in implicit form and in [16] and [15] explicitly. If the string is smooth such that $\rho \in W_{2}^{2}(0, l)$ and $\rho(s) \geq \epsilon>0$, where $\rho$ is the density of the string and $l$ is its length, then one can apply the Liouville transformation ([5, p. 202]) to reduce the equation of the string to the Sturm-Liouville equation. The corresponding boundary value problem with the spectral parameter in the boundary conditions was considered in many publications (see [7, 21, 22, 27] and
references therein). In [10] and [23] conditions on a sequence of complex numbers were given necessary and sufficient to be the spectrum of a smooth inhomogeneous string damped at one end for any value of the damping parameter except for one crucial value.

The opposite case of extremely nonsmooth, so-called Stieltjes string (i.e. a thread bearing point masses) was considered in [9] and [14] assuming absence of damping. The inverse problem, i.e. the problem of recovering the parameters of the string by the spectrum of its vibrations and by the total length of the string, for Stieltjes strings of finite number of masses with point-wise (i.e. one dimensional) damping at the right end and the left end free was solved in [2]. Also this problem can be reduced to the problem of damped oscillators considered in [29,30]. Another approach to the inverse problem for damped finite dimensional systems was developed in [17] where the given data included not only eigenvalues but also the so-called Jordan pairs. An inverse problem for a Stieltjes string with damping at the midpoint was solved in [4]. A nice review can be found in [6]. Inverse problems generated by the Stieltjes string recurrence relations on graph domains were considered in [24, 25] and [18].

The case of distributed damping is not investigated in detail. We should mention [31], where it was proved that two spectra of Dirichlet-Dirichlet boundary value problem (the problem with the Dirichlet boundary conditions at both ends) and Dirichlet-Neumann boundary problem (the problem with the Dirichlet boundary condition at the left end and the Neumann boundary condition at the right end) uniquely determine the density and the stiffness of the string if the damping is constant but the problem with constant damping considered in this paper can be reduced by a change of the spectral parameter to the case of an undamped string. In $[1,12]$ the inverse problems for a damped string were considered on the semi-axis and on the axis, respectively.

We consider the case of a Stieltjes string with finite number of point masses. In Section 2 we describe the spectra of a Stieltjes string small transversal vibrations with the both ends fixed (Dirichlet-Dirichlet boundary value problem) and with the left end fixed and the right end free (Dirichlet-Neumann boundary value problem). In Sections 2 and 3 we consider the case where only the mass neighboring the right end is damped. The Dirichlet-Neumann boundary value problem in this case can be reduced to the problem considered in [29,30] and [3].

In Section 3 we solve the inverse Dirichlet-Dirichlet problem, i.e. the problem of recovering the values of masses, of subintervals and of the damping coefficient of the damped mass by given spectrum, the total length of the string and the length of the last subinterval at the right end.

It is well known that in the undamped case the eigenvalues of these two problems interlace. In Section 4 the analogues of the interlacing conditions for the case of one dimensional damping at the mass neighboring the right end are found. They appear to be the set of equalities and inequalities involving the Dirichlet-Dirichlet and Dirichlet-Neumann eigenvalues. In Section 5 the problem of vibrations of the Stieltjes string is considered for the case where all the masses are damped. The aim is to recover not only the values of the masses and the lengths of the subintervals but the coefficients of damping also. As the given data the spectra of the

Dirichlet-Dirichlet and Dirichlet-Neumann problems are used together with the total length of the string. The conditions on two sequences of complex numbers to be the spectra of the Dirichlet-Dirichlet and Dirichlet-Neumann boundary value problems are established. These conditions are given in implicit form: the ratio of the corresponding characteristic polynomials should admit decomposition in continued fraction of a special type. The decomposition is an analogue to that in [9] based on the results of [28].

## 2. DIRECT PROBLEM FOR A DAMPED STIELTJES STRING

Like in [9] we suppose the string to be a thread (i.e. a string of zero density) bearing a finite number of point masses. Let $l_{k}(k=0,1, \ldots, n)$ be the lengths of the intervals of zero density and let $m_{k}(k=1,2, \ldots, n)$ be the values of the masses separating the intervals ( $l_{k}$ lies between $m_{k}$ and $m_{k+1}$ ), the last mass has only one thread at the left and $m_{k}>0$ for $k=1,2, \ldots, n$. Let us denote $\alpha_{k} \geq 0$ the coefficient of damping (viscous friction) of the point mass $m_{k}$. Denote by $v_{k}(t)$ the transversal displacements of the point masses at the time $t$.

We impose a Dirichlet boundary condition on the left end, that is, the left end is fixed and consider the two cases: 1) the right end is fixed (Dirichlet-Dirichlet problem) and 2) right end is free to move in the direction orthogonal to the equilibrium position of the string (Dirichlet-Neumann problem).

We assume the thread to be stretched by the stretching force which is equal to 1 . Taking into account that on the intervals of zero density the general solution of the differential equation is a linear function of $s$ multiplied by a function of $t$ we obtain

$$
\begin{equation*}
\frac{v_{k}(t)-v_{k+1}(t)}{l_{k}}+\frac{v_{k}(t)-v_{k-1}(t)}{l_{k-1}}+m_{k} v_{k}^{\prime \prime}(t)+\alpha_{k} v_{k}^{\prime}(t)=0 \quad(k=1,2, \ldots, n) \tag{2.1}
\end{equation*}
$$

In this section we assume $\alpha_{k}=0$ for $k=1,2, \ldots, n-1$ and $\alpha_{n}=\alpha>0$. Substituting $v_{k}(t)=u_{k} e^{i \lambda t}$ we obtain

$$
\begin{align*}
\frac{u_{k}-u_{k+1}}{l_{k}}+\frac{u_{k}-u_{k-1}}{l_{k-1}}-m_{k} \lambda^{2} u_{k} & =0 \quad(k=1,2, \ldots, n-1),  \tag{2.2}\\
\frac{u_{n}-u_{n+1}}{l_{n}}+\frac{u_{n}-u_{n-1}}{l_{n-1}}-m_{n} \lambda^{2} u_{n}+i \alpha \lambda u_{n} & =0 \tag{2.3}
\end{align*}
$$

For the Dirichlet-Dirichlet problem we suppose the ends of the thread to be fixed, i.e. $v_{0}(t)=v_{n+1}(t)=0$, what means that

$$
\begin{gather*}
u_{0}=0  \tag{2.4}\\
u_{n+1}=0 \tag{2.5}
\end{gather*}
$$

For the Dirichlet-Neumann problem we assume that the right end is free and bears no point mass. Then the boundary condition at the right looks as follows:

$$
\begin{equation*}
u_{n+1}=u_{n} \tag{2.6}
\end{equation*}
$$

According to [9] we have

$$
\begin{equation*}
u_{k}=R_{2 k-2}\left(\lambda^{2}\right) u_{1} \quad(k=1,2, \ldots, n-1), \tag{2.7}
\end{equation*}
$$

where $R_{2 k-2}\left(\lambda^{2}\right)$ is a polynomial of degree $2 k-2$ obtained by (2.2) and by definition

$$
R_{2 k-1}\left(\lambda^{2}\right)=\frac{R_{2 k}\left(\lambda^{2}\right)-R_{2 k-2}\left(\lambda^{2}\right)}{l_{k}} .
$$

