# Multiblock Problems for Almost Periodic Matrix Functions of Several Variables 

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#### Abstract

In this paper we solve positive and contractive multiblock problems in the Wiener algebra of almost periodic functions of several variables. We thus generalize the classical four block problem that appears in robust control in many ways. The necessary and sufficient conditions are in terms of appropriate Toeplitz (positive case) and Hankel operators (contractive case) on Besikovitch space. In addition, a model matching interpretation is given, and some more general patterns are treated as well.


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## 1. Introduction

The suboptimal four block problem can be stated as follows. We let $L_{\infty}$ and $H_{\infty}$ denote the Lebesgue space and Hardy space, with respect to the essential supremum norm $\|\cdot\|_{\infty}$, on $\mathbb{T}=\{z \in \mathbb{C}:|z|=1\}$, respectively. We denote by $X^{p \times q}$ the set of $p \times q$ matrices with entries in the set $X$. Let $f_{11} \in L_{\infty}^{n_{1} \times m_{1}}, f_{12} \in L_{\infty}^{n_{1} \times m_{2}}$, $f_{21} \in L_{\infty}^{n_{2} \times m_{1}}$, and $f_{22} \in L_{\infty}^{n_{2} \times m_{2}}$ be given. Find, if possible, $\phi \in H_{\infty}^{n_{1} \times m_{2}}$ so that

$$
\sup _{|z|=1}\left\|\left[\begin{array}{cc}
f_{11}(z) & f_{12}(z)+\phi(z)  \tag{1.1}\\
f_{21}(z) & f_{22}(z)
\end{array}\right]\right\|<1
$$

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Here $\|\cdot\|$ is the operator norm (=the largest singular value) of a matrix.
This problem, as well as many of its variations, including the optimal four block problem (allowing equality in (1.1)) is ubiquitous in robust, or $H_{\infty}$, control. Its solution is typically given in terms of appropriately defined Hankel-type operators. The mathematical and engineering literature on the four block problem is extensive, and we quote here only a representative sample of books [6], [13], [2], [7], [8], where further information and references can be found.

The present paper grew out of the authors' desire to provide a solution to a multivariable almost periodic analogue of the four block problem. In the multivariable setting we impose a linear order on $\mathbb{R}^{k}$ through the use of a halfspace. This idea goes back to the seminal papers [11], [12] in which it was recognized that several of the results in one-variable Hardy space theory carry over to the setting of Hardy spaces on linearly ordered groups.

A natural starting point for solving the four block problem in our setting is to consider a related positive definite extension problem. While we treat the problem in great generality in subsequent sections, it may be instructive to state here a new problem that we solve which is most closely related to (1.1). This concerns $4 \times 4$ block matrix valued functions with entries in the Wiener algebra

$$
\mathcal{W}=\left\{f(z)=\sum_{i=-\infty}^{\infty} f_{i} z^{i}: \sum_{i=-\infty}^{\infty}\left\|f_{i}\right\|<\infty\right\}
$$

on the unit circle. Throughout the paper, a (bounded linear) Hilbert space operator $T: \mathcal{B} \rightarrow \mathcal{B}$ is called positive definite (notation: $T>0$ ) if $\langle T x, x\rangle \geq \epsilon\langle x, x\rangle$ for every $x \in \mathcal{B}$, where $\epsilon>0$ is independent of $x$, and where $\langle x, y\rangle$ is the inner product in the Hilbert space $\mathcal{B}$.

Theorem 1.1. Let $k_{i j} \in \mathcal{W}^{n_{i} \times n_{j}}, i, j=1,2,3,4$ be given. There exist $\phi \in \mathcal{W}^{n_{1} \times n_{4}}$ with $\phi(z)=\sum_{i=0}^{\infty} \phi_{i} z^{i}$ so that

$$
k_{\mathrm{ext}}(z)=\left[\begin{array}{cccc}
k_{11}(z) & k_{12}(z) & k_{13}(z) & k_{14}(z)+\phi(z) \\
k_{21}(z) & k_{22}(z) & k_{23}(z) & k_{24}(z) \\
k_{31}(z) & k_{32}(z) & k_{33}(z) & k_{34}(z) \\
k_{41}(z)+\phi(z)^{*} & k_{42}(z) & k_{43}(z) & k_{44}(z)
\end{array}\right]>0,|z|=1
$$

if and only if the following Toeplitz-like operator $T: \mathcal{B}_{2} \rightarrow \mathcal{B}_{2}$ is positive definite. Here the Hilbert space $\mathcal{B}_{2}$ is given by

