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# DENSITY OF POLYNOMIALS IN THE $L^{2}$ SPACE ON THE REAL AND THE IMAGINARY AXES AND IN A SOBOLEV SPACE 

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Abstract. In this paper we consider an $L^{2}$ type space of scalar functions $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ which can be, in particular, the usual $L^{2}$ space of scalar functions on $\mathbb{R} \cup i \mathbb{R}$. We find conditions for density of polynomials in this space using a connection with the $L^{2}$ space of square-integrable matrix-valued functions on $\mathbb{R}$ with respect to a non-negative Hermitian matrix measure. The completness of $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ is also established.

1. Introduction. Let $\mathbb{C}_{n \times n}$ be the algebra of all $n \times n$ matrices with complex elements and $\mathbb{C}_{n \times n}^{\geq}$be the cone of all non-negative Hermitian matrices from $\mathbb{C}_{n \times n}, n \in \mathbb{N}$. Denote by $\mathfrak{B}(\mathbb{R})$ the Borel subsets of $\mathbb{R}$. By a $\mathbb{C}_{n \times n}^{\geq}$-valued measure on $\mathfrak{B}(\mathbb{R})$ we will mean a $\sigma$-additive function from $\mathfrak{B}(\mathbb{R})$ into $\mathbb{C}_{n \times n}^{\geq}$.
[^0]For each $\mathbb{C}_{n \times n}^{\geq}$-valued measure $M$ on $\mathfrak{B}(\mathbb{R})$ one can define a non-decreasing leftcontinuous (in the matrix norm) $\mathbb{C}_{n \times n}$-valued function (distribution function), see [1]:

$$
M(x):=\left\{\begin{array}{cc}
M([0, x)), & x>0  \tag{1}\\
0, & x=0 \\
M([x, 0)), & x<0
\end{array} .\right.
$$

Conversely, each non-decreasing left-continuous $\mathbb{C}_{n \times n}$-valued function $M(x), x \in$ $\mathbb{R}$, defines a $\mathbb{C}_{n \times n}^{\geq}$-valued measure on $\mathfrak{B}(\mathbb{R})[1]$.

Let $M$ be a $\mathbb{C}_{n \times n}^{\geq}$-valued measure on $\mathfrak{B}(\mathbb{R})$ with finite matrix moments

$$
\begin{equation*}
S_{k}=\int_{\mathbb{R}} x^{k} d M, \quad k \in \mathbb{Z}_{+} \tag{2}
\end{equation*}
$$

The set of all $\mathbb{C}_{n \times n}^{\geq}$-valued measures on $\mathfrak{B}(\mathbb{R})$ with moments $\left\{S_{k}\right\}_{k=0}^{\infty}$ we denote by $V$. The problem of the description of $V$ if one knows moments $\left\{S_{k}\right\}_{k=0}^{\infty}$ is the matrix Hamburger moment problem [14, p. 52]. This problem is called determinate if $V$ consists of a unique measure and indeterminate in the opposite case.

Let $L_{M}^{2}=L_{M}^{2}(\mathbb{R})$ be the space of square-integrable with respect to $M$ $\mathbb{C}_{n \times n}$-valued functions (see [12]). Denote by $\mathbb{P}_{n \times n}$ the set of all $(n \times n)$ matrices whose entries are complex polynomials. Elements of $\mathbb{P}_{n \times n}$ we call matrix polynomials. The set of all complex polynomials we denote by $\mathbb{P}\left(=\mathbb{P}_{1 \times 1}\right)$. Is the set $\mathbb{P}_{n \times n}$ dense in $L_{M}^{2}$ ? The theorem of Riesz states that in the scalar case $(n=1)$ polynomials will be dense if and only if $M$ is a $N$-extremal solution of moment problem (2) [2, p. 59]. A similar situation is in the case of an arbitrary $n$. We suppose that for every $P(x)=x^{k} I_{n}+A_{k-1} x^{k-1}+A_{k-2} x^{k-2}+\cdots+A_{0} \in \mathbb{P}_{n \times n}$, where $A_{i} \in \mathbb{C}_{n \times n}, \quad i=0,1, \ldots, k-1 ; k \in \mathbb{Z}_{+} ; I_{n}=\left(\delta_{i, j}\right)_{i, j=1}^{n}$, the following condition holds true:

$$
\begin{equation*}
\int_{\mathbb{R}} P(x) d M P^{*}(x) \quad \text { is an invertible matrix. } \tag{3}
\end{equation*}
$$

Applying the process of pseudo-orthogonalization to the sequence $I_{n}, x I_{n}, x^{2} I_{n}, \ldots$, (see [3, pp. 577-578]) we obtain a sequence of orthonormal matrix polynomials $\left\{P_{k}(x)\right\}_{k=0}^{\infty}$ (where $P_{k}(x)$ has degree $\left.n\right)$ :

$$
\begin{equation*}
\int_{\mathbb{R}} P_{k}(x) d M P_{l}^{*}(x)=I_{n} \delta_{k, l}, \quad k, l \in \mathbb{Z}_{+} \tag{4}
\end{equation*}
$$

For these polynomials one can construct a generalized Jacobi matrix $J$ which is a $(n \times n)$ block tridiagonal matrix. The corresponding linear operator $\mathbf{J}$ in $l^{2}=\left\{\left(x_{l}\right)_{l \in \mathbb{Z}_{+}}: x_{l} \in \mathbb{C}, \sum_{k=0}^{\infty}\left|x_{k}\right|^{2}<\infty\right\}$ is symmetric. Let $\left(m_{-}, m_{+}\right)$be its deficiency indices. The matrix Hamburger moment problem is called completely indeterminate if $m_{+}=m_{-}=n$. In this case $\mathbb{P}_{n \times n}$ will be dense in $L_{M}^{2}$ iff (see [4, Theorem 1.1, p. 249]):

$$
\begin{align*}
& \int_{\mathbb{R}} \frac{d M(x)}{x-\lambda}=  \tag{5}\\
& =-\left\{C^{*}(\lambda)[I+U]-i A^{*}(\lambda)[I-U]\right\}\left\{D^{*}(\lambda)[I+U]-i B^{*}(\lambda)[I-U]\right\}^{-1} \\
& \quad \lambda \in \mathbb{C} \backslash \mathbb{R}
\end{align*}
$$

where $U$ is a constant unitary matrix and $A(\cdot), B(\cdot), C(\cdot), D(\cdot)$ are holomorphic matrix functions which can be computed explicitly using moments $\left\{S_{k}\right\}_{k=0}^{\infty}$.

Measures which satisfy condition (5) are called $N$-extremal.
Notice also that in the scalar case $(n=1)$ Riesz proved that polynomials are dense in $L_{M}^{2}$ iff the measure $(1+|x|)^{-2} d M(x)$ is determinate (i.e. defines a determinate moment problem (2)) [5]. Consequently, several criteria for the determinacy of a measure can be used to study the question of density of polynomials [6].

The main purpose of our present investigation is to study the question of density of $\mathbb{P}$ in a general $L^{2}$ space which in special cases can be the $L^{2}$ space on $\mathbb{R} \cup i \mathbb{R}$ or a discrete Sobolev type space related to the discrete Sobolev orthogonal polynomials [7]. These general $L^{2}$ spaces are related to orthogonal polynomials on rays [8].

If $f=f(\lambda)$ is a $\mathbb{C}$-valued function on $\mathbb{R} \cup i \mathbb{R}$, we set

$$
\overrightarrow{f_{s}}(\lambda)=(f(\lambda), f(-\lambda)), \quad \lambda \in \mathbb{R} \cup i \mathbb{R}
$$

If there exist derivatives of $f(\lambda)$ at $\lambda=0$ in directions of $\mathbb{R}, i \mathbb{R}$ and they coincide, we denote by $f^{\prime}(0)$ their common value.