Due to (2.2) the polynomials $R_{k}$ satisfy the recurrence conditions:

$$
\begin{align*}
R_{2 k-1}\left(\lambda^{2}\right) & =-\lambda^{2} m_{k} R_{2 k-2}\left(\lambda^{2}\right)+R_{2 k-3}\left(\lambda^{2}\right)  \tag{2.8}\\
R_{2 k}\left(\lambda^{2}\right) & =l_{k} R_{2 k-1}\left(\lambda^{2}\right)+R_{2 k-2}\left(\lambda^{2}\right)  \tag{2.9}\\
& \left(k=1,2, \ldots, n-1, \quad R_{-1}\left(\lambda^{2}\right)=\frac{1}{l_{0}}, \quad R_{0}\left(\lambda^{2}\right)=1\right) .
\end{align*}
$$

Using (2.7)-(2.9) and taking into account the boundary condition (2.5) we rewrite (2.3) as follows:

$$
\begin{equation*}
\phi(\lambda):=R_{2 n-3}\left(\lambda^{2}\right)+\left(-m_{n} \lambda^{2}+i \lambda \alpha+l_{n}^{-1}\right) R_{2 n-2}\left(\lambda^{2}\right)=0 \tag{2.10}
\end{equation*}
$$

The spectrum $\left\{\nu_{k}\right\}(k= \pm 1, \pm 2, \ldots, \pm n)$ of problem (2.2)-(2.5) coincides with the set of zeros of $\phi(\lambda)$. We are also interested in problem (2.2)-(2.4), (2.6). Again using (2.7)-(2.9) and taking into account boundary condition (2.6) we obtain from (2.3):

$$
\begin{equation*}
\psi(\lambda):=R_{2 n-3}\left(\lambda^{2}\right)+\left(-m_{n} \lambda^{2}+i \lambda \alpha\right) R_{2 n-2}\left(\lambda^{2}\right)=0 . \tag{2.11}
\end{equation*}
$$

The spectrum $\left\{\mu_{k}\right\}(k= \pm 1, \pm 2, \ldots, \pm n)$ of problem (2.2)-(2.4), (2.6) coincides with the set of zeros of $\psi(\lambda)$.

We call $\phi(\lambda)$ and $\psi(\lambda)$ characteristic polynomials of problems (2.2)-(2.5) and (2.2)-(2.4), (2.6), respectively.

It is known [9] that

$$
\begin{equation*}
\frac{R_{2 n-2}\left(\lambda^{2}\right)}{R_{2 n-3}\left(\lambda^{2}\right)}=l_{n-1}+\frac{1}{-m_{n-1} \lambda^{2}+\frac{1}{l_{n-2}+\frac{1}{-m_{n-2} \lambda^{2}+\ldots+\frac{1}{l_{1}+\frac{1}{-m_{1} \lambda^{2}+\frac{1}{l_{0}}}}}}} . \tag{2.12}
\end{equation*}
$$

Definition 2.1. A function $\omega(\lambda)$ is said to be a Nevanlinna function (or R-function in terms of [13]) if:

1) it is analytic in the half-planes $\operatorname{Im} \lambda>0$ and $\operatorname{Im} \lambda<0$,
2) $\omega(\bar{\lambda})=\overline{\omega(\lambda)}(\operatorname{Im} \lambda \neq 0)$,
3) $\operatorname{Im} \lambda \operatorname{Im} \omega(\lambda) \geq 0$ for $\operatorname{Im} \lambda \neq 0$.

Definition 2.2 ([13]). A Nevanlinna function $\omega(\lambda)$ is said to be an S-function if it is defined and analytic in $C \backslash[0, \infty)$ and $\omega(\lambda)>0$ for $\lambda<0$. A meromorphic $S$-function is said to be an $S_{0}$-function if $\omega(0)<\infty$.

Lemma 2.3 ([13]). The function $\frac{R_{2 n-2}(z)}{R_{2 n-3}(z)}$ is an $S_{0}$-function.
Definition 2.4. A polynomial is said to be Hermite-Biehler (HB) if all its zeros lie in the open upper half-plane.

It should be mentioned that the transformation $\lambda \rightarrow i \lambda$ transforms a HB polynomial into a so-called Hurwitz polynomial.

Theorem 2.5 (Hermite-Biehler theorem, see $[11,19]$ ). In order for the polynomial

$$
\omega(\lambda)=P(\lambda)+i Q(\lambda)
$$

where $P(\lambda)$ and $Q(\lambda)$ are real polynomials, to have no zeros in the closed lower half-plane $\operatorname{Im} \lambda \leq 0$, i.e. belong to $H B$, it is necessary and sufficient that the following conditions be satisfied:

1) the polynomials $P(\lambda)$ and $Q(\lambda)$ have only simple real zeros, while these zeros separate one another, i.e. between two successive zeros of one of these polynomials there lies exactly one zero of the other,
2) at some point $\lambda_{0}$ of the real axis

$$
Q^{\prime}\left(\lambda_{0}\right) P\left(\lambda_{0}\right)-Q\left(\lambda_{0}\right) P^{\prime}\left(\lambda_{0}\right)>0
$$

The fact that the two polynomials satisfy condition 1) will be expressed by saying that "the zeros of the polynomials $P(\lambda)$ and $Q(\lambda)$ are interlaced".

Now Lemma 2.3 and Theorem 2.5 imply the following result.
Corollary 2.6. The polynomial $P\left(\lambda^{2}\right)+i \lambda Q\left(\lambda^{2}\right)$ belongs to the Hermite-Biehler class.
Definition 2.7. The polynomial $\omega(\lambda)$ is said to be symmetric if $\omega(-\bar{\lambda})=\overline{\omega(\lambda)}$ for all $\lambda \in \mathbb{C}$. The polynomial $\omega(\lambda)$ is said to belong to the class SHB if it is symmetric and belongs to the Hermite-Biehler class.

For a symmetric polynomial $\omega(\lambda)$ the following is valid:

$$
\omega(\lambda)=P(\lambda)+i Q(\lambda)=P(\lambda)+i \lambda \hat{Q}(\lambda)=\tilde{P}\left(\lambda^{2}\right)+i \lambda \tilde{\hat{Q}}\left(\lambda^{2}\right)
$$

where $P(\lambda)$ and $\hat{Q}(\lambda)$ are real even functions. Here

$$
\tilde{P}\left(\lambda^{2}\right)=P(\lambda), \quad \tilde{\hat{Q}}\left(\lambda^{2}\right)=\hat{Q}(\lambda)
$$

## 3. INVERSE PROBLEMS I AND II

In this section we consider the following inverse problems.
Inverse problem I. Given the total length of the string $l>0$, the length of the right subinterval $l_{n} \in(0, l)$ and the spectrum $\left\{\nu_{k}\right\}(k= \pm 1, \pm 2, \ldots, \pm n)$ of problem (2.2)-(2.5). Find $\left\{m_{k}\right\}(k=1,2, \ldots, n)$ and $\left\{l_{k}\right\}(k=0,1, \ldots, n-1)$.

Inverse problem II. Given the total length of the string $l>0$, the length of the right subinterval $l_{n} \in(0, l)$ and the spectrum $\left\{\mu_{k}\right\}(k= \pm 1, \pm 2, \ldots, \pm n)$ of problem (2.2)-(2.4), (2.6). Find $\left\{m_{k}\right\}(k=1,2, \ldots, n)$ and $\left\{l_{k}\right\}(k=0,1, \ldots, n-1)$.