$$
\mathcal{B}_{2}=\left[\begin{array}{cccc}
H_{2}^{n_{1} \times n_{1}} & L_{2}^{n_{1} \times n_{2}} & L_{2}^{n_{1} \times n_{3}} & L_{2}^{n_{1} \times n_{4}} \ominus H_{2}^{n_{1} \times n_{4}} \\
0 & H_{2}^{n_{2} \times n_{2}} & L_{2}^{n_{2} \times n_{3}} & L_{2}^{n_{2} \times n_{4}} \\
0 & 0 & H_{2}^{n_{3} \times n_{3}} & L_{2}^{n_{3} \times n_{4}} \\
0 & 0 & 0 & H_{2}^{n_{4} \times n_{4}}
\end{array}\right]
$$

where $L_{2}$ and $H_{2}$ are the Lebesgue and Hardy space on $\mathbb{T}$, respectively, with respect to the Hilbert space norm,

$$
T(g)=P_{\mathcal{B}_{2}}\left(\left(k_{i j}\right)_{i, j=1}^{4} g\right), \quad g=\left(g_{i j}\right)_{i, j=1}^{4} \in \mathcal{B}_{2}
$$

and $P_{\mathcal{B}_{2}}$ is the orthogonal projection of $L_{2}^{\left(n_{1}+n_{2}+n_{3}+n_{4}\right) \times\left(n_{1}+n_{2}+n_{3}+n_{4}\right)}$ onto $\mathcal{B}_{2}$.

The Hilbert space structure in $\mathcal{B}_{2}$ is the standard one: If $\mathcal{X}_{i, j}, i=1, \ldots, p$, $j=1, \ldots, q$ are Hilbert spaces, then the matrix $\mathcal{X}=\left[\mathcal{X}_{i, j}\right]_{i, j=1}^{p, q}$ is a Hilbert space with the inner product

$$
\left\langle\left[x_{i, j}\right]_{i, j=1}^{p, q},\left[y_{i, j}\right]_{i, j=1}^{p, q}\right\rangle=\sum_{i=1}^{p} \sum_{j=1}^{q}\left\langle x_{i, j}, y_{i, j}\right\rangle_{\mathcal{X}_{i, j}}, \quad x_{i, j}, y_{i, j} \in \mathcal{X}_{i, j}
$$

where $\langle\cdot, \cdot\rangle_{\mathcal{X}_{i, j}}$ is the inner product in $\mathcal{X}_{i, j}$.
The above result is a special case of Theorem 2.3 corresponding to $p=q=3$ and periodic functions defined on $\mathbb{R}$ (as opposed to almost periodic functions defined on $\mathbb{R}^{k}$ ). Observe also that Theorem 2.3 deals with arbitrary representations of $\mathbb{R}^{k}$ as unions of two halfspaces that overlap in the origin; for $k=1$ there is essentially one such representation: $\mathbb{R}=[0, \infty) \cup(-\infty, 0]$. Theorem 2.3 also gives (i) a construction of a specific $k_{\text {ext }}(z)$, that enjoys a maximum entropy; and (ii) a linear fractional description of the set of all solutions.

Using Theorem 1.1 one may solve a four block problem by letting

$$
\begin{gathered}
k_{i i}(z)=I_{n_{i}}, \quad i=1,2,3,4 \\
k_{12}(z)=k_{21}(z)=k_{34}(z)=k_{43}(z)=0 \\
k_{13}(z)=k_{31}(z)^{*}=f_{11}(z), \quad k_{14}(z)=k_{41}(z)^{*}=f_{12}(z) \\
k_{23}(z)=k_{32}(z)^{*}=f_{21}(z), \quad k_{24}(z)=k_{42}(z)^{*}=f_{22}(z)
\end{gathered}
$$

Note that this particular four block problem is in the Wiener algebra rather than in $L_{\infty}$.

Our main results are stated and proved in the framework of algebras of almost periodic matrix functions of several variables with Fourier spectrum in a given subgroup (see the next section for definitions and basic properties). This generality allows us to treat at once many particular situations, including the periodic case. Related problems of factorizations, positive and contractive extensions (which may be termed "one-block" problems in the context of the present paper) have been studied and solved in [26], [24], [25] for almost periodic matrix functions of several variables, and in earlier papers [28], [23] for almost periodic scalar and matrix functions of one variable. This development was largely motivated by recent applications of almost periodic functions in convolution equations on finite intervals, inverse scattering, and stochastic processes. We use in the present paper several key results proved in [26], [24] and [25].

