# THE MAXIMAL $\mathscr{J}$-REGULAR PART OF A $q$-VARIATE WEAKLY STATIONARY PROCESS 

By

LUTZ KLOTZ (LEIPZIG)


#### Abstract

Let $x$ be a $q$-variate (weakly) stationary process over a locally compact Abelian group $\boldsymbol{G}$, and $\mathscr{J}$ a family of subsets of $\boldsymbol{G}$ invariant under translation. We show that the set of all regular non-negative Hermitian matrix-valued measures $M$ not exceeding the (non-stochastic) spectral measure of $x$ and such that the Hilbert space $L^{2}(M)$ is $\mathscr{J}$-regular contains a unique maximal element. Moreover, this maximal element coincides with the spectral measure of the $\mathscr{J}$-regular part of the Wold decomposition of $x$.


1991 Mathematics Subject Classification: Primary 60G25; Secondary 15A57.

## 1. INTRODUCTION

Let $N$ be the set of positive integers and $q \in N$. By $M_{q}$ we denote the algebra of $q \times q$-matrices with entries from the field of complex numbers $C$ and by $M_{q}^{\geqslant}$the subset of non-negative Hermitian matrices. The symbol $I$ stands for the unit matrix of $M_{q}$.

Let $\boldsymbol{G}$ be a locally compact Abelian group, $\Gamma$ its dual, and $\langle g, \gamma\rangle$ the value of a character $\gamma \in \boldsymbol{\Gamma}$ on $g \in \boldsymbol{G}$. If $J$ is a subset of $\boldsymbol{G}$, then a (finite) $\boldsymbol{M}_{\boldsymbol{q}}$-linear combination of functions $\langle g, \cdot\rangle I, g \in J$, is called a trigonometric polynomial with frequencies from $J$.

Let $x$ be a $q$-variate (weakly) stationary process over $\boldsymbol{G}$, and $\boldsymbol{H}_{x}$ its time domain, i.e. the left Hilbert- $\boldsymbol{M}_{\boldsymbol{q}}$-module spanned by the values of $x$. If $\mathscr{J}$ is a family of subsets of $\boldsymbol{G}$ invariant under translation, then there exists a unique Wold decomposition of $x$ into an orthogonal sum of $q$-variate stationary processes $y$ and $z$ such that $y$ is $\mathscr{J}$-regular and $z$ is $\mathscr{J}$-singular (cf. [12], Theorem 2.13). It could be expected that, in a certain sense, the process $y$ is the "maximal $\mathscr{J}$-regular part of $x$ ". The aim of this note is to specify this statement. To do this it is more convenient to work with the spectral domain instead of the time domain of $x$.

Let $\mathscr{B}(\Gamma)$ be the $\sigma$-algebra of Borel sets of $\Gamma$. The (non-stochastic) spectral measure $M_{x}$ of $x$ (cf. [12], Definition 3.5) is a regular $M_{q}^{\geqslant}$-valued measure on
$\mathscr{B}(\Gamma)$. Loewner's partial ordering of $M_{q}^{\geqslant}$induces a partial ordering on the set of all regular $M_{q}^{\geqslant}$-valued measures on $\mathscr{B}(\Gamma)$. We will show (see Theorem 3.3) that among all regular $M_{q}^{\geqslant}$-valued measures $M$ on $\mathscr{B}(\Gamma)$, which do not exceed $M_{x}$ and for which the space $L^{2}(M)$ is $\mathscr{J}$-regular, there exists a maximal measure. Moreover, in Section 4 it will be shown that this maximal measure coincides with the spectral measure of the $\mathscr{J}$-regular part $y$ of the Wold decomposition of $x$. Section 5 deals with an application of our results to the case where $\mathscr{J}$ is the family $\mathscr{J}_{0}$ of complements of all singletons of $G$. Using Makagon and Weron's characterization of $\mathscr{J}_{0}$-regular processes (see [7], Theorem 5.3), we compute the spectral measures of the $\mathscr{J}_{0}$-regular and $\mathscr{F}_{0}$-singular parts of the Wold decomposition of $x$.

## 2. PRELIMINARIES

For any matrix $B$ with complex entries, denote by $B^{*}$ its adjoint and by $\mathscr{R}(B)$ its range. For $A \in M_{q}$, let $\operatorname{ker} A, \operatorname{tr} A$, and $A^{+}$be the kernel, trace, and Moore-Penrose inverse of $A$, respectively. Let $P_{A}$ be the orthoprojector in the left Hilbert- $\boldsymbol{M}_{q}$-module $\boldsymbol{C}^{q}$ of column vectors of length $q$ onto $\mathscr{R}(A)$. If $A \in \boldsymbol{M}_{q}^{\geqslant}$, we denote by $A^{1 / 2}$ the unique non-negative Hermitian square root of $A$. We equip $M_{q}^{\geqslant}$with Loewner's partial ordering, i.e. we write $A \leqslant B$ if and only if $B-A$ is a non-negative Hermitian, $A, B \in M_{q}^{\geqslant}$.

We give some more or less known results on $M_{q}^{\geqslant}$and the measurability of $M_{q}$-valued functions, which for ease of reference will be stated as lemmas.