Theorem 3.1. Let $l>0$ and $l_{n} \in(0, l)$ be given together with the set of complex numbers $\left\{\nu_{k}\right\}(k= \pm 1, \pm 2, \ldots \pm n)$ which satisfy the conditions:

1) $\operatorname{Im} \nu_{k}>0$ for $k= \pm 1, \pm 2, \ldots \pm n$,
2) $\nu_{-k}=-\overline{\nu_{k}}$ for not pure imaginary $\nu_{-k}$ and the multiplicities of symmetrically located numbers are equal.
Then there exists a unique Stieltjes string, i.e. a unique set of intervals $l_{k}>0$ ( $k=$ $0,1, \ldots, n-1)$ of total length $\sum_{k=0}^{n-1} l_{k}=l-l_{n}$, a unique set of masses $m_{k}>0$ ( $k=1,2, \ldots, n$ ) and a unique positive number $\alpha$ which generate problem (2.2)-(2.5) with the spectrum coinciding with the set $\left\{\nu_{k}\right\}(k= \pm 1, \pm 2, \ldots, \pm n)$.
Proof. Let us construct the polynomial

$$
\begin{equation*}
\Phi(\lambda)=\prod_{-n, k \neq 0}^{n}\left(1-\frac{\lambda}{\nu_{k}}\right) \tag{3.1}
\end{equation*}
$$

Due to the symmetry of the zeros of this polynomial the following even polynomials are real:

$$
P\left(\lambda^{2}\right)=\frac{\Phi(\lambda)+\Phi(-\lambda)}{2}
$$

and

$$
\begin{equation*}
Q\left(\lambda^{2}\right)=\frac{\Phi(\lambda)-\Phi(-\lambda)}{2 i \lambda} \tag{3.2}
\end{equation*}
$$

Set

$$
\begin{equation*}
\alpha=Q(0)\left(\frac{1}{l-l_{n}}+\frac{1}{l_{n}}\right) . \tag{3.3}
\end{equation*}
$$

Using (3.1) and (3.2) we obtain

$$
\begin{equation*}
\alpha=i\left(\frac{1}{l-l_{n}}+\frac{1}{l_{n}}\right) \sum_{k=-n, k \neq 0}^{n} \frac{1}{\nu_{k}} \tag{3.4}
\end{equation*}
$$

and because of the symmetry in location of the zeros of the polynomial $\Phi(\lambda)$ and conditions 1) and 2) we conclude that $\alpha>0$. Set

$$
\begin{equation*}
m_{n}=-\alpha \lim _{|\lambda| \rightarrow \infty} \frac{P\left(\lambda^{2}\right)}{\lambda^{2} Q\left(\lambda^{2}\right)} \tag{3.5}
\end{equation*}
$$

The limit in the right-hand side of (3.5) exists because the degree of $P\left(\lambda^{2}\right)$ is $2 n$ and the degree of $Q\left(\lambda^{2}\right)$ is $2 n-2$. Moreover,

$$
\frac{P\left(\lambda^{2}\right)}{\lambda^{2} Q\left(\lambda^{2}\right)}=(i \lambda \mid \rightarrow \infty
$$

Conditions 1) and 2) imply

$$
i \sum_{k=-n, k \neq 0}^{n} \nu_{k}<0
$$

and according to (3.5)

$$
m_{n}=-\alpha\left(i \sum_{k=-n, k \neq 0}^{n} \nu_{k}\right)^{-1}>0
$$

Since $P\left(\lambda^{2}\right)+i \lambda Q\left(\lambda^{2}\right)=\Phi(\lambda)$ belongs to SHB, $\phi(\lambda):=P\left(\lambda^{2}\right)+i \lambda \alpha^{-1} Q\left(\lambda^{2}\right)$ is also a SHB polynomial by the Hermite-Biehler theorem. We consider the polynomial

$$
\phi\left(\lambda, m_{n}, l_{n}^{-1}\right):=P\left(\lambda^{2}\right)+\alpha^{-1}\left(m_{n} \lambda^{2}-l_{n}^{-1}\right) Q\left(\lambda^{2}\right)+i \lambda \alpha^{-1} Q\left(\lambda^{2}\right)
$$

as a perturbation of $P\left(\lambda^{2}\right)+i \lambda \alpha^{-1} Q\left(\lambda^{2}\right)$. Since $\phi(\lambda, \eta, \zeta)$ is a polynomial with respect to the variables $\lambda, \eta$ and $\zeta$, the zeros of it in the $\lambda$-plane are piecewise analytic and continuous functions of $\eta$ and of $\zeta$ ([20]). The zeros do not cross the real axis when $\eta$ changes from 0 to $m_{n}$ and $\zeta$ changes from 0 to $l_{n}^{-1}$. Otherwise, we would have $P\left(\lambda^{2}\right)=\lambda Q\left(\lambda^{2}\right)=0$ for some $\eta>0$, some $\zeta>0$ and some real $\lambda$. If this $\lambda \neq 0$ then $Q\left(\lambda^{2}\right)=0$ and consequently, $P\left(\lambda^{2}\right)=0$ and $\Phi(\lambda)=0$ for this real $\lambda$, what contradicts condition 1$)$. If $\phi(0, \eta, \zeta)=0$, then

$$
P(0)-\alpha^{-1} \zeta Q(0)=1-\alpha^{-1} \zeta Q(0)=0 .
$$

This is impossible for $\zeta \in\left(0, l_{n}^{-1}\right)$ due to (3.3). The degree of the polynomial $\phi(\lambda, \eta, \zeta)$ is $2 n$ for each $\eta \in\left[0, m_{n}\right)$ and each $\zeta \in\left[0, l_{n}^{-1}\right]$ and the degree is equal $2 n-1$ for $\eta=m_{n}$ and each $\zeta \in\left[0, l_{n}^{-1}\right]$. This means the zeros do not come from infinity. Therefore, $\phi(\lambda, \eta, \zeta) \in S H B$ for each $\zeta \in\left[0, l_{n}^{-1}\right]$ and each $\eta \in\left[0, m_{n}\right]$. This implies (see [19, p. 308]) that

$$
\frac{\alpha^{-1} Q(z)}{P(z)+\left(m_{n} z-l_{n}^{-1}\right) \alpha^{-1} Q(z)}
$$

is an S-function. Then according to [9] we have

$$
\begin{equation*}
\frac{\alpha^{-1} Q(z)}{P(z)+\left(m_{n} z-l_{n}^{-1}\right) \alpha^{-1} Q(z)}=a_{n-1}+\frac{1}{-b_{n-1} z+\frac{1}{a_{n-2}+\frac{1}{-b_{n-2} z+\ldots+\frac{1}{a_{1}+\frac{1}{-b_{1} z+\frac{1}{a_{0}}}}}}} \tag{3.6}
\end{equation*}
$$

with $a_{k}>0$ and $b_{k}>0$ for each $k$.
We identify $a_{k}$ with the length of $k$-th interval and $b_{k}$ with the $k$-th mass of a Stieltjes string, i.e.