Abstract band methods served as key technical tools in [28], [23], [26], [24], [25]. A standard exposition of the abstract band method is found in Part IX of [10]. Another version of the abstract band method was developed and used in [25]. In the present paper the abstract band method plays a key role as well. In Sections 2 and 3 we use the standard abstract band method (see, e.g., Part IX of [10]), while in Section 5 we use the more general version constructed in [15].

The paper is organized as follows. In Section 2 we formulate and solve a general banded almost periodic several variables positive extension problem of which Theorem 1.1 is a very particular case. In Section 3 we apply the results of Section 2 to obtain solutions to a general contractive extension problem. The relation of the results of Section 3 to those of Section 2 are based on the same principle that relates the four block problem to the setting of Theorem 1.1. A class of problems of model matching type is studied in Section 4. In Section 5 we shall solve positive extensions
problems for a non-banded pattern, namely, a specific example of a whole class of positive extension problems for almost periodic functions of several variables.

We use the standard notation $\mathbb{C}, \mathbb{R}, \mathbb{T}, \mathbb{Z}$, and $\mathbb{N}$ for the sets of complex numbers, of real numbers, of unimodular complex numbers, of integers, and of positive integers, respectively.

## 2. The positive multiblock problem

We first introduce the concepts and notation concerning almost periodic functions.

We let $\left(A P^{k}\right)$ denote the algebra of complex valued almost periodic functions of $k$ real variables, i.e., the closed subalgebra of $L^{\infty}\left(\mathbb{R}^{k}\right)$ (with respect to the standard Lebesgue measure) generated by all the functions $e_{\lambda}(t)=e^{i\langle\lambda, t\rangle}$, where $\lambda=\left(\lambda_{1}, \cdots, \lambda_{k}\right) \in \mathbb{R}^{k}$. Here the variable $t=\left(t_{1}, \cdots, t_{k}\right) \in \mathbb{R}^{k}$, and

$$
\langle\lambda, t\rangle=\sum_{j=1}^{k} \lambda_{j} t_{j}
$$

is the standard inner product of $\lambda$ and $t$. The norm in $\left(A P^{k}\right)$ will be denoted by $\|\cdot\|_{\infty}$. Recall that for any $f \in\left(A P^{k}\right)$ its Fourier series is defined by the formal sum

$$
\begin{equation*}
\sum_{\lambda} f_{\lambda} e^{i\langle\lambda, t\rangle} \tag{2.1}
\end{equation*}
$$

where

$$
\begin{equation*}
f_{\lambda}=\lim _{T \rightarrow \infty} \frac{1}{(2 T)^{k}} \int_{[-T, T]^{k}} e^{-i\langle\lambda, t\rangle} f(t) d t, \lambda \in \mathbb{R}^{k} \tag{2.2}
\end{equation*}
$$

and the sum in (2.1) is taken over the set $\sigma(f)=\left\{\lambda \in \mathbb{R}^{k}: f_{\lambda} \neq 0\right\}$, called the Fourier spectrum of $f$. The Fourier spectrum of every $f \in\left(A P^{k}\right)$ is at most a countable set. The mean $M\{f\}$ of $f \in\left(A P^{k}\right)$ is defined by $M\{f\}=f_{0}=$ $\lim _{T \rightarrow \infty} \frac{1}{(2 T)^{k}} \int_{[-T, T]^{k}} f(t) d t$. The Wiener algebra $\left(A P W^{k}\right)$ is defined as the set of all $f \in\left(A P^{k}\right)$ such that the Fourier series of $f$ converges absolutely. The Wiener algebra is a Banach $*$-algebra with respect to the Wiener norm $\|f\|_{W}=\sum_{\lambda \in \mathbb{R}^{k}}\left|f_{\lambda}\right|$ (the multiplication in $\left(A P W^{k}\right)$ and the involution are defined pointwise). Note that $\left(A P W^{k}\right)$ is dense in $\left(A P^{k}\right)$.