If there exists $f^{\prime}(0)$, we set $\vec{f}_{d}=\left(f(0), f^{\prime}(0)\right)$.
Let $M$ and $M_{i}$ be $\mathbb{C}_{n \times n}^{\geq}$-valued measures on $\mathfrak{B}(\mathbb{R})$ and $M(x), M_{i}(x)$ be their distribution functions, respectively. Denote by $\mathfrak{B}(i \mathbb{R})$ the set $\{i A, A \in$ $\mathfrak{B}(\mathbb{R})\}$. For a set $\hat{A} \in \mathfrak{B}(i \mathbb{R})$ we set $M_{I}(\hat{A})=M_{i}\left(\frac{1}{i} \hat{A}\right)$. Then $M_{I}$ will be a $\sigma$-additive function from $\mathfrak{B}(i \mathbb{R})$ into $\mathbb{C}_{n \times n}^{\geq}$. Such functions we call $\mathbb{C}_{n \times n}^{\geq}$-valued measures on $\mathfrak{B}(i \mathbb{R})$. The function

$$
M_{I}(x):=\left\{\begin{array}{cc}
M([0, x)), & \frac{1}{i} x>0  \tag{6}\\
0, & x=0 \\
M([x, 0)), & \frac{1}{i} x<0
\end{array}\right.
$$

will be called a distribution function of $M_{I}$.
The Radon-Nikodym derivatives of the measures $M, M_{I}$ we will denote by $M^{\prime}$ and $M_{I}^{\prime}$, respectively [12]. The corresponding trace measures we denote by $\tau_{M}$ and $\tau_{M_{I}}$, respectively. We set $\mathbf{M}=\left(M, M_{I}\right)$. Let $A \in \mathbb{C}_{n \times n}^{\geq}$.

Definition 1.1. Consider a set $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ of $\mathbb{C}$-valued functions on $\mathbb{R} \cup i \mathbb{R}, f \in L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R}) \Leftrightarrow$
(i) $f$ is $\tau_{M}$-measurable on $\mathbb{R}$, $f$ is $\tau_{M_{I}}$-measurable on $i \mathbb{R}$;
(ii) $\vec{f}_{s}(\lambda) M^{\prime} \vec{f}_{s}^{*}(\lambda) \in L_{\tau_{M}}^{1}(\mathbb{R}), \vec{f}_{s}(\lambda) M_{I}^{\prime}(\lambda) \overrightarrow{f_{s}^{*}}(\lambda) \in L_{\tau_{M_{I}}}^{1}(i \mathbb{R})$;
(iii) if $A \neq 0$ then $f(0)$ is finite and there exists $f^{\prime}(0)$.

Obviously, $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ is a linear vector space. We define the inner product in $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ in the following way:

$$
\begin{align*}
&\langle f, g\rangle_{2 ; \mathbf{M}, A}=\int_{\mathbb{R}} \vec{f}_{s}(\lambda) M^{\prime}(\lambda) \vec{g}_{s}^{*}(\lambda) d \tau_{M}+\int_{i \mathbb{R}} \vec{f}_{s}(\lambda) M_{I}^{\prime}(\lambda) \vec{g}_{s}^{*}(\lambda) d \tau_{M_{I}}+  \tag{7}\\
&+\vec{f}_{d} A \vec{g}_{d}^{*}, \quad f, g \in L_{\mathbf{M}, A}^{2}
\end{align*}
$$

For $f \in L_{\mathbf{M}, A}^{2}$ the norm is defined as

$$
\begin{equation*}
\|f\|_{2 ; \mathbf{M}, A}^{2}=\sqrt{\langle f, f\rangle_{2 ; \mathbf{M}, A}} \tag{8}
\end{equation*}
$$

As usual, we shall consider equivalence classes of such functions with respect to $\langle\cdot, \cdot\rangle_{2 ; \mathbf{M}, A}$. Then $L_{\mathbf{M}, A}^{2}$ becomes a unitary space with inner product $\langle\cdot, \cdot\rangle_{2 ; \mathbf{M}, A}$.

To study density of $\mathbb{P}$ in $L_{\mathbf{M}, A}^{2}$ we will construct a non-decreasing $\mathbb{C}_{2 \times 2^{-}}$ valued function $M_{1}$ on $\mathbb{R}$ depending on $\mathbf{M}, A$ and establish an isometric isomorphism $\mathbf{U}$ between $L_{\mathbf{M}, A}^{2}$ and the Hilbert space $L_{v}^{2}\left(M_{1}\right)$ of $\mathbb{C}^{2}$-valued functions square integrable with respect to $M_{1}$, see [12]. Then we use a connection between $L_{v}^{2}\left(M_{1}\right)$ and $L_{M_{1}}^{2}$ and apply the mentioned above results about density to $L_{M_{1}}^{2}$. As a by-product, we obtain that the space $L_{\mathbf{M}, A}^{2}$ is complete.

Notations. As usual, we denote by $\mathbb{R}, \mathbb{C}, \mathbb{N}, \mathbb{Z}, \mathbb{Z}_{+}$the sets of real, complex, positive integer, integer, non-negative integer numbers, respectively, and $i \mathbb{R}$ stands for the imaginary axis in the complex plane. Besides the definitions given above we should note the following notations. The set of complex vectors $\left(c_{1}, c_{2}\right), c_{1}, c_{2} \in \mathbb{C}$ we denote by $\mathbb{C}^{2}$. We identify $\mathbb{C}_{1 \times 1}$ with $\mathbb{C}$ and $\mathbb{C}_{1 \times 1}^{\geq}$with the set of non-negative real numbers. If $A \in \mathbb{C}_{n \times n}$ then $A^{*}$ stands for its adjoint, $n \in \mathbb{N}$. If $A \in \mathbb{C}_{n \times n}$ is non-degenerate, $A^{-1}$ means its inverse. The set of all vectors $\left(p_{1}, p_{2}\right)$ where $p_{1}, p_{2} \in \mathbb{P}$ we denote by $\mathbb{P}_{v}$. These vectors will be called vector polynomials.
2. Density of polynomials. Let $C_{0}$ be a subset of $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ which consists of continuous functions which have the first derivative at zero (in the directions of $\mathbb{R}, i \mathbb{R}$ and they coincide). For $f, g \in C_{0}$ we can write

$$
\begin{align*}
& \langle f, g\rangle_{2 ; \mathbf{M}, A}=\int_{\mathbb{R}} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \tau_{M}+\int_{i \mathbb{R}} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \tau_{M_{I}}+\vec{f}_{d} A \vec{g}_{d}^{*}= \\
& =\int_{\mathbb{R} \backslash\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \tau_{M}+\int_{(i \mathbb{R}) \backslash\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \tau_{M_{I}}+\int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \tau_{M}+ \\
& 9) \quad+\int_{\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \tau_{M_{I}}+\vec{f}_{d} A \vec{g}_{d}^{*}, \tag{9}
\end{align*}
$$

where $\varepsilon_{i}>0, i=\overline{1,4}$.
Define

$$
\hat{\tau}_{M}(\lambda)= \begin{cases}\tau_{M}(\lambda), & \lambda \leq 0 \\ \tau_{M}(\lambda)-\Delta \tau_{M}(0), & \lambda>0\end{cases}
$$

where $\Delta \tau_{M}(0):=\tau_{M}(+0)-\tau_{M}(0)$ is a jump of $\tau_{M}(\lambda)$ at zero;

$$
\hat{\tau}_{M_{I}}(\lambda)= \begin{cases}\tau_{M_{I}}(\lambda), & \lambda \in(-i \infty, 0] \\ \tau_{M_{I}}(\lambda)-\Delta_{i} \tau_{M_{I}}(0), & \lambda \in(0, i \infty)\end{cases}
$$

where $\Delta_{i} \tau_{M_{I}}(0):=\tau_{M_{I}}(+i 0)-\tau_{M_{I}}(0)=\lim _{\varepsilon \rightarrow+0} \tau_{M_{I}}(\varepsilon i)-\tau_{M_{I}}(0)$, is a jump of $\tau_{M_{I}}(\lambda)$ at zero.

Then

$$
\int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \tau_{M}=\int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \hat{\tau}_{M}=\vec{f}_{s}(0) M^{\prime}(0) \vec{g}_{s}^{*}(0) \Delta \tau_{M}(0)+
$$

$$
=\int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \hat{\tau}_{M}+f(0) \overline{g(0)} \tilde{a}
$$

where $\tilde{a}:=(1,1) M^{\prime}(0)\binom{1}{1} \Delta \tau_{M}(0)$.
Since functions $f, g$ and entries of $M^{\prime}(\lambda)$ are $\tau_{M^{-}}$-a.e. bounded functions in any neighbourhood of 0 , we have

$$
\int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \hat{\tau}_{M} \rightarrow 0, \quad \text { as } \varepsilon_{1}, \varepsilon_{2} \rightarrow 0
$$

Consequently,

$$
\begin{equation*}
\int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \tau_{M} \rightarrow f(0) \overline{g(0)} \tilde{a}, \quad \text { as } \varepsilon_{1}, \varepsilon_{2} \rightarrow 0 \tag{10}
\end{equation*}
$$