$$
\begin{equation*}
\frac{\alpha^{-1} Q\left(\lambda^{2}\right)}{P\left(\lambda^{2}\right)+\left(m_{n} \lambda^{2}-l_{n}^{-1}\right) \alpha^{-1} Q\left(\lambda^{2}\right)}=\frac{R_{2 n-2}\left(\lambda^{2}\right)}{R_{2 n-3}\left(\lambda^{2}\right)}, \tag{3.7}
\end{equation*}
$$

where $R_{2 n-2}\left(\lambda^{2}\right)$ and $R_{2 n-3}\left(\lambda^{2}\right)$ are the corresponding polynomials for this Stieltjes string. Consequently,

$$
\begin{gathered}
\alpha^{-1} Q\left(\lambda^{2}\right)=T R_{2 n-2}\left(\lambda^{2}\right) \\
P\left(\lambda^{2}\right)+\left(m_{n} \lambda^{2}-l_{n}^{-1}\right) \alpha^{-1} Q\left(\lambda^{2}\right)=T R_{2 n-3}\left(\lambda^{2}\right)
\end{gathered}
$$

where $T$ is a positive constant. Therefore,

$$
\Phi(\lambda)=P\left(\lambda^{2}\right)+i \lambda Q\left(\lambda^{2}\right)=T\left(R_{2 n-3}\left(\lambda^{2}\right)+\left(-m_{n} \lambda^{2}+l_{n}^{-1}+i \lambda \alpha\right) R_{2 n-2}\left(\lambda^{2}\right)\right)
$$

According to (2.10), this means that the set $\left\{\nu_{k}\right\}$ is the spectrum of problem (2.2)-(2.5) with the masses $b_{k}(k=1,2, \ldots, n-1)$ and $m_{n}$ and the lengths $a_{k}$ ( $k=0,1, \ldots, n-1$ ) and $l_{n}$ damped at the mass $m_{n}$ with the coefficient of damping $\alpha$. The length of the interval between the left end of the string and the mass $m_{n}$ according to [9] is equal to $\frac{R_{2 n-2}(0)}{R_{2 n-3}(0)}$. From (3.7) we obtain

$$
\begin{equation*}
\frac{R_{2 n-2}(0)}{R_{2 n-3}(0)}=\frac{\alpha^{-1} Q(0)}{P(0)-l_{n}^{-1} \alpha^{-1} Q(0)} \tag{3.8}
\end{equation*}
$$

Using (3.3) and the evident identity $P(0)=1$ we obtain

$$
\begin{equation*}
\frac{R_{2 n-2}(0)}{R_{2 n-3}(0)}=l-l_{n} \tag{3.9}
\end{equation*}
$$

Let us prove uniqueness of the solution to our inverse problem. Suppose there exists another Stieljes string with the same total length $l$, the same length of the right interval $l_{n}$ between the fixed end and the first mass from the right (which is damped) having the same spectrum $\left\{\nu_{k}\right\}_{k=-n, k \neq 0}^{n}$. In other words, we suppose that there exist sequences of positive numbers $\left\{\tilde{m}_{k}\right\}_{k=1}^{n}$ and $\left\{\tilde{l}_{k}\right\}_{k=0}^{n-1}\left(\sum_{k=0}^{n-1} \tilde{l}_{k}=l-l_{n}\right)$, not identical with sets $\left\{m_{k}\right\}_{k=1}^{n}$ and $\left\{l_{k}\right\}_{k=0}^{n-1}$, which together with $\tilde{l}_{n}=l_{n}$ generate problem (2.2)-(2.5) having the same spectrum $\left\{\nu_{k}\right\}_{k=-n, k \neq 0}^{n}$.

All the quantities related to problem (2.2)-(2.5) generated by $\left\{\tilde{m}_{k}\right\}_{k=1}^{n}$ and $\left\{\tilde{l}_{k}\right\}_{k=0}^{n-1}$ and $\tilde{l}_{n}=l_{n}$ will have the sign tilde. Then we have the following analogue of (2.12):

$$
\begin{equation*}
\frac{\tilde{R}_{2 n-2}\left(\lambda^{2}\right)}{\tilde{R}_{2 n-3}\left(\lambda^{2}\right)}=\tilde{l}_{n-1}+\frac{1}{-\tilde{m}_{n-1} \lambda^{2}+\frac{1}{\overline{l_{n-2}+\frac{\tilde{m}_{n-2} \lambda^{2}+\ldots+}{\tilde{l}_{1}+\frac{1}{-\tilde{m}_{1} \lambda^{2}+\frac{1}{\zeta_{0}}}}}}} \tag{3.10}
\end{equation*}
$$

The analogue of (2.10) is

$$
\begin{equation*}
\tilde{\phi}(\lambda):=\tilde{R}_{2 n-3}\left(\lambda^{2}\right)+\left(-\tilde{m}_{n} \lambda^{2}+i \lambda \tilde{\alpha}+l_{n}^{-1}\right) \tilde{R}_{2 n-2}\left(\lambda^{2}\right)=0 \tag{3.11}
\end{equation*}
$$

Since the sets of zeros of $\phi(\lambda)$ and $\tilde{\phi}(\lambda)$ coincide with the spectrum $\left\{\nu_{k}\right\}_{k=-n, k \neq 0}^{n}$, we conclude that

$$
\begin{equation*}
\tilde{\phi}(\lambda)=C \phi(\lambda) \tag{3.12}
\end{equation*}
$$

with a constant $C$ which is positive, because the following are positive

$$
\tilde{\phi}(0)=\tilde{R}_{2 n-3}(0)+l_{n}^{-1} \tilde{R}_{2 n-2}(0) \quad \text { and } \quad \phi(0)=R_{2 n-3}(0)+l_{n}^{-1} R_{2 n-2}(0) .
$$

Using (2.10) and (3.11) we obtain from (3.12):
$\tilde{R}_{2 n-3}\left(\lambda^{2}\right)+\left(-\tilde{m}_{n} \lambda^{2}+l_{n}^{-1}\right) \tilde{R}_{2 n-2}\left(\lambda^{2}\right)=C\left(R_{2 n-3}\left(\lambda^{2}\right)+\left(-m_{n} \lambda^{2}+l_{n}^{-1}\right) R_{2 n-2}\left(\lambda^{2}\right)\right)$
and

$$
\begin{equation*}
\tilde{\alpha} \tilde{R}_{2 n-2}\left(\lambda^{2}\right)=C \alpha R_{2 n-2}\left(\lambda^{2}\right) \tag{3.14}
\end{equation*}
$$

Equations (3.13) and (3.14) imply

$$
\begin{equation*}
\frac{\tilde{R}_{2 n-3}\left(\lambda^{2}\right)}{\tilde{\alpha} \tilde{R}_{2 n-2}\left(\lambda^{2}\right)}-\frac{\tilde{m}_{n} \lambda^{2}}{\tilde{\alpha}}+\frac{1}{\tilde{\alpha} l_{n}}=\frac{R_{2 n-3}\left(\lambda^{2}\right)}{\alpha R_{2 n-2}\left(\lambda^{2}\right)}-\frac{m_{n} \lambda^{2}}{\alpha}+\frac{1}{\alpha l_{n}} \tag{3.15}
\end{equation*}
$$

Since the polynomials $R_{2 n-2}\left(\lambda^{2}\right), R_{2 n-3}\left(\lambda^{2}\right), \tilde{R}_{2 n-2}\left(\lambda^{2}\right)$ and $\tilde{R}_{2 n-3}\left(\lambda^{2}\right)$ are of the same degree, (3.15) implies

$$
\begin{equation*}
\frac{\tilde{m}_{n}}{\tilde{\alpha}}=\frac{m_{n}}{\alpha} \tag{3.16}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\tilde{R}_{2 n-3}\left(\lambda^{2}\right)}{\tilde{\alpha} \tilde{R}_{2 n-2}\left(\lambda^{2}\right)}+\frac{1}{\tilde{\alpha} l_{n}}=\frac{R_{2 n-3}\left(\lambda^{2}\right)}{\alpha R_{2 n-2}\left(\lambda^{2}\right)}+\frac{1}{\alpha l_{n}} \tag{3.17}
\end{equation*}
$$

and, consequently,

$$
\begin{equation*}
\frac{\tilde{R}_{2 n-3}(0)}{\tilde{\alpha} \tilde{R}_{2 n-2}(0)}+\frac{1}{\tilde{\alpha} l_{n}}=\frac{R_{2 n-3}(0)}{\alpha R_{2 n-2}(0)}+\frac{1}{\alpha l_{n}} \tag{3.18}
\end{equation*}
$$