Lemma 2.1. Let $\mathscr{D}$ be a directed subset of $\mathbb{M}_{q}^{\geqslant}$, which has an upper bound. Then there exists a least upper bound $C$ of $\mathscr{D}$ and we have

$$
u^{*} C u=\sup \left\{u^{*} D u: D \in \mathscr{D}\right\}, \quad u \in C^{q}
$$

Proof. For $u \in C^{q}$, set $t(u):=\sup \left\{u^{*} D u: D \in \mathscr{D}\right\}$. Obviously, if $\lambda \in C$, we have

$$
\begin{equation*}
t(\lambda u)=|\lambda|^{2} t(u) \tag{2.1}
\end{equation*}
$$

and if $u, v \in C^{q}$, we obtain

$$
\begin{equation*}
\sup \left\{u^{*} D u+v^{*} D v: D \in \mathscr{D}\right\} \leqslant t(u)+t(v) \tag{2.2}
\end{equation*}
$$

Since $\mathscr{D}$ is directed, for $D_{1}, D_{2} \in \mathscr{D}$ there exists $D_{3} \in \mathscr{D}$ such that

$$
u^{*} D_{1} u+v^{*} D_{2} v \leqslant u^{*} D_{3} u+v^{*} D_{3} v .
$$

This yields

$$
\begin{equation*}
t(u)+t(v) \leqslant \sup \left\{u^{*} D u+v^{*} D v: D \in \mathscr{D}\right\} . \tag{2.3}
\end{equation*}
$$

The parallelogram identity implies that

$$
\begin{align*}
\sup \left\{(u+v)^{*} D(u+v)+(u-v)^{*} D(u-v)\right. & : D \in \mathscr{D}\}  \tag{2.4}\\
& =2 \sup \left\{u^{*} D u+v^{*} D v: D \in \mathscr{D}\right\} .
\end{align*}
$$

Combining (2.4), (2.2), and (2.3), we get

$$
\begin{equation*}
t(u+v)+t(u-v)=2 t(u)+2 t(v) \tag{2.5}
\end{equation*}
$$

From (2.1) and (2.5) it follows that there exists $C \in M_{q}^{\geqslant}$such that $t(u)=u^{*} C u$, $u \in \boldsymbol{C}^{q}$. From the definition of $t$ it is clear that $C$ is the least upper bound of $\mathscr{D}$.

Lemma 2.2 (cf. [1], Theorem 1). Let $p, q \in N$. A block matrix

$$
X=\left(\begin{array}{ll}
X_{11} & X_{12} \\
X_{12}^{*} & X_{22}
\end{array}\right) \in M_{p+q}
$$

belongs to $M_{p+q}^{\geqslant}$if and only if
(i) $\mathscr{R}\left(X_{12}^{*}\right) \subseteq \mathscr{R}\left(X_{22}\right)$,
(ii) $X_{22} \in M_{q}^{\geq}$,
(iii) $X_{11}-X_{12} X_{22}^{+} X_{12}^{*}=:\left(X / X_{22}\right)$ is a non-negative Hermitian.

Lemma 2.3 (cf. [3], p. 391). If $F$ is a (Borel) measurable $M_{q}$-valued function on $\Gamma$, then $P_{F}$ is measurable. If $W$ is a measurable $M_{q}^{\geqslant}$-valued function on $\Gamma$, then $W^{1 / 2}$ and $W^{+}$are measurable.

Let $M$ be a regular $M_{q}^{\geq}$-valued measure on $\mathscr{B}(\Gamma)$ and $\tau$ a regular non--negative $\sigma$-finite measure on $\mathscr{B}(\Gamma)$ such that $M$ is absolutely continuous with respect to $\tau$. For example, one can take $\tau=\operatorname{tr} M$. Let $W:=d M / d \tau$ be the Radon-Nikodym derivative of $M$ with respect to (abbreviated to "w.r.t.") $\tau$. By definition, the left Hilbert- $M_{q}$-module $L^{2}(M)$ consists of (equivalence classes of) measurable $M_{q}$-valued functions $F$ on $\Gamma$ such that

$$
\operatorname{tr} \int_{\boldsymbol{\Gamma}} F(\gamma) W(\gamma) F(\gamma)^{*} \tau(d \gamma)<\infty
$$

The corresponding scalar product of $L^{2}(M)$ is defined by

$$
\operatorname{tr} \int_{\boldsymbol{F}} F(\gamma) W(\gamma) G(\gamma)^{*} \tau(d \gamma), \quad F, G \in L^{2}(M) .
$$

The definition does not depend on the choice of $\tau$ (cf. [10]).
Lemma 2.4. Let $F \in L^{2}(M)$. Then $F=0$ in $L^{2}(M)$ if and only if $\mathscr{R}(W) \subseteq \operatorname{ker} F \quad \tau$-a.e.

Proof. Since $F W F^{*}=F W^{1 / 2}\left(F W^{1 / 2}\right)^{*}$, we have $F=0$ in $L^{2}(M)$ if and only if $\mathscr{R}\left(W^{1 / 2}\right) \subseteq \operatorname{ker} F \tau$-a.e. Since $\mathscr{R}\left(W^{1 / 2}\right)=\mathscr{R}(W)$, the result follows.

If $M_{x}$ is the spectral measure of a $q$-variate stationary process $x$ over $G$, the corresponding space $L^{2}\left(M_{x}\right)$ is called the spectral domain of $x$. There exists an isometric and isomorphic map $V_{x}$ of $H_{x}$ onto $L^{2}\left(M_{x}\right)$ such that $V_{x} x_{g}=\langle g, \cdot\rangle I$,
$g \in \boldsymbol{G}$. The map $V_{x}$ is called Kolmogorov's isomorphism. It enables us to formulate $\mathscr{J}$-regularity and $\mathscr{\mathscr { L }}$-singularity of $x$ in terms of $L^{2}\left(M_{x}\right)$. According to this we call a space $L^{2}(M) \mathscr{J}$-regular or $\mathscr{J}$-singular if and only if

$$
\bigcap_{J \in g} \bigvee_{M}\{\langle g, \cdot\rangle I: g \in J\}=\{0\} \quad \text { or } \quad \bigvee_{M}\{\langle g, \cdot\rangle I: g \in J\}=L^{2}(M)
$$

for all $J \in \mathscr{J}$, respectively. The symbol $\bigvee_{M}$ stands for the closed $M_{q}$-linear hull in $L^{2}(M)$. We simply write $\bigvee$ if $M=M_{x}$ is the spectral measure of the process $x$.