In a similar manner we can write

$$
\begin{aligned}
\int_{\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \tau_{M_{I}} & =\int_{\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \hat{\tau}_{M_{I}}+\vec{f}_{s}(0) M_{I}^{\prime}(0) \vec{g}_{s}^{*}(0) \Delta_{i} \tau_{M_{I}}(0)= \\
& =\int_{\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \hat{\tau}_{M}+f(0) \overline{g(0)} \tilde{a}_{I}
\end{aligned}
$$

where $\tilde{a}_{I}:=(1,1) M_{I}^{\prime}(0)\binom{1}{1} \Delta_{i} \tau_{M_{i}}(0)$.
Since functions $f, g$ and entries of $M_{I}^{\prime}(\lambda)$ are $\tau_{M}$-a.e. bounded functions in any neighbourhood of 0 , we have

$$
\int_{\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \hat{\tau}_{M_{I}} \rightarrow 0, \quad \text { as } \varepsilon_{3}, \varepsilon_{4} \rightarrow 0
$$

Thus, we have

$$
\begin{equation*}
\int_{\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \tau_{M_{I}} \rightarrow f(0) \overline{g(0)} \tilde{a}_{I}, \quad \text { as } \varepsilon_{3}, \varepsilon_{4} \rightarrow 0 \tag{11}
\end{equation*}
$$

Set $J_{\lambda}:=\left(\begin{array}{cc}1 & \frac{1}{\lambda} \\ 1 & -\frac{1}{\lambda}\end{array}\right)$. Notice that $J_{\lambda}^{-1}=\frac{1}{2}\left(\begin{array}{cc}1 & 1 \\ \lambda & -\lambda\end{array}\right)$.
Then we can write

$$
\int_{\mathbb{R} \backslash\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} M^{\prime} \vec{g}_{s}^{*} d \tau_{M}=\int_{\mathbb{R} \backslash\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} J_{\lambda} J_{\lambda}^{-1} M^{\prime}\left(J_{\lambda}^{*}\right)^{-1} J_{\lambda}^{*} \vec{g}_{s}^{*} d \tau_{M}=
$$

$$
\begin{gather*}
=\int_{\mathbb{R} \backslash\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{s} J_{\lambda} \tilde{M}(\lambda) J_{\lambda}^{*} \vec{g}_{s}^{*} d \tau_{M} \rightarrow \\
\rightarrow \int_{\mathbb{R}}^{\prime} \vec{f}_{s} J_{\lambda} \tilde{M}(\lambda) J_{\lambda}^{*} \vec{g}_{s}^{*} d \tau_{M}, \quad \text { as } \varepsilon_{1}, \varepsilon_{2} \rightarrow 0, \tag{12}
\end{gather*}
$$

where $\tilde{M}(\lambda):=J_{\lambda}^{-1} M^{\prime}(\lambda)\left(J_{\lambda}^{*}\right)^{-1}=\frac{1}{4}\left(\begin{array}{cc}1 & 1 \\ \lambda & -\lambda\end{array}\right) M^{\prime}(\lambda)\left(\begin{array}{cc}1 & \bar{\lambda} \\ 1 & -\bar{\lambda}\end{array}\right)$ and the prime in $\int^{\prime}$ means that the integral is understood in the improper sense at zero.

In an analogous manner we write

$$
\begin{aligned}
\int_{(i \mathbb{R}) \backslash\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} M_{I}^{\prime} \vec{g}_{s}^{*} d \tau_{M_{I}} & =\int_{(i \mathbb{R}) \backslash\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} J_{\lambda} J_{\lambda}^{-1} M_{I}^{\prime}\left(J_{\lambda}^{*}\right)^{-1} J_{\lambda}^{*} \vec{g}_{s}^{*} d \tau_{M_{I}}= \\
& =\int_{(i \mathbb{R}) \backslash\left(-i \varepsilon_{3}, i \varepsilon_{4}\right)} \vec{f}_{s} J_{\lambda} \tilde{M}_{I}(\lambda) J_{\lambda}^{*} \vec{g}_{s}^{*} d \tau_{M_{I}} \rightarrow \\
& \rightarrow \int_{i \mathbb{R}}^{\prime} \vec{f}_{s} J_{\lambda} \tilde{M}_{I}(\lambda) J_{\lambda}^{*} \vec{g}_{s}^{*} d \tau_{M_{I}}, \quad \text { as } \varepsilon_{3}, \varepsilon_{4} \rightarrow 0
\end{aligned}
$$

where $\tilde{M}_{I}(\lambda):=J_{\lambda}^{-1} M_{I}^{\prime}(\lambda)\left(J_{\lambda}^{*}\right)^{-1}=\frac{1}{4}\left(\begin{array}{cc}1 & 1 \\ \lambda & -\lambda\end{array}\right) M_{I}^{\prime}(\lambda)\left(\begin{array}{cc}1 & \bar{\lambda} \\ 1 & -\bar{\lambda}\end{array}\right)$ and the prime in $\int^{\prime}$ means that the integral is understood in the improper sense at zero.

Passing to the limit in (9) and using (10), (11), (12), (13) we get

$$
\begin{align*}
\langle f, g\rangle_{2 ; \mathbf{M}, A}= & \int_{\mathbb{R}}^{\prime} \vec{f}_{s} J_{\lambda} \tilde{M} J_{\lambda}^{*} \vec{g}_{s}^{*} d \tau_{M}+\int_{i \mathbb{R}}^{\prime} \vec{f}_{s} J_{\lambda} \tilde{M}_{I} J_{\lambda}^{*} \vec{g}_{s}^{*} d \tau_{M_{I}}+ \\
& +\vec{f}_{d}\left(A+\left(\begin{array}{cc}
\tilde{a}+\tilde{a}_{I} & 0 \\
0 & 0
\end{array}\right)\right) \vec{g}_{d}^{*} \tag{14}
\end{align*}
$$

For $f \in C_{0}$ we set

$$
\begin{aligned}
f^{+}(\lambda) & :=\frac{f(\lambda)+f(-\lambda)}{2}, \lambda \in \mathbb{R} \cup i \mathbb{R} \\
f^{-}(\lambda) & :=\frac{f(\lambda)-f(-\lambda)}{2 \lambda}, \lambda \in(\mathbb{R} \cup i \mathbb{R}) \backslash\{0\} .
\end{aligned}
$$

Functions $f^{+}, f^{-}$are continuous. Moreover,

$$
\begin{aligned}
\lim _{\lambda \rightarrow 0} f^{-}(\lambda)= & \lim _{\lambda \rightarrow 0} \frac{f(\lambda)-f(-\lambda)}{2 \lambda}=\lim _{\lambda \rightarrow 0} \frac{f(\lambda)-f(0)}{2 \lambda}+ \\
& +\lim _{\lambda \rightarrow 0} \frac{f(-\lambda)-f(0)}{-2 \lambda}=f^{\prime}(0)
\end{aligned}
$$

and also $\lim _{\lambda \rightarrow i 0} f^{-}(\lambda)=f^{\prime}(0)$. Then we set

$$
f^{-}(0):=f^{\prime}(0)
$$

and get a continuous function $f^{-}$on $\mathbb{R} \cup i \mathbb{R}$.
We set

$$
\vec{f}_{p}(\lambda):=\left(f^{+}(\lambda), f^{-}(\lambda)\right), \lambda \in \mathbb{R} \cup i \mathbb{R}
$$

Using these definitions we can rewrite (14) in the following form:

$$
\begin{align*}
& \langle f, g\rangle_{2 ; \mathbf{M}, A}=4 \int_{\mathbb{R}}^{\prime} \overrightarrow{f_{p}} \tilde{M} \vec{g}_{p}^{*} d \tau_{M}+4 \int_{i \mathbb{R}}^{\prime} \vec{f}_{p} \tilde{M}_{I} \vec{g}_{p}^{*} d \tau_{M_{I}}+  \tag{15}\\
& \quad+\vec{f}_{d}\left(A+\left(\begin{array}{cc}
\tilde{a}+\tilde{a}_{I} & 0 \\
0 & 0
\end{array}\right)\right) \vec{g}_{d}
\end{align*}
$$

Consider the following integrals:

$$
I_{1}:=4 \int_{\mathbb{R}} \vec{f}_{p} \tilde{M} \vec{g}_{p}^{*} d \tau_{M}, \quad I_{1, I}:=4 \int_{i \mathbb{R}} \vec{f}_{p} \tilde{M}_{I} \vec{g}_{p}^{*} d \tau_{M_{I}}
$$

which exist in view of the definitions of $\tilde{M}$ and $\tilde{M}_{I}$.
Let us study the integral $I_{1}$ :

$$
\begin{equation*}
I_{1}=4 \int_{\mathbb{R} \backslash\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \overrightarrow{f_{p}} \tilde{M} \vec{g}_{p}^{*} d \tau_{M}+4 \int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \overrightarrow{f_{p}} \tilde{M} \vec{g}_{p}^{*} d \tau_{M} \tag{16}
\end{equation*}
$$

where $\varepsilon_{1}, \varepsilon_{2}>0$.
We can write

$$
4 \int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \vec{f}_{p} \tilde{M} \vec{g}_{p}^{*} d \tau_{M}=4 \int_{\left(-\varepsilon_{1}, \varepsilon_{2}\right)} \overrightarrow{f_{p}} \tilde{M} \vec{g}_{p}^{*} d \hat{\tau}_{M}+4 \vec{f}_{p}(0) \tilde{M}(0) \vec{g}_{p}^{*}(0) \Delta \tau_{M}(0)
$$

Since functions $f^{+}, f^{-}, g^{+}, g^{-}$and entries of $\tilde{M}$ are $\tau_{M^{-}}$a.e. bounded in any neighbourhood of 0 , the first addend on the right of the latter equality tends to zero as $\varepsilon_{1}, \varepsilon_{2} \rightarrow 0$.