Setting $\lambda=0$ in (2.12) and (3.10) we obtain

$$
\begin{equation*}
\frac{R_{2 n-2}(0)}{R_{2 n-3}(0)}=\frac{\tilde{R}_{2 n-2}(0)}{\tilde{R}_{2 n-3}(0)}=\sum_{k=0}^{n-1} l_{k}=\sum_{k=0}^{n-1} \tilde{l}_{k}=l-l_{n} \tag{3.19}
\end{equation*}
$$

Combining (3.19) with (3.18) we arrive at $\tilde{\alpha}=\alpha$. Now (3.16) implies $m_{n}=\tilde{m}_{n}$ and it follows from (3.17) that

$$
\frac{\tilde{R}_{2 n-2}\left(\lambda^{2}\right)}{\tilde{R}_{2 n-3}\left(\lambda^{2}\right)}=\frac{R_{2 n-2}\left(\lambda^{2}\right)}{R_{2 n-3}\left(\lambda^{2}\right)}
$$

Since the left-hand sides of (2.12) and (3.10) coincide, we conclude that $m_{k}=\tilde{m}_{k}$ for $k=1,2, \ldots, n-1$ and $l_{k}=\tilde{l}_{k}$ for $k=0,1,2, \ldots, n-1$. The theorem is proved.

Theorem 3.1 gives a solution of Inverse problem I. A solution of Inverse problem II is similar. It is necessary just to delete the summand $l_{n}^{-1}$ in (3.3), (3.4), set $l_{n}^{-1}=0$ in (3.6), (3.7), (3.8), (3.11), (3.13), (3.15), (3.17), (3.18).

## 4. COMPARISON OF PROBLEMS I AND II

In this section we compare the functions $\phi(\lambda)$ and $\psi(\lambda)$ defined by (2.10) and (2.11) and the sets of their zeros $\left\{\nu_{k}\right\}$ and $\left\{\mu_{k}\right\}$.

Comparison of (2.10) with (2.11) gives

$$
\phi(\lambda)=\psi(\lambda)+\frac{\psi(\lambda)-\psi(-\lambda)}{2 i \lambda \alpha l_{n}}
$$

$$
\begin{align*}
\frac{\phi(\lambda)}{\psi(\lambda)} & =1+\frac{l_{n}^{-1}}{-m_{n} \lambda^{2}+i \alpha \lambda+\frac{1}{R_{2 n-2}\left(\lambda^{2}\right) / R_{2 n-3}\left(\lambda^{2}\right)}} \\
& =1+\frac{l_{n}^{-1}}{-m_{n} \lambda^{2}+i \alpha \lambda+\frac{1}{l_{n-1}+\frac{1}{-m_{n-1} \lambda^{2}+\ldots+\frac{1}{l_{1}+\frac{1}{-m_{1} \lambda^{2}+\frac{1}{l_{0}}}}}}} \tag{4.1}
\end{align*}
$$

It follows from (4.1) that

$$
\begin{equation*}
\frac{\phi(0)}{\psi(0)}=\frac{l}{l_{n}} \tag{4.2}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{R_{2 n-2}(\lambda)}{R_{2 n-3}(\lambda)}=\left(m_{n} \lambda-i \alpha \lambda^{1 / 2}+\left(l_{n}\left(\frac{\phi\left(\lambda^{1 / 2}\right)}{\psi\left(\lambda^{1 / 2}\right)}-1\right)\right)^{-1}\right)^{-1} \tag{4.3}
\end{equation*}
$$

is an $S_{0}$-function.
Theorem 4.1. For the polynomials $\phi(\lambda)$ and $\psi(\lambda)$ of degree $2 n$ each to be the characteristic polynomials of problems I and II, respectively, normalized by (4.2) it is necessary and sufficient that the function

$$
\begin{equation*}
\left(m_{n} \lambda-i \alpha \lambda^{1 / 2}+\left(l_{n}\left(\frac{\phi\left(\lambda^{1 / 2}\right)}{\psi\left(\lambda^{1 / 2}\right)}-1\right)\right)^{-1}\right)^{-1} \tag{4.4}
\end{equation*}
$$

with

$$
\begin{equation*}
\alpha=:\left(\frac{1}{l-l_{n}}+\frac{1}{l_{n}}\right) \lim _{\lambda \rightarrow 0} \frac{\phi(\lambda)-\phi(-\lambda)}{2 i \lambda} \tag{4.5}
\end{equation*}
$$

and

$$
\begin{equation*}
m_{n}=:-i \alpha \lim _{|\lambda| \rightarrow \infty} \frac{\phi(\lambda)+\phi(-\lambda)}{\lambda(\phi(\lambda)-\phi(-\lambda)} \tag{4.6}
\end{equation*}
$$

be a rational $S_{0}$-function.
Proof. Necessity follows from (4.3) and (2.12). Now let the function (4.4) be an $S_{0}$-function. Then it can be expanded into a continued fraction:

$$
\begin{aligned}
& \left(m_{n} \lambda-i \alpha \lambda^{1 / 2}+\left(l_{n}\left(\frac{\phi\left(\lambda^{1 / 2}\right)}{\psi\left(\lambda^{1 / 2}\right)}-1\right)\right)^{-1}\right)^{-1} \\
& =a_{n-1}+\frac{1}{-b_{n-1} \lambda^{2}+\frac{1}{a_{n-2}+\frac{1}{-b_{n-2} \lambda^{2}+\ldots+\frac{1}{a_{1}+\frac{1}{-b_{1} \lambda^{2}+\frac{1}{b_{0}}}}}}}
\end{aligned}
$$

where $a_{k}>0$ for $k=0,1, \ldots, n-1$ and $b>0$ for $k=1,2, \ldots, n-1$.

We identify $b_{k}$ with masses on a Stieltjes string and $a_{k}$ with the subintervals into which the masses divide the length $l-l_{n}$. This data together with the given $l_{n}$ and $\alpha$ obtained from (4.5) and $m_{n}$ obtained from (4.6) generate problems (2.2)-(2.5) and (2.2)-(2.4), (2.6) which have spectra $\left\{\nu_{k}\right\}$ and $\left\{\mu_{k}\right\}$, respectively. The proof is complete.