## 3. THE MAXIMAL $\mathscr{J}$-REGULAR PART

Let $M_{x}$ be the spectral measure of a $q$-variate stationary process over $\boldsymbol{G}$, $\tau_{x}:=\operatorname{tr} M_{x}$, and $W_{x}:=d M_{x} / d \tau_{x}$. In the sequel, all relations between measurable functions on $\Gamma$ are to be understood as relations which hold true $\tau_{x}$-a.e.

Let $\mathscr{W}_{x}$ be the set of all measurable $M_{q}^{\geqslant}$-valued functions $W$ on $\Gamma$ such that $W \leqslant W_{x}$ and let $\widetilde{\mathscr{W}}_{x}$ be the set of all $M_{q}^{\geqslant}$-valued measures of the form $W d \tau_{x}$, $W \in \mathscr{W}_{x}$. The partial ordering on $\mathscr{W}_{x}$ induces a partial ordering on $\tilde{W}_{x}$ : define $W_{1} d \tau_{x} \leqslant W_{2} d \tau_{x}$ if and only if $W_{1} \leqslant W_{2}, W_{1}, W_{2} \in \mathscr{W}_{x}$. Note that for $M_{1}, M_{2} \in \tilde{W}_{x}$ we have $M_{1} \leqslant M_{2}$ if and only if $M_{1}(\Delta) \leqslant M_{2}(\Delta), \Delta \in \mathscr{B}(\Gamma)$.

Lemma 3.1. For any directed subset $\mathscr{D}$ of $\mathscr{W}_{x}$, there exists a least upper bound.

Proof. According to the remarks preceding the lemma it is enough to show that the subset $\tilde{\mathscr{D}}:=\left\{W d \tau_{x}: W \in \mathscr{D}\right\}$ of $\tilde{\mathscr{W}}_{x}$ has the least upper bound. For $\Delta \in \mathscr{B}(\Gamma)$, let $\mathscr{D}_{\Delta}$ be the set of matrices of the form

$$
\begin{equation*}
\sum_{j=1}^{n} M_{j}\left(\Delta_{j}\right) \tag{3.1}
\end{equation*}
$$

where $M_{1}, \ldots, M_{n} \in \tilde{\mathscr{D}}$, and $\left\{\Delta_{1}, \ldots, \Delta_{n}\right\}$ is a partition of $\Delta, n \in N$. The matrix $M_{x}(\Delta)$ is an upper bound of $\mathscr{D}_{\Delta}$. Moreover, $\mathscr{D}_{\Delta}$ is a directed set. In fact, if (3.1) and

$$
\begin{equation*}
\sum_{k=1}^{m} M_{k}^{\prime}\left(\Delta_{k}^{\prime}\right) \tag{3.2}
\end{equation*}
$$

are two elements of $\mathscr{D}_{\Delta}$, consider $M_{j k} \in \tilde{\mathscr{D}}$ such that $M_{j} \leqslant M_{j k}, M_{k}^{\prime} \leqslant M_{j k}$, $j=1, \ldots, n, k=1, \ldots, m$. Then $\sum_{j=1}^{n} \sum_{k=1}^{m} M_{j k}\left(\Delta_{j} \cap \Delta_{k}^{\prime}\right)$ belongs to $\mathscr{D}_{\Delta}$ and exceeds both matrices (3.1) and (3.2). From Lemma 2.1 it follows that $\mathscr{D}_{\Delta}$ has the least upper bound $N(\Delta)$ and that

$$
\begin{equation*}
u^{*} N(\Delta) u=\sup \left\{u^{*} D u: D \in \mathscr{D}_{\Delta}\right\}, \quad u \in C^{q} . \tag{3.3}
\end{equation*}
$$

Standard measure-theoretic arguments (cf. the proof of Theorem 5 of Section III. 7 of [2]) show that, for $u \in C^{q}, u^{*} N u$ is an additive function on $\mathscr{B}(\Gamma)$. Hence
$N$ is additive. Since $N \leqslant M_{x}$, it even belongs to $\tilde{\mathscr{W}}_{x}$. Finally, from (3.3) it follows easily that $N$ is the least upper bound of $\mathscr{\mathscr { D }}$.

If $W \in \mathscr{W}_{x}$, set $L^{2}(W):=L^{2}\left(W d \tau_{x}\right)$. Moreover, we define

$$
\mathscr{W}_{x}^{(r)}:=\left\{W \in \mathscr{W}_{x}: L^{2}(W) \text { is } \mathscr{J} \text {-regular }\right\}
$$

Lemma 3.2. The set $\mathscr{W}_{x}^{(r)}$ is directed.
Proof. Let $W_{1}, W_{2} \in \mathscr{W}_{x}^{(r)}$ and let $Q(\gamma)$ be the orthogonal projection in $C^{q}$ onto the algebraic sum $\mathscr{R}\left(W_{1}(\gamma)\right)+\mathscr{R}\left(W_{2}(\gamma)\right), \gamma \in \Gamma$. From von Neumann's alternating projections theorem (cf. [4], Problem 96) we can conclude the measurability of the function $Q$. Let

$$
W_{x}=\left(\begin{array}{ll}
W_{x, 11} & W_{x, 12} \\
W_{x, 12}^{*} & W_{x, 22}
\end{array}\right)
$$

be the block partition of $W_{x}$ w.r.t. the orthogonal decomposition

$$
C^{q}=Q C^{q} \oplus(I-Q) C^{q}
$$

Let us set

$$
W_{3}:=\left(\begin{array}{cc}
W_{x, 11}-W_{x, 12} W_{x, 22}^{+} W_{x, 12}^{*} & 0 \\
0 & 0
\end{array}\right)
$$