The first addend on the right-hand side of (16) tends to $4 \int_{\mathbb{R}}^{\prime} \vec{f}_{p} \tilde{M} \vec{g}_{p}^{*} d \tau_{M}$, as $\varepsilon_{1}, \varepsilon_{2} \rightarrow 0$, by the definition of the improper integral.

Consequently, passing to the limit in (16) we get

$$
I_{1}=4 \int_{\mathbb{R}}^{\prime} \overrightarrow{f_{p}} \tilde{M} \vec{g}_{p}^{*} d \tau_{M}+4 \vec{f}_{p}(0) \tilde{M}(0) \vec{g}_{p}^{*}(0) \Delta \tau_{M}(0)
$$

In a similar way we get that

$$
I_{1, I}=4 \int_{i \mathbb{R}}^{\prime} \vec{f}_{p} \tilde{M}_{I} \vec{g}_{p}^{*} d \tau_{M_{I}}+4 \vec{f}_{p}(0) \tilde{M}_{I}(0) \vec{g}_{p}^{*}(0) \Delta_{i} \tau_{M_{I}}(0)
$$

Therefore

$$
\begin{aligned}
& 4 \int_{\mathbb{R}}^{\prime} \vec{f}_{p} \tilde{M} \vec{g}_{p}^{*} d \tau_{M}+4 \int_{i \mathbb{R}}^{\prime} \vec{f}_{p} \tilde{M}_{I} \vec{g}_{p}^{*} d \tau_{M_{I}}=I_{1}+I_{1, I}- \\
& -4 \vec{f}_{p}(0) \tilde{M}(0) \vec{g}_{p}^{*}(0) \Delta \tau_{M}(0)-4 \vec{f}_{p}(0) \tilde{M}_{I}(0) \vec{g}_{p}^{*}(0) \Delta_{i} \tau_{M_{I}}(0)= \\
& =I_{1}+I_{1, I}-4 \vec{f}_{d} \tilde{M}(0) \vec{g}_{d}^{*} \Delta \tau_{M}(0)-4 \vec{f}_{d} \tilde{M}_{I}(0) \vec{g}_{d}^{*} \Delta_{i} \tau_{M_{I}}(0) .
\end{aligned}
$$

Substituting the last expression in (15) we get

$$
\begin{equation*}
\langle f, g\rangle_{2 ; \mathbf{M}, A}=4 \int_{\mathbb{R}} \overrightarrow{f_{p}} \tilde{M} \vec{g}_{p}^{*} d \tau_{M}+4 \int_{i \mathbb{R}} \vec{f}_{p} \tilde{M}_{I} \vec{g}_{p}^{*} d \tau_{M_{I}}+\vec{f}_{d} A \vec{g}_{d}^{*} \tag{17}
\end{equation*}
$$

since

$$
\left(\begin{array}{cc}
\tilde{a}+\tilde{a}_{I} & 0 \\
0 & 0
\end{array}\right)-4 \tilde{M}(0) \Delta \tau_{M}(0)-4 \tilde{M}_{I}(0) \Delta_{i} \tau_{M_{I}}(0)=0
$$

Set

$$
\widehat{M}(\lambda)= \begin{cases}4 \int_{0}^{\lambda} \tilde{M}(\lambda) d \tau_{M}(\lambda), & \lambda \in \mathbb{R} \\ 4 \int_{0}^{\lambda} \tilde{M}_{I}(\lambda) d \tau_{M_{I}}(\lambda), & \lambda \in i \mathbb{R}\end{cases}
$$

Since functions $f^{+}, f^{-}, g^{+}, g^{-}$are continuous, we can rewrite (17) in the following form:

$$
\begin{equation*}
\langle f, g\rangle_{2 ; \mathbf{M}, A}=\int_{\mathbb{R}} \overrightarrow{f_{p}} d \widehat{M}(\lambda) \vec{g}_{p}^{*}+\int_{i \mathbb{R}} \overrightarrow{f_{p}} d \widehat{M}(\lambda) \vec{g}_{p}^{*}+\vec{f}_{d} A \vec{g}_{d}^{*} \tag{18}
\end{equation*}
$$

where integrals are understood as a sum of integrals of each scalar addend.
Denote the sum of integrals on the right by $I_{2}$. Then we can write

$$
\begin{equation*}
I_{2}=\int_{(-\infty, 0]} \vec{f}_{p} d \widehat{M}(\lambda) \vec{g}_{p}^{*}+\int_{[0, \infty)}+\cdots+\int_{(-i \infty, 0]}+\cdots+\int_{[0, i \infty)}+\cdots \tag{19}
\end{equation*}
$$

In the first two terms on the right we make the changes of variable: $x=\lambda^{2}, \lambda=$ $-\sqrt{x}$ and $x=\lambda^{2}, \lambda=\sqrt{x}$, respectively. Here $\sqrt{x}=|x|^{\frac{1}{2}}$.

In the last two terms on the right of (19) we make the changes of variable: $x=\lambda^{2}, \lambda=-\sqrt{x}$ and $x=\lambda^{2}, \lambda=\sqrt{x}$, respectively. Here $\sqrt{x}=|x|^{\frac{1}{2}} i$. Note that functions $f^{+}, f^{-}, g^{+}, g^{-}$are even and do not depend on the choice of a branch of the square root. Define the following matrix-function:

$$
\widehat{M}_{1}(x)=\left\{\begin{array}{ll}
\widehat{M}(-\sqrt{x})-\widehat{M}(\sqrt{x}), & x \leq 0  \tag{20}\\
\widehat{M}(\sqrt{x})-\widehat{M}(-\sqrt{x})+A, & x>0
\end{array} .\right.
$$

This function is a non-decreasing $\mathbb{C}_{2 \times 2}$-valued function on $\mathbb{R}$. Then from (18) and (19) we get

$$
\langle f, g\rangle_{2 ; \mathbf{M}, A}=\int_{\mathbb{R}} \vec{f}_{p}(\sqrt{x}) d \widehat{M}_{1}(x) \vec{g}_{p}^{*}(\sqrt{x})
$$

We set $M_{1}(x):=\widehat{M}_{1}(x-0), x \in \mathbb{R}$, to get a left-continuous function and

$$
\begin{equation*}
\langle f, g\rangle_{2 ; \mathbf{M}, A}=\int_{\mathbb{R}} \vec{f}_{p}(\sqrt{x}) d M_{1}(x) \vec{g}_{p}^{*}(\sqrt{x}) \tag{21}
\end{equation*}
$$

The function $M_{1}(x)$ defines a $\mathbb{C}_{2 \times 2}^{\geq}$-valued measure $M_{1}$ on $\mathfrak{B}(\mathbb{R})$. The space of square-integrable vector functions $\vec{v}=\left(v_{1}(x), v_{2}(x)\right)$ on $\mathbb{R}$ with respect to the measure $M_{1}$ we shall denote $L_{v}^{2}\left(M_{1}\right)$. The space of square-integrable matrixfunctions $v=\left(\begin{array}{ll}v_{1}(x) & v_{2}(x) \\ v_{3}(x) & v_{4}(x)\end{array}\right)$ on $\mathbb{R}$ with respect to the measure $M_{1}$ we denote $L_{M_{1}}^{2}$. For the general definition of such spaces we refer to [12]. Notice that in fact

$$
\begin{equation*}
L_{M_{1}}^{2}=L_{v}^{2}\left(M_{1}\right) \oplus L_{v}^{2}\left(M_{1}\right) \tag{22}
\end{equation*}
$$

according to the map

$$
L_{M_{1}}^{2} \ni\left(\begin{array}{ll}
v_{1} & v_{2} \\
v_{3} & v_{4}
\end{array}\right) \rightarrow\left(v_{1}, v_{2}\right) \oplus\left(v_{3}, v_{4}\right) \in L_{v}^{2}\left(M_{1}\right) \oplus L_{v}^{2}\left(M_{1}\right)
$$

From this fact it follows a corollary:
Corollary 2.1. The set $\mathbb{P}_{v}$ is dense in $L_{v}^{2}\left(M_{1}\right)$ iff the set $\mathbb{P}_{2 \times 2}$ is dense in $L^{2}\left(M_{1}\right)$.