Comparing (2.10) with (2.11), which can be rewritten as

$$
\begin{equation*}
\phi(\lambda)=(-1)^{n} m_{n} \prod_{1}^{n-1} m_{k} l_{k} \prod_{-n, k \neq 0}^{n}\left(\nu_{k}-\lambda\right) \tag{4.7}
\end{equation*}
$$

and

$$
\begin{equation*}
\psi(\lambda)=(-1)^{n} m_{n} \prod_{1}^{n-1} m_{k} l_{k} \prod_{-n, k \neq 0}^{n}\left(\mu_{k}-\lambda\right) \tag{4.8}
\end{equation*}
$$

we obtain

$$
\phi(\lambda)-\psi(\lambda)=l_{n}^{-1} R_{2 n-2}\left(\lambda^{2}\right)
$$

Let us introduce the following notation:

$$
\begin{gather*}
M_{p}=\sum_{k=-n, k<k^{\prime}<\ldots<k^{(p-1)},}^{k=n, k_{k \neq 0,}^{k^{\prime}=n, \ldots k^{(p-1)}=n} \sum_{k^{\prime} \neq 0, \ldots, k^{(p-1)} \neq 0} \nu_{k} \nu_{k^{\prime}} \ldots \nu_{k^{(p-1)}},} \begin{array}{l}
k=n, k^{\prime}=n, \ldots k^{(p-1)}=n \\
N_{p}=\sum_{k=-n, k<k^{\prime}<\ldots<k^{(p-1)},} \ldots \sum_{k \neq 0, k^{\prime} \neq 0, \ldots, k^{(p-1)} \neq 0} \mu_{k} \mu_{k^{\prime}} \ldots \mu_{k^{(p-1)}}
\end{array} . l \tag{4.9}
\end{gather*}
$$

Then

$$
\begin{aligned}
& R_{2 n-2}\left(\lambda^{2}\right)= l_{n}(-1)^{n} m_{n} \prod_{1}^{n-1} m_{k} l_{k}\left(\lambda^{2 n-2}\left(N_{2}-M_{2}\right)+\lambda^{2 n-4}\left(N_{4}-M_{4}\right)+\ldots\right. \\
&\left.\ldots+N_{2 n}-M_{2 n}\right) . \\
& R_{2 n-2}\left(\lambda^{2}\right)=\frac{\phi(\lambda)-\phi(-\lambda)}{2 i \alpha \lambda} \\
&=i(-1)^{n} m_{n} \prod_{1}^{n-1} m_{k} l_{k} \alpha^{-1}\left(\lambda^{2 n-2} N_{1}+\lambda^{2 n-4} N_{3}+\ldots+N_{2 n-1}\right) .
\end{aligned}
$$

Comparing (4.9) with (4.10) we obtain

$$
\begin{aligned}
l_{n}\left(N_{2}-M_{2}\right) & =\left(\frac{1}{l-l_{n}}+\frac{1}{l_{n}}\right)^{-1}\left(\sum_{k=-n, k \neq 0}^{n} \frac{1}{\nu_{k}}\right)^{-1} N_{1}, \\
l_{n}\left(N_{4}-M_{4}\right) & =\left(\frac{1}{l-l_{n}}+\frac{1}{l_{n}}\right)^{-1}\left(\sum_{k=-n, k \neq 0}^{n} \frac{1}{\nu_{k}}\right)^{-1} N_{3}, \\
& \ldots \\
l_{n}\left(N_{2 n}-M_{2 n}\right) & =\left(\frac{1}{l-l_{n}}+\frac{1}{l_{n}}\right)^{-1}\left(\sum_{k=-n, k \neq 0}^{n} \frac{1}{\nu_{k}}\right)^{-1} N_{2 n-1} .
\end{aligned}
$$

Definitions (4.7) and (4.8) imply

$$
\begin{equation*}
\phi(\lambda)-\phi(-\lambda)=\psi(\lambda)-\psi(-\lambda) \tag{4.11}
\end{equation*}
$$

Substituting (4.7) and (4.8) into (4.11) we obtain

$$
\begin{equation*}
M_{2 k-1}=N_{2 k-1}, \quad k=1,2, \ldots, n \tag{4.12}
\end{equation*}
$$

Using (2.11) we derive

$$
\begin{equation*}
R_{2 n-3}\left(\lambda^{2}\right)=\frac{\psi(\lambda)+\psi(-\lambda)}{2}+m_{n} \lambda^{2} R_{2 n-2}\left(\lambda^{2}\right) \tag{4.13}
\end{equation*}
$$

Substituting (4.8) and (4.10) into (4.13) we obtain

$$
\begin{align*}
R_{2 n-3}\left(\lambda^{2}\right)= & (-1)^{n} m_{n} \prod_{k=1}^{n-1} m_{k} l_{k}\left(\lambda^{2 n-2}\left(M_{2}-N_{1}^{-1} N_{3}\right)+\lambda^{2 n-4}\left(M_{4}-N_{1}^{-1} N_{5}\right)+\ldots\right.  \tag{4.14}\\
& \left.+\lambda^{2}\left(M_{2 n-2}-N_{1}^{-1} N_{2 n-1}\right)+M_{2 n}\right) .
\end{align*}
$$

Taking into account that $R_{2 n-3}(z)$ has zeros only on the positive half-axis and comparing (4.14) with (2.8) and (2.9) we obtain

$$
(-1)^{k}\left(M_{2 k}-N_{1}^{-1} N_{2 k+1}\right)>0, k=1,2, \ldots, n-1, \quad(-1)^{n} M_{2 n}>0
$$

Using interlacing of the zeros of $R_{2 n-3}(z)$ with the zeros of $R_{2 n-2}(z)$ we obtain

$$
\begin{gathered}
(-1)^{k-1} \frac{M_{2 k}-N_{1}^{-1} N_{2 k+1}}{M_{2}-N_{1}^{-1} N_{3}}>(-1)^{k-1} \frac{N_{2 k-1}}{N_{1}}, \quad k=2,3, \ldots, n-1, \\
(-1)^{n-1} \frac{M_{2 n}}{M_{2}-N_{1}^{-1} N_{3}}>(-1)^{n-1} \frac{N_{2 n-1}}{N_{1}}
\end{gathered}
$$

Equation (4.12) for $k=1$ is equivalent to

$$
\sum_{k=-n, k \neq 0}^{n} \operatorname{Im} \mu_{k}=\sum_{k=-n, k \neq 0}^{n} \operatorname{Im} \nu_{k}
$$

## 5. INVERSE PROBLEM FOR A DAMPED STIELTJES STRING

Now we come back to the problems generated by equation (2.1) where all the masses are damped. By an inverse problem we mean recovering the parameters of problems generated by equation

$$
\begin{equation*}
\frac{u_{k}-u_{k+1}}{l_{k}}+\frac{u_{k}-u_{k-1}}{l_{k-1}}-m_{k} \lambda^{2} u_{k}+i \alpha_{k} \lambda u_{k}=0 \quad(k=1,2, \ldots, n) \tag{5.1}
\end{equation*}
$$

and conditions (2.4), (2.5) and (2.4), (2.6), i.e. $\left\{m_{k}\right\},\left\{\alpha_{k}\right\}(k=1,2, \ldots, n),\left\{l_{k}\right\}$ ( $k=0,1, \ldots, n$ ) using the spectra of these problems and the total length of the string $l=l_{0}+l_{1}+\ldots+l_{n}$.