The measurability of $Q$ and Lemmas 2.3 and 2.2 imply that $W_{3} \in \mathscr{W}_{x}$. Moreover, from Lemma 2.2 it follows that $W_{1} \leqslant W_{3}$ and $W_{2} \leqslant W_{3}$. To complete the proof it is enough to show that $L^{2}\left(W_{3}\right)$ is $\mathscr{\mathscr { I }}$-regular. Let $F \in L^{2}\left(W_{3}\right)$ be such that for each $J \in \mathscr{J}$ it can be approximated by trigonometric polynomials with frequencies from $J$ in $L^{2}\left(W_{3}\right)$. Since $W_{1} \leqslant W_{3}$, an analogous approximation exists in $L^{2}\left(W_{1}\right)$. The $\mathscr{J}$-regularity of $L^{2}\left(W_{1}\right)$ yields $F=0$ in $L^{2}\left(W_{1}\right)$. Similarly, $F=0$ in $L^{2}\left(W_{2}\right)$. Using Lemma 2.4 , we can conclude that $\mathscr{R}\left(W_{1}\right)+\mathscr{R}\left(W_{2}\right) \subseteq \operatorname{ker} F$. Since $\mathscr{R}\left(W_{3}\right) \subseteq \mathscr{R}\left(W_{1}\right)+\mathscr{R}\left(W_{2}\right)$, it follows that $F=0$ in $L^{2}\left(W_{3}\right)$.

Theorem 3.3. The set $\mathscr{W}_{x}^{(r)}$ has a unique maximal element.
Proof. By Lemmas 3.1 and 3.2, the set $\mathscr{W}_{x}^{(r)}$ has the least upper bound $W^{(r)} \in \mathscr{W}_{x}$. Assume that $L^{2}\left(W^{(r)}\right)$ is not $\mathscr{J}$-regular. Then there exists $F \in L^{2}\left(W^{(r)}\right)$, $F \neq 0$, such that, for each $J \in \mathscr{J}, F$ can be approximated by trigonometric polynomials with frequencies from $J$. Let $W \in \mathscr{W}_{x}^{(r)}$. Then, in particular, $W \leqslant W^{(r)}$, and similar arguments to those in the proof of Lemma 3.2 show that

$$
\begin{equation*}
\mathscr{R}(W) \subseteq \operatorname{ker} F \tag{3.4}
\end{equation*}
$$

Let

$$
W^{(r)}=\left(\begin{array}{ll}
W_{11}^{(r)} & W_{12}^{(r)} \\
W_{12}^{(r) *} & W_{22}^{(r)}
\end{array}\right)
$$

be the block partition of $W^{(r)}$ w.r.t. the orthogonal decomposition $C^{q}=$ $=\mathscr{R}\left(F^{*}\right) \oplus \operatorname{ker} F$. Let us set

$$
W^{(\varrho)}:=\left(\begin{array}{cc}
W_{11}^{(r)}-W_{12}^{(r)} W_{22}^{(r)+} W_{12}^{(r) *} & 0 \\
0 & 0
\end{array}\right)
$$

It is not hard to see (cf. the proof of Lemma 3.2) that

$$
\begin{equation*}
W^{(\varrho)} \in \mathscr{W}_{x}, \quad W^{(0)} \leqslant W^{(r)}, \quad \text { and } \quad W \leqslant W^{(\varrho)} \tag{3.5}
\end{equation*}
$$

On the other hand, since $F \neq 0$ in $L^{2}\left(W^{(r)}\right)$, Lemma 2.4 implies that there exists $\Delta \in \mathscr{B}(\Gamma)$ such that $\tau_{x}(\Delta)>0$ and $\mathscr{R}\left(W^{(r)}\right)$ is not a subspace of $\operatorname{ker} F$ on $\Delta$. It follows that $W_{22}^{(r)} \neq 0$ on $\Delta$, and hence $W^{(r)}-W^{(e)} \neq 0$ on $\Delta$. Combining this with (3.5), we obtain a contradiction to the definition of a least upper bound. Thus, $L^{2}\left(W^{(r)}\right)$ is $\mathscr{J}$-regular and $W^{(r)}$ is a maximal element of $\mathscr{W}_{x}^{(r)}$. Its uniqueness follows from Lemma 3.2.

## 4. CONCORDANCE OF THE MAXIMAL REGULAR PART and the regular part of the wold decomposition

Let $x$ be a $q$-variate stationary process over $\boldsymbol{G}$ and $\mathscr{J}$ a family of subsets of $\boldsymbol{G}$ invariant under translation. Let $x_{g}=y_{g}+z_{g}, g \in \boldsymbol{G}$, be the Wold decomposition of $x$, where $y$ is $\mathscr{J}$-regular and $z$ is $\mathscr{J}$-singular. If we set $\widetilde{\mathscr{W}}_{y}:=d M_{y} / d \tau_{x}$ and $\tilde{W}_{z}:=d M_{z} / d \tau_{x}$, we have (cf. [9], Lemmas 4.3 and 4.4)

$$
\begin{equation*}
\tilde{W}_{y}+\tilde{W}_{z}=W_{x}, \quad \mathscr{R}\left(\tilde{W}_{y}\right) \cap \mathscr{R}\left(\tilde{W}_{z}\right)=\{0\} . \tag{4.1}
\end{equation*}
$$