By $\langle\cdot, \cdot\rangle_{2, v ; M_{1}}$ and $\|\cdot\|_{2, v ; M_{1}}$ we denote the inner product and the norm in $L_{v}^{2}\left(M_{1}\right)$, respectively. We can define an operator $\mathbf{V}: C_{0} \rightarrow L_{v}^{2}\left(M_{1}\right)$ :

$$
\begin{equation*}
(\mathbf{V} f)(x)=\vec{f}_{p}(\sqrt{x}), x \in \mathbb{R}, f \in C_{0} \tag{23}
\end{equation*}
$$

From (21) it follows that

$$
\begin{equation*}
\langle\mathbf{V} f, \mathbf{V} g\rangle_{2, v ; M_{1}}=\langle f, g\rangle_{2 ; \mathbf{M}, A}, f, g \in C_{0} \tag{24}
\end{equation*}
$$

Let $R_{0}$ be a subset of $L_{v}^{2}\left(M_{1}\right)$ which consists of functions $\vec{v}(x)=\left(v_{1}(x), v_{2}(x)\right)$, such that $v_{1}, v_{2}$ are continuous and $w(\lambda):=v_{1}\left(\lambda^{2}\right), \lambda \in \mathbb{R} \cup i \mathbb{R}$ has the first derivative at zero (in the directions of $\mathbb{R}$ and $i \mathbb{R}$ and they coincide). We shall show that $\mathbf{V} C_{0}=R_{0}$.

Notice that for $f \in C_{0}$ the functions $f^{+}(\sqrt{x}), f^{-}(\sqrt{x})$ are continuous and the function $f^{+}(\lambda)=\frac{1}{2}(f(\lambda)+f(-\lambda))$ has the derivative $\left.\frac{d}{d \lambda} f^{+}\right|_{\lambda=0}(=0)$. Hence, $\mathbf{V} C_{0} \subseteq R_{0}$. On the other hand, take $\vec{v}(x)=\left(v_{1}(x), v_{2}(x)\right) \in R_{0}$ and set $f_{v}(\lambda):=v_{1}\left(\lambda^{2}\right)+\lambda v_{2}\left(\lambda^{2}\right), \lambda \in \mathbb{R} \cup i \mathbb{R}$. Then $f_{v}^{+}(\lambda)=v_{1}\left(\lambda^{2}\right), f_{v}^{-}(\lambda)=v_{2}\left(\lambda^{2}\right)$ and

$$
\begin{equation*}
\left(f_{v}^{+}(\sqrt{x}), f_{v}^{-}(\sqrt{x})\right)=\left(v_{1}(x), v_{2}(x)\right), x \in \mathbb{R} \tag{25}
\end{equation*}
$$

Moreover,

$$
\begin{gathered}
\exists f_{v}^{\prime}(0)=\lim _{\Delta \lambda \rightarrow 0} \frac{f_{v}(\Delta \lambda)-f_{v}(0)}{\Delta \lambda}= \\
=\lim _{\Delta \lambda \rightarrow 0} \frac{v_{1}\left((\Delta \lambda)^{2}\right)+\Delta \lambda v_{2}\left((\Delta \lambda)^{2}\right)-v_{1}(0)}{\Delta \lambda}=v_{2}(0)+\left.\left(v_{1}\left(\lambda^{2}\right)\right)^{\prime}\right|_{\lambda=0},
\end{gathered}
$$

Using (25) we can write

$$
\|\vec{v}\|_{2, v ; M_{1}}=\int_{\mathbb{R}}\left(\overrightarrow{f_{v}}\right)_{p}(\sqrt{x}) d M_{1}(x)\left(\overrightarrow{f_{v}}\right)_{p}^{*}(\sqrt{x}) .
$$

Moving backwards from (21) to (9) we obtain

$$
\|\vec{v}\|_{2, v ; M_{1}}=\left\|f_{v}\right\|_{2 ; \mathbf{M}, A}<\infty
$$

Hence, $f_{v} \in C_{0}$. Moreover,

$$
\begin{equation*}
\left(\mathbf{V} f_{v}\right)(x)=\left(\overrightarrow{f_{v}}\right)_{p}(\sqrt{x})=\vec{v}(x) \tag{26}
\end{equation*}
$$

Consequently, $\mathrm{VC}_{0}=R_{0}$.
Define an operator $\mathbf{G}: R_{0} \rightarrow C_{0}$ :

$$
\begin{equation*}
(\mathbf{G} \vec{v})(\lambda)=v_{1}\left(\lambda^{2}\right)+\lambda v_{2}\left(\lambda^{2}\right), \lambda \in \mathbb{R} \cup i \mathbb{R}, \vec{v}=\left(v_{1}, v_{2}\right) \in R_{0} \tag{27}
\end{equation*}
$$

From (27) it follows that $\mathbf{V G} \vec{v}=\vec{v}, \vec{v} \in R_{0}$. On the other hand, $\mathbf{G V} f=\mathbf{G} \vec{f}_{p}(\sqrt{x})=f^{+}(\lambda)+\lambda f^{-}(\lambda)=f, f \in C_{0}$. Hence, $\mathbf{G}=\mathbf{V}^{-1}$.

Lemma 2.2. The subset $C_{0}$ is dense in $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$.
Proof. 1) Case $A=0$. Here we shall follow the ideas of Kats in [13]. Take $f \in L_{\mathbf{M}, A}^{2}$. Set

$$
f_{N}(\lambda)=\left\{\begin{array}{ll}
f(\lambda), & \text { if }|f(\lambda)|<N,|f(-\lambda)|<N  \tag{28}\\
0, & \text { otherwise }
\end{array}, \lambda \in \mathbb{R} \cup i \mathbb{R}, N \in \mathbb{N}\right.
$$

Let $E_{N}:=\left\{\lambda \in \mathbb{R} \cup i \mathbb{R}: f_{N}(\lambda) \neq f(\lambda)\right\}$. Since $f(\lambda)$ has almost everywhere with respect to $\tau_{M}$ on $\mathbb{R}$ and with respect to $\tau_{M_{I}}$ on $i \mathbb{R}$ finite values, we get that

$$
\begin{equation*}
\tau_{M}\left(E_{N} \cap \mathbb{R}\right) \rightarrow 0 \text { and } \tau_{M_{I}}\left(E_{N} \cap(i \mathbb{R})\right) \rightarrow 0 \tag{29}
\end{equation*}
$$

as $N \rightarrow \infty$. We shall denote

$$
\begin{gathered}
\Psi(\lambda ; g):=\vec{g}_{s}(\lambda) M^{\prime}(\lambda) \vec{g}_{s}^{*}(\lambda), \quad \lambda \in \mathbb{R} \\
\Psi_{I}(\lambda ; g):=\vec{g}_{s}(\lambda) M_{I}^{\prime}(\lambda) \vec{g}_{s}^{*}(\lambda), \lambda \in i \mathbb{R}, \quad g \in L_{\mathbf{M}, A}^{2}
\end{gathered}
$$

Then

$$
\left\|f_{N}-f\right\|_{2 ; \mathbf{M}, A}^{2}=\int_{E_{N} \cap \mathbb{R}} \Psi(\lambda ; f) d \tau_{M}+\int_{E_{N} \cap(i \mathbb{R})} \Psi_{I}(\lambda ; f) d \tau_{M_{I}} \rightarrow 0
$$

as $N \rightarrow \infty$, since $f$ is integrable and (29) holds.
Hence, bounded functions are dense in $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$.
Consider a bounded function $f_{b}(\lambda) \in L_{\mathbf{M}, A}^{2}:\left|f_{b}(\lambda)\right| \leq N, N \in \mathbb{N}$. Take an arbitrary $\varepsilon: 0<\varepsilon<1$. Choose $K>0$ such that $\tau_{M}$ is continuous at points $-K, K$ and

$$
\left|\int_{\mathbb{R} \backslash[-K, K]} \Psi(\lambda ; f) d \tau_{M}\right|<\varepsilon
$$