Suppose we know $\left\{\mu_{k}\right\}(k= \pm 1, \pm 2, \ldots \pm n)$ eigenvalues of problem (5.1), (2.4), (2.5), i.e. the zeros of the polynomial $R_{2 n-1}(\lambda)$ which here are obtained from the recurrence relations

$$
\begin{gathered}
R_{2 k-1}\left(\lambda^{2}\right)=\left(-\lambda^{2} m_{k}+i \alpha \lambda\right) R_{2 k-2}\left(\lambda^{2}\right)+R_{2 k-3}\left(\lambda^{2}\right) \\
R_{2 k}\left(\lambda^{2}\right)=l_{k} R_{2 k-1}\left(\lambda^{2}\right)+R_{2 k-2}\left(\lambda^{2}\right), \quad\left(k=1,2, \ldots, n, R_{-1}\left(\lambda^{2}\right)=\frac{1}{l_{0}}, R_{0}\left(\lambda^{2}\right)=1\right)
\end{gathered}
$$

and $\left\{\nu_{k}\right\}(k= \pm 1, \pm 2, \ldots \pm n)$ eigenvalues of problem (5.1), (2.4), (2.6), i.e. the zeros of the polynomial $R_{2 n}(\lambda)$ and the total length of the string $l$. We construct the polynomials

$$
p(\lambda)=\prod_{k=1}^{n}\left(1-\frac{\lambda}{\mu_{k}}\right)\left(1-\frac{\lambda}{\mu_{-k}}\right), \quad q(\lambda)=\prod_{k=1}^{n}\left(1-\frac{\lambda}{\nu_{k}}\right)\left(1-\frac{\lambda}{\nu_{-k}}\right)
$$

These polynomials have the same sets of zeros as $R_{2 n}(\lambda)$ and $R_{2 n-1}(\lambda)$, respectively. Therefore,

$$
R_{2 n-1}(\lambda)=T_{1} p(\lambda), \quad R_{2 n}(\lambda)=T_{2} q(\lambda)
$$

and

$$
\begin{equation*}
\frac{R_{2 n}(\lambda)}{R_{2 n-1}(\lambda)}=\frac{T_{2}}{T_{1}} \frac{q(\lambda)}{p(\lambda)} \tag{5.2}
\end{equation*}
$$

We construct $p(\lambda)$ and $q(\lambda)$ using $\left\{\mu_{k}\right\}(k= \pm 1, \pm 2, \ldots \pm n)$ and $\left\{\nu_{k}\right\}(k=$ $\pm 1, \pm 2, \ldots \pm n)$. To find $\frac{T_{2}}{T_{1}}$ we substitute $\lambda=0$ in (5.2):

$$
\frac{R_{2 n}(0)}{R_{2 n-1}(0)}=\frac{T_{2}}{T_{1}} \frac{q(0)}{p(0)}=\frac{T_{2}}{T_{1}} .
$$

Substituting $\lambda=0$ into (2.8) we obtain

$$
\frac{T_{2}}{T_{1}}=l_{0}+l_{1}+\ldots+l_{n}=l
$$

Thus the sets $\left\{\mu_{k}\right\}$ and $\left\{\nu_{k}\right\}$ together with given $l$ uniquely determine the rational function

$$
\frac{R_{2 n}(\lambda)}{R_{2 n-1}(\lambda)}=l \frac{q(\lambda)}{p(\lambda)}
$$

By expanding $\frac{R_{2 n}(\lambda)}{R_{2 n-1}(\lambda)}$ into continued fraction according to (2.8) we can find $\left\{m_{k}\right\}$, $\left\{\alpha_{k}\right\}(k=1,2, \ldots, n)$ and $\left\{l_{k}\right\}(k=0,1, \ldots, n)$. Hence, we have proved the following theorem.

Theorem 5.1. The eigenvalues $\left\{\mu_{k}\right\}(k= \pm 1, \pm 2, \ldots \pm n)$ and $\left\{\nu_{k}\right\}(k=$ $\pm 1, \pm 2, \ldots \pm n)$ together with the given total length $l$ uniquely determine $\left\{m_{k}\right\},\left\{\alpha_{k}\right\}$ $(k=1,2, \ldots, n)$ and $\left\{l_{k}\right\}(k=0,1, \ldots, n)$.

Now let us find conditions which must be satisfied by sets of complex numbers $\left\{\nu_{k}\right\}$ and $\left\{\mu_{k}\right\}(k= \pm 1, \pm 2, \ldots, \pm n)$ to be the spectra of problems (5.1), (2.4), (2.5) and (5.1) (2.4), (2.6), respectively.

Theorem 5.2. Let $\left\{\nu_{k}\right\}$ and $\left\{\mu_{k}\right\}(k= \pm 1, \pm 2, \ldots, \pm n)$ be two sequences of complex numbers and let $l$ be a positive number. In order $\left\{\nu_{k}\right\}$ and $\left\{\mu_{k}\right\}(k= \pm 1, \pm 2, \ldots, \pm n)$ be the spectra of problems (5.1), (2.4), (2.5) and (5.1) (2.4), (2.6), respectively, with $m_{k}>0, \alpha_{k}>0(k=1,2, \ldots, n)$, it is necessary and sufficient that:

1) $\left\{\mu_{k}\right\} \cap\left\{\nu_{k}\right\}=\emptyset$,
2) the product $l \prod_{k=-n, k \neq 0}^{n}\left(1-\frac{\lambda}{\nu_{k}}\right)\left(1-\frac{\lambda}{\mu_{k}}\right)^{-1}$ can be presented as a continued fraction of the form:

$$
\begin{align*}
& l \prod_{k=-n, k \neq 0}^{n}\left(1-\frac{\lambda}{\nu_{k}}\right)\left(1-\frac{\lambda}{\mu_{k}}\right)^{-1} \\
& =a_{n}+\frac{1}{b_{n} \lambda^{2}+i c_{n} \lambda+\frac{1}{a_{n-1}+\frac{1}{b_{n-1} \lambda^{2}+i c_{n-1} \lambda+\ldots+\frac{1}{a_{1}+\frac{1}{b_{1} \lambda^{2}+i c_{1} \lambda+\frac{1}{a_{0}}}}}}} \tag{5.3}
\end{align*}
$$

with $a_{k}>0(k=0,1,2, \ldots, n), b_{k}<0, c_{k} \geq 0(k=0,1,2, \ldots, n)$.
Proof. Let us use the continued fraction (5.3) to construct another continued fraction:

$$
a_{n}+\frac{1}{b_{n} \lambda^{2}+\frac{1}{a_{n-1}+\frac{1}{b_{n-1} \lambda^{2}+\ldots \frac{1}{a_{1}+\frac{1}{b_{1} \lambda^{2}+\frac{1}{a_{0}}}}}}} .
$$

It is clear from (5.3) that $\sum_{k=0}^{n} a_{k}=l$. According to [9] this fraction can be identified as the ratio of two polynomials the zeros of the numerator are the eigenvalues of an undamped Stieltjes string with fixed ends and the zeros of the denominator are the eigenvalues of the same string with the left end fixed and the right end free. The masses of this string are $\left|b_{k}\right|$ at distances $a_{k}$. The total length of the string is equal $l$. Now let us consider the same string (the same masses $\left|b_{k}\right|$ and the same lengths of intervals $a_{k}$ ) but with damping proportional to $c_{k}$ at mass $\left|b_{k}\right|$. This new damped string generates the continued fraction (5.3). The theorem is proved.