Let $V_{x}$ be Kolmogorov's isomorphism of $\boldsymbol{H}_{x}$ onto $L^{2}\left(W_{x}\right)$ and set

$$
V_{x} y_{0}=: F_{y} \quad \text { and } \quad V_{x} z_{0}=: F_{z}
$$

where 0 is the neutral element of $\boldsymbol{G}$. It is not hard to see that

$$
\begin{equation*}
F_{y} W_{x} F_{y}^{*}=\tilde{W}_{y} \quad \text { and } \quad F_{z} W_{x} F_{z}^{*}=\tilde{W}_{z} \tag{4.2}
\end{equation*}
$$

$$
\begin{equation*}
V_{x} \boldsymbol{H}_{y}=\bigvee\left\{\langle g, \cdot\rangle F_{y}: g \in \boldsymbol{G}\right\} \quad \text { and } \quad V_{x} \boldsymbol{H}_{z}=\bigvee\left\{\langle g, \cdot\rangle F_{z}: g \in \boldsymbol{G}\right\}, \tag{4.3}
\end{equation*}
$$ and hence

$$
\begin{equation*}
\bigvee\left\{\langle g, \cdot\rangle F_{y}: g \in \boldsymbol{G}\right\} \oplus \bigvee\left\{\langle g, \cdot\rangle F_{z}: g \in \boldsymbol{G}\right\}=L^{2}\left(W_{x}\right) . \tag{4.4}
\end{equation*}
$$

From the relation (4.4) it follows that

$$
\int_{\boldsymbol{\Gamma}}\langle g, \gamma\rangle F_{y}(\gamma) W_{x}(\gamma) F_{z}(\gamma)^{*} \tau_{x}(d \gamma)=0, \quad g \in \boldsymbol{G},
$$

which yields

$$
\begin{equation*}
F_{y} W_{x} F_{z}^{*}=0 \tag{4.5}
\end{equation*}
$$

We can assume and we will do so in the sequel that

$$
\begin{equation*}
\operatorname{ker} W_{x} \subseteq\left(\operatorname{ker} F_{y} \cap \operatorname{ker} F_{z}\right) \tag{4.6}
\end{equation*}
$$

Then we have

$$
\begin{equation*}
P_{W_{x}}=F_{y}+F_{z} \tag{4.7}
\end{equation*}
$$

as well as
(4.8) $\mathscr{R}\left(F_{y} W_{x}\right)=\mathscr{R}\left(F_{y} W_{x}^{1 / 2}\right)=\mathscr{R}\left(F_{y}\right)$ and $\mathscr{R}\left(F_{z} W_{x}\right)=\mathscr{R}\left(F_{z} W_{x}^{1 / 2}\right)=\mathscr{R}\left(F_{z}\right)$.

Moreover, from (4.2) it follows that $\mathscr{R}\left(\tilde{W}_{y}^{1 / 2}\right)=\mathscr{R}\left(F_{y} W_{x}^{1 / 2}\right)$ and $\mathscr{R}\left(\tilde{W}_{z}^{1 / 2}\right)=$ $=\mathscr{R}\left(F_{z} W_{x}^{1 / 2}\right)$. Combining this with (4.8), we obtain

$$
\begin{equation*}
\mathscr{R}\left(\tilde{W}_{y}\right)=\mathscr{R}\left(F_{y}\right) \quad \text { and } \quad \mathscr{R}\left(\tilde{W}_{z}\right)=\mathscr{R}\left(F_{z}\right) \tag{4.9}
\end{equation*}
$$

Let $W^{(r)}$ be the maximal element of $\mathscr{W}_{x}^{(r)}$. We wish to show that $W^{(r)}$ coincides with $\tilde{W}_{y}$. In order to prove this we first derive some properties of $F_{z}$, which eventually lead to the conclusion that $P_{F_{z}^{*}}=0$ in $L^{2}\left(W^{(r)}\right)$. Then we will see that the assumption $W^{(r)} \neq \tilde{W}_{y}$ would imply that $P_{F_{z}^{*}} \neq 0$ in $L^{2}\left(W^{(r)}\right)$.

Lemma 4.1. The values of $F_{z}$ are diagonalizable matrices.
Proof. From (4.5) and (4.7) it follows that

$$
\begin{equation*}
F_{z} W_{x} F_{z}^{*}=W_{x} F_{z}^{*} \tag{4.10}
\end{equation*}
$$

Since $\mathscr{R}\left(F_{z} W_{x} F_{z}^{*}\right)=\mathscr{R}\left(F_{z} W_{x}^{1 / 2}\right)$, from (4.8) and (4.10) we obtain

$$
\begin{equation*}
\mathscr{R}\left(F_{z}\right)=\mathscr{R}\left(W_{x} F_{z}^{*}\right) \subseteq \mathscr{R}\left(W_{x}\right) . \tag{4.11}
\end{equation*}
$$

On the other hand, (4.6) gives

$$
\begin{equation*}
\mathscr{R}\left(F_{z}^{*}\right) \subseteq \mathscr{R}\left(W_{x}\right) \tag{4.12}
\end{equation*}
$$

The relations (4.11) and (4.12) show that it is enough to prove that the restrictions $\bar{F}_{z}$ of $F_{z}$ to $\mathscr{R}\left(W_{x}\right)$ are diagonalizable. Denoting by $\bar{W}_{x}$ the restrictions of $W_{x}$ to $\mathscr{R}\left(W_{x}\right)$, from (4.10)-(4.12) we get $\bar{F}_{z} \bar{W}_{x} \bar{F}_{z}^{*}=\bar{W}_{x} \bar{F}_{z}^{*}$, which yields

$$
\bar{W}_{x}^{-1 / 2} \bar{F}_{z} \bar{W}_{x} \bar{F}_{z}^{*} \bar{W}_{x}^{-1 / 2}=\bar{W}_{x}^{1 / 2} \bar{F}_{z}^{*} \bar{W}_{x}^{-1 / 2}
$$

This shows that the values of $\bar{F}_{z}^{*}$, and hence of $\bar{F}_{z}$, are similar to self-adjoint matrices, which implies that they are diagonalizable.