Let $\tau_{M}(\lambda), \lambda \in \mathbb{R}$ be the distribution function of the measure $\tau_{M}$. Consider the function

$$
\hat{\tau}_{M}(\lambda):=\left\{\begin{array}{ll}
\tau_{M}(\lambda), & \lambda \leq 0 \\
\tau_{M}(\lambda)+1, & \lambda>0
\end{array} .\right.
$$

The function $\hat{\tau}_{M}(\lambda)$ defines a $\mathbb{C}^{\geq}$-valued measure $\hat{\tau}_{M}$ on $\mathfrak{B}(\mathbb{R})$. It is important for us that

$$
\begin{gather*}
\hat{\tau}_{M}(A)=\tau_{M}(A), \quad A \in \mathfrak{B}(\mathbb{R}): 0 \notin A  \tag{30}\\
\hat{\tau}_{M}(A)=\tau_{M}(A)+1 \geq 1, \quad A \in \mathfrak{B}(\mathbb{R}): 0 \in A \tag{31}
\end{gather*}
$$

Applying Luzin's theorem (see [9, pp. 227-229]) we can find a continuous function $g_{N}(\lambda)$ on $[-K, K]$ such that for $Q_{N}:=\left\{\lambda \in[-K, K]: g_{N}(\lambda) \neq f_{b}(\lambda)\right\}$ we have $\hat{\tau}_{M}\left(Q_{N}\right)<\varepsilon$. Also $g_{N}(\lambda)$ can be chosen such that $\left|g_{N}(\lambda)\right|<N$.

From (31) it follows that $0 \notin Q_{N}$, therefore $g_{N}(0)=f_{b}(0)$. From (31) it follows also that $\tau_{M}\left(Q_{N}\right)<\varepsilon$.

There exists $\delta>0$ such that $\tau_{M}([-K,-K-\delta])<\varepsilon$ and $\tau_{M}([K, K+\delta])<$ $\varepsilon$, since $\tau_{M}$ is continuous at points $-K, K$.

We continue the function $g_{N}(\lambda)$ on intervals $[-K,-K-\delta]$ and $[K, K+\delta]$ by intervals of a straight line $a \lambda+b$ connecting the known value of $g_{N}(\lambda)$ with 0 . Namely,

$$
g_{N}(\lambda):=\left\{\begin{array}{ll}
\frac{g_{N}(-K)}{\delta}(\lambda+K+\delta), & \lambda \in[-K,-K-\delta] \\
-\frac{g_{N}(K)}{\delta}(\lambda-K-\delta), & \lambda \in[K, K+\delta]
\end{array} .\right.
$$

Thus, we have

$$
\left|g_{N}(\lambda)\right| \leq\left|g_{N}(-K)\right|<N, \lambda \in[-K,-K-\delta]
$$

$$
\left|g_{N}(\lambda)\right| \leq\left|g_{N}(K)\right|<N, \lambda \in[K, K-\delta]
$$

On the rest of $\mathbb{R}$ we set $g_{N}(\lambda)=0$.
Consequently,

$$
\begin{aligned}
& \left|\int_{\mathbb{R}} \Psi\left(\lambda ; f-g_{N}\right) d \tau_{M}\right|<\varepsilon+\left|\int_{[-K-\delta,-K] \cup[K, K+\delta]} \Psi\left(\lambda ; f-g_{N}\right) d \tau_{M}\right|+ \\
& \quad+\left|\int_{[-K, K]} \Psi\left(\lambda ; f-g_{N}\right) d \tau_{M}\right|<\varepsilon+16 N^{2} \tau_{M}([-K-\delta,-K])+ \\
& \quad+16 N^{2} \tau_{M}([K, K+\delta])+16 N^{2} \tau_{M}\left(Q_{N}\right)<\left(48 N^{2}+1\right) \varepsilon
\end{aligned}
$$

In a similar way we can construct a continuous function $g_{N, I}(\lambda)$ on $i \mathbb{R}$ such that $g_{N, I}(0)=f_{b}(0)$ and

$$
\left|\int_{i \mathbb{R}} \Psi_{I}\left(\lambda ; f-g_{N, I}\right) d \tau_{M_{I}}\right|<\left(48 N^{2}+1\right) \varepsilon
$$

We set

$$
d_{N}(\lambda):= \begin{cases}g_{N}(\lambda), & \lambda \in \mathbb{R} \\ g_{N, I}(\lambda), & \lambda \in i \mathbb{R}\end{cases}
$$

It is not hard to see that $d_{N}(\lambda)$ will approximate $f_{b}(\lambda)$. Consequently, continuous functions are dense in $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$.

Take a continuous function $f_{c} \in L_{\mathbf{M}, A}^{2}$. For $a=-1 ; 1 ;-i ; i$ set

$$
\begin{equation*}
\hat{f}(\lambda)=(\lambda-a) g(\lambda)+f_{c}(a) \tag{32}
\end{equation*}
$$

where

$$
g(\lambda)=\frac{1}{a}\left(f_{c}(a)-f_{c}(0)\right)+\lambda\left(\frac{1}{a^{2}}\left(f_{c}(a)-f_{c}(0)\right)-\frac{d}{a}\right), \lambda \in[0, a]
$$

and $d$ is an arbitrary complex number. It is not hard to see that $\hat{f}(\lambda)$ is a continuos function on $[-1,1] \cup[-i, i]$ such that

$$
\begin{equation*}
\hat{f}(-1)=f_{c}(-1), \hat{f}(1)=f_{c}(1), \hat{f}(-i)=f_{c}(-i), \hat{f}(i)=f_{c}(i) \tag{33}
\end{equation*}
$$

$$
\begin{equation*}
\hat{f}(0)=f_{c}(0), \quad \exists \hat{f}^{\prime}(0)=d \tag{34}
\end{equation*}
$$

Since $f_{c}(\lambda) \rightarrow f_{c}(0)$ and $\hat{f}(\lambda) \rightarrow f_{c}(0)$ as $\lambda \rightarrow 0$, we get that

$$
\begin{gathered}
\forall \varepsilon>0 \exists \delta_{1}>0\left(\delta_{1}<1\right): \lambda \in[-1,1] \cup[-i, i],|\lambda| \leq \delta \Rightarrow \\
\left|\hat{f}(\lambda)-f_{c}(\lambda)\right|<\varepsilon .
\end{gathered}
$$

Set
(35) $\quad f_{1}(\lambda)= \begin{cases}\hat{f}(\lambda), & |\lambda|<\delta_{1} \\ f_{c}(\lambda)+\frac{|\lambda|-1}{\delta_{1}-1}\left(\hat{f}\left(\delta_{1}\right)-f_{c}\left(\delta_{1}\right)\right), & \delta_{1} \leq|\lambda| \leq 1 \quad, ~ \lambda \in \mathbb{R} \cup i \mathbb{R} . \\ f_{c}(\lambda), & |\lambda|>1\end{cases}$

Then $f_{1}(\lambda)$ is a continuous function on $\mathbb{R} \cup i \mathbb{R}$ which has the first derivative at zero and such that:

$$
\begin{equation*}
\left|f_{1}(\lambda)-f_{c}(\lambda)\right|<\varepsilon, \lambda \in[-1,1] \cup[-i, i] \tag{36}
\end{equation*}
$$

Then

$$
\begin{align*}
\left\|f_{1}-f_{c}\right\|_{2 ; \mathbf{M}, A}^{2}= & \int_{[-1,1]} \Psi\left(\lambda ; f_{1}-f_{c}\right) d \tau_{M}+\int_{[-i, i]} \Psi_{I}\left(\lambda ; f_{1}-f_{c}\right) d \tau_{M_{I}} \leq \\
& \leq 4 \varepsilon^{2} \tau_{M}([-1,1])+4 \varepsilon^{2} \tau_{M_{I}}([-i, i]) \tag{37}
\end{align*}
$$