## Remark 5.3.

a) Due to conditions 1), 2) $\mu_{-k}=-\bar{\mu}_{k}$ for each not pure imaginary $\mu_{k}$ and $\nu_{-k}=-\bar{\nu}_{k}$ for each not pure imaginary $\nu_{k}$.
b) Condition 2) in explicit form consists of rather involved relations. The first, the second and the third of them, however, are

$$
\begin{aligned}
& \sum_{k=-n, k \neq 0}^{n} \operatorname{Im} \mu_{k}=\sum_{k=-n, k \neq 0}^{n} \operatorname{Im} \nu_{k}, \\
& \sum_{k=-n, k^{\prime}=-n, k \neq 0, k^{\prime} \neq 0}^{k=n, k^{\prime}=n} \mu_{k} \mu_{k^{\prime}}>\sum_{k=-n, k^{\prime}=-n, k \neq 0, k^{\prime} \neq 0}^{k=n, k^{\prime}=n} \nu_{k} \nu_{k^{\prime}}, \\
& \sum_{-n, k \neq 0}^{n} \operatorname{Im} \mu_{k}+ \\
& +i \sum_{k=-n, k^{\prime}=-n, k^{\prime \prime}=-n, k \neq 0, k^{\prime} \neq 0, k^{\prime \prime} \neq 0}^{k=n, k^{\prime}=n, k^{\prime \prime}=n}\left(\mu_{k} \mu_{k^{\prime}} \mu_{k} "-\nu_{k} \nu_{k^{\prime}} \nu_{k^{\prime}}\right) . \\
& \left(\sum_{k=-n, k^{\prime}=-n, k \neq 0, k^{\prime} \neq 0,}^{k=n, k^{\prime}=n,}\left(\mu_{k} \mu_{k^{\prime}}-\nu_{k} \nu_{k^{\prime}}\right)^{-1}>0 .\right.
\end{aligned}
$$

## REFERENCES

[1] T. Aktosun, M. Klaus, C. van der Mee, Wave scattering in one dimension with absorption, J. Math. Phys. 39 (1998) 4, 1957-1992.
[2] D.Z. Arov, Realization of a canonical system with a dissipative boundary condition at one end of the segment in terms of the coefficient of dynamical compliance, Sibirsk. Math. Z. 16 (1975), 440-463 [in Russian]; English trans.: Sibirian Math. J. 16 (1975), 335-352.
[3] O. Boyko, V. Pivovarchik, Inverse problem for Stieltjes string damped at one end, Methods Funct. Anal. Topology 14 (2008) 1, 10-19.
[4] O. Boyko, V. Pivovarchik, The inverse three-spectral problem for a Stieltjes string and the inverse problem with one-dimensional damping, Inverse Problems 24 (2008) 1, $015019,13 \mathrm{pp}$.
[5] R. Courant, D. Hilbert, Methods of Mathematical Physics, Vol. 1, Interscience, New York, 1953.
[6] S. Cox, M. Embree, J. Hokanson, One can hear the composition of a string: experiments with an inverse eigenvalue problem, SIAM Rev. 54 (2012) 1, 157-178.
[7] S. Cox, E. Zuazua, The rate at which energy decays in a string damped at one end, Indiana Univ. Math. J. 44 (1995) 2, 545-573.
[8] R. Duffing, A Minimax theory for overdamped networks, Arch. Rat. Mech. Anal. 4 (1955), 221-233.
[9] F.R. Gantmakher, M.G. Krein, Oscillating Matrices and Kernels and Vibrations of Mechanical Systems, GITTL, Moscow-Leningrad, 1950 [in Russian]; German transl. Akademie Verlag, Berlin, 1960; Revised edition, AMS Chelsea Publishing, Providence, RI, 2002.
[10] G. Gubreev, V. Pivovarchik, Spectral analysis of T. Regge problem with parameters, Funktsional. Anal. i Prilozhen. 31 (1997) 1, 70-74 [in Russian]; English transl.: Funct. Anal. Appl. 31 (1997) 1, 54-57.
[11] C. Hermite, Extract d'une lettre de M.Ch. Hermite de Paris a' Mr. Borchardt de Berlin sur le nombre des racines d'une e'quation alge'brique comprises entre des limites donne'es, J. Reine Angew. Math. 52 (1856), 39-51; reprinted in his Oeuvres, Vol. 1, Gauthier-Villars, Paris, 1905, 397-414.
[12] M. Jaulent, Inverse scattering problems in absorbing media, J. Math. Phys. 17 (1976) 7, 1351-1360.
[13] I.S. Kac, M.G. Krein, On spectral functions of a string, [in:] F.V. Atkinson, Discrete and Continuous Boundary Problems (Russian translation), Moscow, Mir, 1968, 648-737 (Addition II); I.C. Kac, M.G. Krein, On the spectral function of the string, Amer. Math. Soc., Translations, Ser. 2, 103 (1974), 19-102.
[14] M.G. Krein, On some new problems of the theory of vibrations of Sturm systems, Prikladnaya Matematika i Mekhanika 16 (1952) 5, 555-568 [in Russian].
[15] M.G. Krein, A.A. Nudelman, On direct and inverse problems for frequencies of boundary dissipation of inhomogeneous string, Doklady AN SSSR 247 (1979) 5, 1046-1049 [in Russian].
[16] M.G. Krein, A.A. Nudelman, On some spectral properties of an inhomogeneous string with dissipative boundary condition, J. Operator Theory 22 (1989), 369-395 [in Russian].
[17] P. Lancaster, J. Maroulas, Inverse eigenvalue problems for damped vibrating Systems, J. Math. Anal. Appl. 123 (1987), 238-261.
[18] C.K. Law, V. Pivovarchik, W.C. Wang, A polynomial identity and its application to inverse problems in Stieltjes strings, Oper. Matrices 7 (2013) 3, 603-617.
[19] B.Ja. Levin, Distribution of Zeros of Entire Functions, Transl. Math. Monographs, Amer. Math. Soc., Providence, R.I., 1980.
[20] A.I. Markushevich, Theory of Analytic Functions, Vol. 1, Nauka, Moscow, 1968 [in Russian]; Revised English edition, translated and edited by R.A. Silverman, Prentice Hall Inc., Englewood Cliffs, N.Y., 1967.
[21] V. Pivovarchik, Inverse problem for a smooth string damped with damping at one end, J. Operator Theory 38 (1997), 243-263.
[22] V. Pivovarchik, Direct and inverse problems for a damped string, J. Operator Theory 42 (1999), 189-220.
[23] V. Pivovarchik, C. van der Mee, The inverse generalized Regge problem, Inverse Problems 17 (2000), 1831-1845.
[24] V. Pivovarchik, Existence of a tree of Stieltjes strings corresponding to two given spectra, J. Phys. A: Math. Theor. 42 (2009), 375 213, 16 pp.
[25] V. Pivovarchik, N. Rozhenko, C. Tretter, Dirichlet-Neumann inverse spectral problem for a star graph of Stieltjes strings, Linear Algebra Appl. 439 (2013) 8, 2263-2292.
[26] V. Pivovarchik, O. Taistruk, On charateristic functions of operators on equilateral graph, Methods Funct. Anal. Topology 18 (2012) 2, 189-197.
[27] M.A. Shubov, Asymptotics of resonances and geometry of resonance states in the problem of scattering of acoustic waves by a spherically symmetric inhomogeneity of the density, Differential Integral Equations 8 (1995) 5, 1073-1115.
[28] T.-L. Stieltjes, Recherches sur les fractions continues, Ann. Fac. Sci. Toulouse 8 (1894), 1-122, 9 (1895), 1-47.
[29] K. Veselić, On linear vibrational systems with one dimensional damping, Appl. Anal. 29 (1988), 1-18.
[30] K. Veselić, On linear vibrational systems with one dimensional damping II, Integral Equations Operator Theory 13 (1990), 883-897.
[31] M. Yamamoto, Inverse eigenvalue problem for a vibration of a string with viscous drag, J. Math. Anal. Appl. 152 (1990), 20-34.

## Olga Boyko

South-Ukrainian National Pedagogical University
Staroportofrankovskaya Str. 26
65020, Odessa, Ukraine
Vyacheslav Pivovarchik
v.pivovarchik@paco.net

South-Ukrainian National Pedagogical University
Staroportofrankovskaya Str. 26
65020, Odessa, Ukraine
Received: October 19, 2013.
Accepted: June 18, 2014.