Lemma 4.2. We have $\operatorname{ker} P_{F_{z}^{*}} \cap \mathscr{R}\left(F_{z}\right)=\{0\}$.
Proof. Let $\gamma \in \Gamma$ and $u \in\left(\operatorname{ker} P_{F_{z}(\gamma)^{*}}\right) \cap \mathscr{R}\left(F_{z}(\gamma)\right)$. Then $u \in \operatorname{ker} F_{z}(\gamma)$, $u=F_{z}(\gamma) v$ for some $v \in C^{q}$, and hence $F_{z}(\gamma)^{2} v=0$. If $u \neq 0$ were true, this would contradict Lemma 4.1.

Lemma 4.3. We have $P_{F_{z}^{*}}=0$ in $L^{2}\left(W^{(r)}\right)$.

Proof. From (4.5) we get $F_{y} W_{x} P_{F_{z}^{*}}=0$. This implies that the function $P_{F_{z}^{*}}$ is orthogonal (in $L^{2}\left(W_{x}\right)$ ) to $\bigvee\left\{\langle g, \cdot\rangle F_{y}: g \in \boldsymbol{G}\right\}$. Examining the proof of the Wold decomposition (cf. the proof of Theorem 2.13 of [12]) and taking into account Kolmogorov's isomorphism, we obtain

$$
\bigvee\left\{\langle g, \cdot\rangle F_{z}: g \in \boldsymbol{G}\right\}=\bigcap_{J \in \mathscr{g}} \bigvee\{\langle g, \cdot\rangle I: g \in J\}
$$

It follows that, for $J \in \mathscr{J}, P_{F_{z}^{*}}$ can be approximated by trigonometric polynomials with frequencies from $J$ in $L^{2}\left(W_{x}\right)$. Since $W^{(r)} \leqslant W_{x}$, an analogous result is true for $L^{2}\left(W^{(r)}\right)$. But since $L^{2}\left(W^{(r)}\right)$ is $\mathscr{J}$-regular, we conclude that $P_{F_{z}^{*}}=0$ in $L^{2}\left(W^{(r)}\right)$.

Theorem 4.4. The functions $W^{(r)}$ and $\tilde{W}_{y}$ coincide.
Proof. Since $\tilde{W}_{y} \in \mathscr{W}_{x}^{(r)}$, it follows that $\tilde{W}_{y} \leqslant W^{(r)}$. Assume that $\tilde{W}_{y} \neq W^{(r)}$ on a set $\Delta \in \mathscr{B}(\Gamma)$ such that $\tau_{x}(\Delta)>0$. First note that $\mathscr{R}\left(\tilde{W}_{y}\right) \neq \mathscr{R}\left(W^{(r)}\right)$ on $\Delta$. For if $\mathscr{R}\left(\tilde{W}_{y}\right)=\mathscr{R}\left(W^{(r)}\right)$ and $\tilde{W}_{y} \neq W^{(r)}$ were true on a set of positive measure $\tau_{x}$, we would get $\mathscr{R}\left(W^{(r)}-\tilde{W}_{y}\right) \cap \mathscr{R}\left(\tilde{W}_{y}\right) \neq\{0\}$, and because of $\mathscr{R}\left(\widetilde{W}_{z}\right)=$ $=\mathscr{R}\left(W_{x}-\tilde{W}_{y}\right) \supseteq \mathscr{R}\left(W^{(r)}-\tilde{W}_{y}\right)$ also $\mathscr{R}\left(\tilde{W}_{z}\right) \cap \mathscr{R}\left(\tilde{W}_{y}\right) \neq\{0\}$, which contradicts (4.1). Thus, $\mathscr{R}\left(\tilde{W}_{y}\right)$ is a proper subspace of $\mathscr{R}\left(W^{(r)}\right)$ on $\Delta$. Then from (4.1) and (4.9) it follows that $\mathscr{R}\left(W^{(r)}\right) \cap \mathscr{R}\left(F_{z}\right) \neq\{0\}$ on $\Delta$. Combining this with Lemma 4.2, we infer that $\mathscr{R}\left(W^{(r)}\right)$ is not a subspace of $\operatorname{ker} P_{F_{z}^{*}}$ on $\Delta$. Applying Lemma 2.4, we conclude that $P_{F_{z}^{*}} \neq 0$ in $L^{2}\left(W^{(r)}\right)$, which is a contradiction to Lemma 4.3.

Let us mention the following consequence of Theorem 4.4.
Corollary 4.5. If $L^{2}\left(W_{x}\right)$ is $\mathscr{J}$-singular, then for $W \in \mathscr{W}_{x}$ so is $L^{2}(W)$.
Proof. The $\mathscr{J}$-singularity of $L^{2}\left(W_{x}\right)$ and Theorem 4.4 imply that $\mathscr{W}_{x}^{(r)}=\{0\}$. For $W \in W_{x}$, consider the Wold decomposition of the corresponding stationary process over $\boldsymbol{G}$. Since the spectral measure of its $\mathscr{J}$-regular part belongs to $\mathscr{W}_{x}^{(r)}$, it is zero measure. Thus, $L^{2}(W)$ is $\mathscr{J}$-singular.