Consequently, $C_{0}$ is dense in $L_{\mathbf{M}, A}^{2}$ in this case.
2) Case $A \neq 0$. Take a function $f \in L_{\mathbf{M}, A}^{2}$. Consider the space $L_{\mathbf{M}, 0}^{2}$. Evidently, $f \in L_{\mathbf{M}, 0}^{2}$ as well. Let us apply constructions from the previous case. Since $f(0)$ is finite, we can for an arbitrary $\varepsilon>0$ construct a bounded function $f_{N}(\lambda)(N>|f(0)|)$ such that $f_{N}(0)=f(0)$ and $\left\|f-f_{N}\right\|_{2 ; \mathbf{M}, 0}<\varepsilon$. Then we construct a continuous function $f_{c}(\lambda)$ such that $f_{c}(0)=f_{N}(0)=f(0)$ and $\left\|f_{N}-f_{c}\right\|_{2 ; \mathbf{M}, 0}<\varepsilon$. For the continuous function $f_{c}$ we construct a function $\hat{f}(\lambda)$ as in (32) with $d=f^{\prime}(0)$. Then we repeat considerations after (32). Define a function $f_{1}(\lambda)$ as in (35). The function $f_{1}$ will posess property (36) and (37) holds with $A=0$. Notice that $f_{1} \in L_{\mathbf{M}, A}^{2}$ by the definition of $L_{\mathbf{M}, A}^{2}$. Moreover, $f_{1} \in C_{0}$. Consequently,

$$
\left\|f-f_{1}\right\|_{2 ; \mathbf{M}, A}=\left\|f-f_{1}\right\|_{2 ; \mathbf{M}, 0} \leq\left\|f-f_{N}\right\|_{2 ; \mathbf{M}, 0}+\left\|f_{N}-f_{c}\right\|_{2 ; \mathbf{M}, 0}+
$$

$$
+\left\|f_{c}-f_{1}\right\|_{2 ; \mathbf{M}, 0} \leq 2 \varepsilon+4 \varepsilon^{2} \tau_{M}([-1,1])+4 \varepsilon^{2} \tau_{M_{I}}([-i, i])
$$

Hence, $C_{0}$ is dense in $L_{\mathbf{M}, A}^{2}$ in this case as well.
Lemma 2.2. The subset $R_{0}$ is dense in $L_{v}^{2}\left(M_{1}\right)$.
Proof. Since continuous matrix-functions are dense in $L^{2}\left(M_{1}\right)$ (see [13]), continuous vector functions are dense in $L_{v}^{2}\left(M_{1}\right)$ as follows from (22). Let $\vec{f}=$ $\left(f_{1}, f_{2}\right) \in L_{v}^{2}\left(M_{1}\right)$ be continuous.

For $a=-1,1$ we set

$$
\hat{f}_{k}(x)=(x-a) g_{k}(x)+f_{k}(a)
$$

where

$$
\begin{gathered}
g_{k}(x)=\frac{1}{a}\left(f_{k}(a)-f_{k}(0)\right)+x\left(\frac{1}{a^{2}}\left(f_{k}(a)-f_{k}(0)\right)-\frac{1}{a}\right), \quad x \in[0, a] \\
k=1,2 .
\end{gathered}
$$

It is not hard to see that $\hat{f}_{1}(x)$ and $\hat{f}_{2}(x)$ are continuous functions on $[-1,1]$ such that

$$
\begin{gather*}
\hat{f}_{1}(-1)=f_{1}(-1), \quad \hat{f}_{1}(1)=f_{1}(1), \hat{f}_{2}(-1)=f_{2}(-1), \quad \hat{f}_{2}(1)=f_{2}(1)  \tag{38}\\
\hat{f}_{1}(0)=f_{1}(0), \hat{f}_{2}(0)=f_{2}(0), \quad \exists \hat{f}_{1}^{\prime}(0)=1 \tag{39}
\end{gather*}
$$

Since $f_{k}(x) \rightarrow f_{k}(0)$ and $\hat{f}_{k}(x) \rightarrow f_{k}(0)$ as $x \rightarrow 0$, we obtain that

$$
\begin{gathered}
\forall \varepsilon>0 \exists \delta_{1}>0\left(\delta_{1}<1\right): x \in[-1,1],|x| \leq \delta_{1} \Rightarrow \\
\left|\hat{f}_{k}(x)-f_{k}(x)\right|<\varepsilon, k=1,2 .
\end{gathered}
$$

Set

$$
\tilde{f}_{k}(x)= \begin{cases}\hat{f}_{k}(x), & |x|<\delta_{1} \\ f_{k}(x)+\frac{|x|-1}{\delta_{1}-1}\left(\hat{f}_{k}\left(\delta_{1}\right)-f_{k}\left(\delta_{1}\right)\right), & \delta_{1} \leq|x| \leq 1 \quad, x \in \mathbb{R} ; k=1,2 \\ f_{k}(x), & |x|>1\end{cases}
$$

Then $\overrightarrow{\tilde{f}}(x):=\left(\tilde{f}_{1}(x), \tilde{f}_{2}(x)\right)$ is a continuous function on $\mathbb{R}$ such that $\tilde{f}_{1}\left(\lambda^{2}\right), \lambda \in$ $\mathbb{R} \cup i \mathbb{R}$ has the first derivative at zero and

$$
\begin{equation*}
\left|\tilde{f}_{k}(x)-f_{k}(x)\right|<\varepsilon, x \in[-1,1], k=1,2 \tag{40}
\end{equation*}
$$

Let $M_{1}^{\prime}(x)$ be the matrix of the Radon-Nikodym derivatives of the measure $M_{1}$ with respect to the trace measure $\tau_{M_{1}}$. Then

$$
\begin{aligned}
& \|\vec{f}-\overrightarrow{\tilde{f}}\|_{2, v ; M_{1}}^{2}=\int_{[-1,1]}\left(f_{1}(x)-\tilde{f}_{1}(x), f_{2}(x)-\tilde{f}_{2}(x)\right) M_{1}^{\prime}(x) * \\
& \quad *\left(f_{1}(x)-\tilde{f}_{1}(x), f_{2}(x)-\tilde{f}_{2}(x)\right)^{*} d \tau_{M_{1}} \leq 4 \varepsilon^{2} \tau_{M_{1}}([-1,1])
\end{aligned}
$$

Consequently, $R_{0}$ is dense in $L_{v}^{2}\left(M_{1}\right)$.
We can extend the operator $\mathbf{V}$ by continuity to an isometric isomorphism $\mathbf{U}$ between $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ and $L_{v}^{2}\left(M_{1}\right)$. Hence, we have obtained the following result:

Lemma 2.3. There exists an isometric isomorphism $\mathbf{U}$ between $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ and $L_{v}^{2}\left(M_{1}\right)$.

Since $L_{v}^{2}\left(M_{1}\right)$ is a complete space (see [12, Theorem, p. 295]), from Lemma 2.3 we get the following theorem.

Theorem 2.1. $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$ is a complete Hilbert space under the inner product $\langle\cdot, \cdot\rangle_{2 ; \mathbf{M}, A}$.

Special choices of $\mathbf{M}$ and $A$ in the definition of $L_{\mathbf{M}, A}^{2}$ lead to the following important $L^{2}$ spaces:
(a) $M=\left(\begin{array}{cc}m & 0 \\ 0 & 0\end{array}\right), M_{I}=0, A=0$. This leads to the usual $L_{m}^{2}(\mathbb{R})$.
(b) $M=\left(\begin{array}{cc}m & 0 \\ 0 & 0\end{array}\right), M_{I}=0, A \neq 0$. That leads to a discrete Sobolev space related to the discrete Sobolev orthogonal polynomials [7].
(c) $M=\left(\begin{array}{cc}m & 0 \\ 0 & 0\end{array}\right), M_{I}=\left(\begin{array}{cc}m_{I} & 0 \\ 0 & 0\end{array}\right): m_{I}(\{0\})=0, A=0$. Define a measure $\mathbf{m}$ on sets $[a, b] \cup[i c, i d], a, b, c, d>0$ :

$$
\mathbf{m}([a, b] \cup[i c, i d]):=m([a, b])+m_{I}([c i, d i])
$$

and in a similar manner for sets with $(a, b],[a, b),(a, b),(c i, d i],[c i, d i),(c i, d i)$. Then we define a measure $\mu$ on rectangles in $\mathbb{C}$ (with sides parallel to the axes):

$$
\mu(R)=\mathbf{m}((\mathbb{R} \cup i \mathbb{R}) \cap R), \quad \text { where } R \text { is a rectangle. }
$$

Finaly, we apply Lebesque's extention to $\mu$ to obtain a positive measure $\tilde{\mu}$ defined on the Borel subsets $\mathfrak{B}(\mathbb{C})$ in the complex plane. It is not hard to see that $L_{2 ; \mathbf{M}, A}=L_{\tilde{\mu}}^{2}$ in this case. The measure $\tilde{\mu}$ has support on $\mathbb{R} \cup i \mathbb{R}$.