Remark 4.6. It would be of interest to have generalizations of Theorem 4.4 to the infinite-variate case. Treil' ([13], Theorem 3.1) gave such a result if $\boldsymbol{G}$ is the group of integers and $\mathscr{J}$ is the family of translates of the set of non-negative integers.

## 5. THE MAXIMAL $\mathscr{J}_{0}$-REGULAR PART

Let $\boldsymbol{G}$ be a discrete Abelian group, $\mathscr{J}_{0}$ the family of complements of all singletons of $\boldsymbol{G}$, and $\sigma$ the normalized Haar measure of $\Gamma$. Let $M_{x}$ be the spectral measure of a $q$-variate stationary process over $\boldsymbol{G}$.

Theorem 5.1 ([7], Theorem 5.3). The space $L^{2}\left(M_{x}\right)$ is $\mathscr{J}_{0}$-regular if and only if
(i) $M_{x}$ is absolutely continuous w.r.t. $\sigma$,
(ii) $\mathscr{R}\left(d M_{x} / d \sigma\right)=$ const $\sigma$-a.e.,
(iii) $\left(d M_{x} / d \sigma\right)^{+}$is integrable w.r.t. $\sigma$.

It follows that the maximal $\mathscr{F}_{0}$-regular parts of $M_{x}$ and of the absolutely continuous part of $M_{x}$ coincide. Thus we can assume that $M_{x}$ is absolutely continuous w.r.t. $\sigma$ and replace the measure $\tau_{x}$ of the preceding sections by $\sigma$. For simplicity, now denote by $W_{x}$ the function $W_{x}=d M_{x} / d \sigma$ and according to this notation define the corresponding objects $\mathscr{W}_{x}$ etc. of Sections 3 and 4.

Let us set

$$
\begin{aligned}
L_{1} & :=\left\{u \in C^{q}: u^{*} W_{x}^{+} u \text { is integrable w.r.t. } \sigma\right\} \\
L_{2} & :=\left\{u \in C^{q}: u \in \mathscr{R}\left(W_{x}\right) \sigma \text {-a.e. }\right\} \\
L & :=L_{1} \cap L_{2}
\end{aligned}
$$

Remark 5.2. Note that the space $L$ coincides with the space $\mathscr{M}$ which appeared in Theorem 4.5 of [6] and was identified there as the range of the Grammian interpolation error matrix. Note further that $L$ is the orthogonal complement of the space $H$ of Lemma 9 of [5].

Let

$$
W_{x}=\left(\begin{array}{ll}
W_{x, 11} & W_{x, 12} \\
W_{x, 12}^{*} & W_{x, 22}
\end{array}\right)
$$

be the block representation of $W_{x}$ w.r.t. the orthogonal decomposition $C^{q}=L \oplus L^{\perp}$. Set

$$
W^{(r)}:=\left(\begin{array}{cc}
\left(W_{x} / W_{x, 22}\right) & 0 \\
0 & 0
\end{array}\right), \quad W^{(s)}:=\left(\begin{array}{cc}
W_{x, 12} W_{x, 22}^{+} W_{x, 12}^{*} & W_{x, 12} \\
W_{x, 12}^{*} & W_{x, 22}
\end{array}\right) .
$$

Using Theorem 4.4 we will show that $W^{(r)} d \sigma$ is the spectral measure of the $\mathscr{J}_{0}$-regular part of the Wold decomposition of $x$, and hence $W^{(s)}=W_{x}-W^{(r)}$ is the spectral measure of the $\mathscr{F}_{0}$-singular part.

Lemma 5.3. The spaces $\mathscr{R}\left(W^{(r)}\right)$ are equal to $L \sigma$-a.e.
Proof. Clearly, $\mathscr{R}\left(W^{(r)}\right) \subseteq L$. On the other hand, $L \subseteq \mathscr{R}\left(W_{x}\right)=$ $\mathscr{R}\left(W^{(r)}\right)+\mathscr{R}\left(W^{(s)}\right)$. Thus, if $L$ were not a subspace of $\mathscr{R}\left(W^{(r)}\right)$, we would have $\mathscr{R}\left(W^{(s)}\right) \cap L \neq\{0\}$. However, using (i) of Lemma 2.2 we easily get $\mathscr{R}\left(W^{(s)}\right) \cap L=\{0\}$. It follows that $\mathscr{R}\left(W^{(r)}\right)=L \sigma$-a.e.

Lemma 5.4. The function $W^{(r)+}$ is integrable w.r.t. $\sigma$.
Proof. Since $L \subseteq \mathscr{R}\left(W_{x}\right)$, we have ker $W_{x} \subseteq L^{\perp}$, and taking into account (i) of Lemma 2.2 we easily obtain $\operatorname{ker} W_{x}=\operatorname{ker} W_{22}$. Thus the generalized Banachiewicz inversion formula (cf. [8], formula (3.32)) is applicable, which implies that the left upper corner of $W_{x}^{+}$is equal to $\left(W_{x} / W_{x, 22}\right)^{-1}$. From the definition of $L$ it follows that $\left(W_{x} / W_{x, 22}\right)^{-1}$ is integrable w.r.t. $\sigma$ and so is $W^{(r)+}$.

Lemma 5.5. The space $L^{2}\left(W^{(r)}\right)$ is $\mathscr{J}_{0}$-regular.
Proof. The result follows immediately from Theorem 5.1 and Lemmas 5.3 and 5.4.