Conversely, each positive measure $\tilde{\mu}$ on $\mathfrak{B}(\mathbb{C})$ with support on $\mathbb{R} \cup i \mathbb{R}$ defines measures $\tilde{M}$ and $\tilde{M}_{I}$ on intervals:

$$
\begin{gather*}
\tilde{M}([a, b]):=\lim _{\varepsilon \rightarrow+0} R(a-\varepsilon i, a+\varepsilon i, b-\varepsilon i, b+\varepsilon i)  \tag{41}\\
\tilde{M}_{I}([c i, d i]):=\lim _{\varepsilon \rightarrow+0} R(c i-\varepsilon, c i+\varepsilon, d i-\varepsilon, d i+\varepsilon) \tag{42}
\end{gather*}
$$

where $R\left(z_{1}, z_{2}, z_{3}, z_{4}\right)$ is a closed rectangle with vertices at points $z_{i}, i=1,2,3,4$; $a, b, c, d>0$.

For intervals $[a, b),(a, b],(a, b),[c i, d i),(c i, d i],(c i, d i)$ we use the obvious modifications of definitions (41),(42). Applying Lebesgue's extension we obtain $\mathbb{C}^{\geq}$-valued measures $\hat{M}$ and $\hat{M}_{I}$ on $\mathfrak{B}(\mathbb{R})$ and $\mathfrak{B}(i \mathbb{R})$, respectively. Finally, we set $M:=\hat{M}$ and for $A \in \mathfrak{B}(i \mathbb{R})$

$$
M_{I}(A):=\left\{\begin{array}{ll}
\hat{M}_{I}(A), & \text { if } 0 \notin A \\
\hat{M}_{I}(A)-\hat{M}(\{0\}), & \text { if } 0 \in A
\end{array} .\right.
$$

It is not hard to see that $L_{\tilde{\mu}}^{2}=L_{2 ; \mathbf{M}, A}$.
Let us now turn to the question of density of polynomials in $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$. It is natural to assume that complex polynomials belong to $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$. This condition is equivalent to existance of the following integrals:

$$
\begin{equation*}
\int_{\mathbb{R}}\left(\overrightarrow{\lambda^{n}}\right)_{s} M^{\prime}(\lambda)\left(\overrightarrow{\lambda^{n}}\right)_{s}^{*} d \tau_{M}+\int_{i \mathbb{R}}\left(\overrightarrow{\lambda^{n}}\right)_{s} M_{I}^{\prime}(\lambda)\left(\overrightarrow{\lambda^{n}}\right)_{s}^{*} d \tau_{M_{I}}<\infty, n \in \mathbb{Z}_{+} \tag{43}
\end{equation*}
$$

It is not hard to see that $\mathbf{V P}=\mathbb{P}_{v}$. Let polynomials be dense in $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$. Therefore they are dense in $C_{0}$. Thus, $\forall \vec{v}=\left(v_{1}, v_{2}\right) \in R_{0}$ we can find a polynomial $p(\lambda)$ such that

$$
\varepsilon>\left\|\mathbf{V}^{-1} v-p\right\|_{2 ; \mathbf{M}, A}=\|v-\mathbf{V} p\|_{2, v ; M_{1}}
$$

Hence, vector polynomials are dense in $R_{0}$ and from Lemma 2.2 it follows that they are dense in $L_{v}^{2}\left(M_{1}\right)$.

Let vector polynomials be dense in $L_{v}^{2}\left(M_{1}\right)$. Therefore they are dense in $R_{0}$. Thus, $\forall f(\lambda) \in C_{0}$ we can find a vector $\vec{p}=\left(p_{1}, p_{2}\right)$, where $p_{1}, p_{2}$ are polynomials such that

$$
\varepsilon>\|\mathbf{V} f-\vec{p}\|_{2, v ; M_{1}}=\left\|f-\mathbf{V}^{-1} \vec{p}\right\|_{2 ; \mathbf{M}, A}
$$

Hence, polynomials are dense in $C_{0}$ and from Lemma 2.2 it follows that they are dense in $L_{\mathbf{M}, A}^{2}(\mathbb{R} \cup i \mathbb{R})$.

We get that polynomials are dense in $L_{\mathbf{M}, A}^{2}$ iff vector polynomials are dense in $L_{v}^{2}\left(M_{1}\right)$.

Taking into account Corollary 2.1 we obtain the following theorem:
Theorem 2.2. Consider an $L_{\mathbf{M}, A}^{2}$ space such that for the measure $M$ condition (43) is satisfied. Then polynomials are dense in $L_{\mathbf{M}, A}^{2}$ iff matrix polynomials are dense in $L_{M_{1}}^{2}$ where $M_{1}$ is constructed as above.

We shall suppose that the measure $M$ is not degenerate in the following sense:

$$
\begin{equation*}
\langle p, p\rangle_{2 ; \mathbf{M}, A}>0, \tag{44}
\end{equation*}
$$

for any complex polynomial $p(\lambda) \neq 0$.
Applying the Gramm-Schmidt orthogonalization method to the sequence $1, \lambda, \lambda^{2}, \ldots$, with respect to $\langle\cdot, \cdot\rangle_{2 ; \mathbf{M}, A}$ we get a system of orthonormal polynomials $\left\{p_{n}\right\}_{n=0}^{\infty}$. They satisfy the orthonormality property

$$
\begin{gather*}
\int_{\mathbb{R}}\left(\overrightarrow{p_{n}}\right)_{s}(\lambda) M^{\prime}(\lambda)\left(\overrightarrow{p_{m}}\right)_{s}^{*}(\lambda) d \tau_{M}+\int_{i \mathbb{R}}\left(\overrightarrow{p_{n}}\right)_{s}(\lambda) M_{I}^{\prime}(\lambda)\left(\overrightarrow{p_{m}}\right)_{s}^{*}(\lambda) d \tau_{M_{I}}+ \\
+\left(\overrightarrow{p_{n}}\right)_{d} A\left(\overrightarrow{p_{m}}\right)_{d}^{*}=\delta_{n, m}, \quad n, m \in \mathbb{Z}_{+} \tag{45}
\end{gather*}
$$

Special choices of $\mathbf{M}$ and $A$ in the definition of $L_{\mathbf{M}, A}^{2}$, see cases (a),(b),(c) after Theorem 2.1, lead to the following important systems of orthogonal polynomials:
(a): Orthogonal polynomials on the real line, see [10];
(b): Discrete Sobolev orthogonal polynomials, see [7];
(c): Orthogonal polynomials on radial rays in the complex plane, see [11].

For orthonormal polynomials from (45) we define

$$
\vec{p}_{n}(x):=\left(\mathbf{V} p_{n}\right)(x)=\left(p_{n}^{+}(\sqrt{x}), p_{n}^{-}(\sqrt{x})\right), x \in \mathbb{R}, n \in \mathbb{Z}_{+} .
$$

Then $\left\{\vec{p}_{n}(x)\right\}_{n \in \mathbb{Z}_{+}}$will be an orthonormal system in $L_{v}^{2}\left(M_{1}\right)$. Define matrix polynomials

$$
P_{k}(x):=\left(\begin{array}{cc}
p_{2 k}^{+}(\sqrt{x}) & p_{2 k}^{-}(\sqrt{x})  \tag{46}\\
p_{2 k+1}^{+}(\sqrt{x}) & p_{2 k+1}^{-}(\sqrt{x})
\end{array}\right)=\binom{\vec{p}_{2 k}(x)}{\vec{p}_{2 k+1}(x)}, k \in \mathbb{Z}_{+}
$$

Notice that $P_{k}(x)$ has degree exactly $k$. Moreover, we have

$$
\begin{equation*}
\int_{\mathbb{R}} P_{k}(x) d M_{1}(x) P_{l}^{*}(x)=I \delta_{k, l}, k, l \in \mathbb{Z}_{+} \tag{47}
\end{equation*}
$$

where $I=\left(\delta_{i, j}\right)_{i, j=1}^{2}$.
Hence, $\left\{P_{k}(x)\right\}_{k=0}^{\infty}$ are matrix orthonormal polynomials. Denote by $J$ the corresponding Jacobi matrix and let $\mathbf{J}$ be the corresponding symmetric operator in $l^{2}$. Suppose that $\mathbf{J}$ has indices of deficiency $(2,2)$. In other words that means that the matrix Hamburger moment problem corresponding to $M_{1}$ is completely indeterminate.

In this case we get the following theorem:
Theorem 2.3. Consider an $L_{\mathbf{M}, A}^{2}$ space such that for the measure $M$ conditions (43), (44) are satisfied. Construct the measure $M_{1}$ as above. Suppose that the corresponding symmetric operator $\mathbf{J}$ has indices of deficiency (2, 2). Then polynomials are dense in $L_{\mathbf{M}, A}^{2}$ iff condition (5) holds true for $M=M_{1}$ with a constant unitary matrix $U$.

Note that there are several criteria for the completely indeterminate case of the matrix moment problem. For example, one can use a generalization of Hamburger's criterion [14, p. 56], see also [15]-[17] and references therein.

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