Lemma 5.6. Let $W \in \mathscr{W}_{x}^{(r)}$. Then $W \leqslant W^{(r)}$.
Proof. According to Theorem 5.1 there exists a subspace $L_{0}$ of $C^{q}$ such that $\mathscr{R}(W)=L_{0} \sigma$-a.e. Assume that $u \in L_{0} \cap L^{\perp}, u \neq 0$. Then $u$ can be written as $u=u_{1}+u_{2}$ for some $u_{1} \in L_{1}^{\frac{1}{1}}, u_{2} \in L_{2}^{\frac{1}{2}}$. If $u_{1}=0$, there exists $\Delta \in \mathscr{B}(\Gamma)$ such that $\sigma(\Delta)>0$ and $u=u_{2} \notin \mathscr{R}\left(W_{x}(\gamma)\right)$ for $\sigma$-a.a. $\gamma \in \Delta$. This contradicts the inclusion $\mathscr{R}(W) \subseteq \mathscr{R}\left(W_{x}\right) \sigma$-a.e. It follows that $u_{1} \neq 0$, and hence $u \notin L_{1}$. From the definition of $L_{1}$ we infer that $u^{*} W_{x}^{+} u$ is not integrable w.r.t. $\sigma$. Let $W_{0}$ be the restriction of $W$ to $L_{0}$ and let

$$
W_{x}=\left(\begin{array}{ll}
W_{x, 11}^{(0)} & W_{x, 12}^{(0)} \\
W_{x, 12}^{(0) *} & W_{x, 22}^{(0)}
\end{array}\right)
$$

be the block representation of $W_{x}$ w.r.t. the orthogonal decomposition $\boldsymbol{C}^{q}=L_{0} \oplus L_{0}^{\perp}$. From the definition of $\mathscr{W}_{x}^{(r)}$ and (iii) of Lemma 2.2 we obtain $W_{0} \leqslant\left(W_{x} / W_{x, 22}^{(0)}\right)$, and hence $\left(W_{x} / W_{x, 22}^{(0)}\right)^{-1} \leqslant W_{0}^{-1}$. By the generalized Banachiewicz inversion formula it follows that

$$
u^{*} W_{x}^{+} u \leqslant u^{*}\left(W_{x} / W_{x, 22}^{(0)}\right)^{-1} u \leqslant u^{*} W_{0}^{-1} u=u^{*} W^{+} u .
$$

Thus $u^{*} W^{+} u$ is not integrable w.r.t. $\sigma$, which contradicts Theorem 5.1. We conclude that $L_{0} \subseteq L$. Then again the definition of $\mathscr{W}_{x}^{(r)}$ and (iii) of Lemma 2.2 imply that the restriction of $W$ to $L$ does not exceed ( $W_{x} / W_{x, 22}$ ), which yields $W \leqslant W^{(r)} \sigma$-a.e.

Combining Lemmas 5.5 and 5.6 with Theorem 4.4 we get the following result.

Theorem 5.7. The measures $W^{(r)} d \sigma$ and $W^{(s)} d \sigma$ are the spectral measures of the $\mathscr{J}_{0}$-regular and $\mathscr{\mathscr { O }}_{0}$-singular parts of the Wold decomposition of $x$, respectively.

## REFERENCES

[1] A. Albert, Conditions for positive and nonnegative definiteness in terms of pseudoinverses, SIAM J. Appl. Math. 17 (1969), pp. 434-440.
[2] N. Dunford and J. T. Schwartz, Linear Operators. Part I: General Theory, Pure Appl. Math., Vol. 7, Interscience Publishers, 2nd edition, New York 1964.
[3] W. Fieger, Die Anwendung einiger mass- und integrationstheoretischer Sätze auf matrizielle Riemann-Stieltjes-Integrale, Math. Ann. 150 (1963), pp. 387-410.
[4] P. R. Halmos, A Hilbert Space Problem Book, Van Nostrand, Princeton 1967.
[5] L. Klotz and F. Schmidt, Some remarks on $J_{0}$-regularity and $J_{0}$-singularity of $q$-variate stationary processes, Probab. Math. Statist. 18 (1998), pp. 351-357.
[6] A. Makagon and A. Weron, q-variate minimal stationary processes, Studia Math. 59 (1976), pp. 41-52.
[7] A. Makagon and A. Weron, Wold-Cramér concordance theorems for interpolation of $q$-variate stationary processes over locally compact abelian groups, J. Multivariate Anal. 6 (1976), pp. 123-137.
[8] R. M. Pringle and A. A. Rayner, Generalized Inverse Matrices with Applications to Statistics, Griffin, London 1971.
[9] J. B. Robertson, Orthogonal decomposition of multivariate weakly stationary stochastic processes, Canad. J. Math. 20 (1968), pp. 368-383.
[10] M. Rosenberg, The square-integrability of matrix-valued functions with respect to a non--negative Hermitian measure, Duke Math. J. 31 (1964), pp. 291-298.
[11] Yu. A. Rozanov, Stationary Random Processes (in Russian), Fizmatgiz, Moscow 1963.
[12] H. Salehi and J. K. Scheidt, Interpolation of q-variate weakly stationary stochastic processes over a locally compact abelian group, J. Multivariate Anal. 2 (1972), pp. 307-331.
[13] S. R. Treil', Geometric methods in spectral theory of operator-valued functions: Some recent results, Operator Theory Vol. 42, Toeplitz Operators and Spectral Function Theory, N. K. Nikolskiĭ (Ed.), Birkhäuser, Basel-Boston-Berlin 1989, pp. 209-280.

Fakultät für Mathematik und Informatik
Universität Leipzig
04109 Leipzig, Germany