Denote by $\left(A P^{k}\right)^{m \times n}$ (resp., $\left(A P W^{k}\right)^{m \times n}$ ) the set, which is an algebra if $m=n$, of $m \times n$ matrices with entries in $\left(A P^{k}\right)$ (resp., $\left(A P W^{k}\right)$ ). The infinity norm in $\left(A P^{k}\right)^{m \times n}$ is

$$
\|f\|_{\infty}=\sup _{t \in \mathbb{R}^{k}}\|f(t)\|, \quad f \in\left(A P^{k}\right)^{m \times n}
$$

where $\|\cdot\|$ is the operator norm. The mean $M\{F\}$ for $F=\left[F_{p, q}\right]_{p=1, q=1}^{m, n} \in$ $\left(A P^{k}\right)^{m \times n}$ is defined by $M\{F\}=\left[M\left\{F_{p, q}\right\}\right]_{p=1, q=1}^{m, n}$. The Fourier spectrum $\sigma(F)$ of $F \in\left(A P^{k}\right)^{m \times n}$ is, by definition, the set of all $\lambda \in \mathbb{R}^{k}$ such that $M\left\{e_{-\lambda} F\right\} \neq 0$. In other words, $\sigma\left(\left[F_{p, q}\right]_{p=1, q=1}^{m, n}\right)=\cup_{p=1, q=1}^{m, n} \sigma\left(F_{p, q}\right)$.

A matrix function $f \in\left(A P^{k}\right)^{n \times n}$ is called positive definite if $f(t)$ is Hermitian for every $t \in \mathbb{R}^{k}$, and there exists an $\epsilon>0$ such that $f(t) \geq \epsilon I_{n}$ for all $t \in \mathbb{R}^{k}$, where $I_{n}$ is the $n \times n$ identity matrix. For the general theory of almost periodic
functions of one and several variables we refer the reader to the books $[3,19,20]$ and to Chapter 1 in [22].

Let $\Delta$ be a non-empty subset of $\mathbb{R}^{k}$. Denote

$$
\begin{equation*}
\left(A P^{k}\right)_{\Delta}=\left\{f \in\left(A P^{k}\right): \sigma(f) \subseteq \Delta\right\},\left(A P W^{k}\right)_{\Delta}=\left\{f \in\left(A P W^{k}\right): \sigma(f) \subseteq \Delta\right\} \tag{2.3}
\end{equation*}
$$

If $\Delta$ is an additive subgroup of $\mathbb{R}^{k}$, then $\left(A P^{k}\right)_{\Delta}\left(\right.$ resp. $\left.\left(A P W^{k}\right)_{\Delta}\right)$ is a unital subalgebra of $\left(A P^{k}\right)$ (resp. $\left(A P W^{k}\right)$ ).

Introduce an inner product on $\left(A P^{k}\right)$ by the formula

$$
\begin{equation*}
\langle f, g\rangle=M\left\{f g^{*}\right\}, f, g \in\left(A P^{k}\right) \tag{2.4}
\end{equation*}
$$

Here and elsewhere we denote by $g^{*}$ the function $g^{*}(t)=\overline{g(t)}$; if $g \in\left(A P^{k}\right)^{m \times n}$, then $g^{*} \in\left(A P^{k}\right)^{n \times m}$ is defined by $g^{*}(t)=\overline{g(t)^{T}}$, where ${ }^{T}$ designates the transposed matrix. The completion of $\left(A P^{k}\right)$ with respect to inner product (2.4) is called the Besikovitch space and is denoted by $\left(B^{k}\right)$. Thus $\left(B^{k}\right)$ is a Hilbert space. For a nonempty set $\Lambda \subseteq \mathbb{R}^{k}$, define the projection

$$
P_{\Lambda}\left(\sum_{\lambda \in \sigma(f)} f_{\lambda} e^{i\langle\lambda, t\rangle}\right)=\sum_{\lambda \in \sigma(f) \cap \Lambda} f_{\lambda} e^{i\langle\lambda, t\rangle}
$$

where $f \in\left(A P W^{k}\right)$. The projection $P_{\Lambda}$ extends by continuity to the orthogonal projection (also denoted $P_{\Lambda}$ ) on $\left(B^{k}\right)$. We denote by $\left(B^{k}\right)_{\Lambda}$ the range of $P_{\Lambda}$, or, equivalently, the completion of $\left(A P^{k}\right)_{\Lambda}$ with respect to the inner product (2.4). The matrix valued Besikovitch space $\left(B^{k}\right)^{n \times m}$ consists of $n \times m$ matrices with components in $\left(B^{k}\right)$, with the standard Hilbert space structure:

$$
\begin{equation*}
\left\langle\left(f_{i, j}\right)_{i=1, j=1}^{n, m},\left(g_{i, j}\right)_{i=1, j=1}^{n, m}\right\rangle=\sum_{i=1}^{n} \sum_{j=1}^{m}\left\langle f_{i, j}, g_{i, j}\right\rangle . \tag{2.5}
\end{equation*}
$$

Similarly, $\left(B^{k}\right)_{\Lambda}^{n \times m}$ is the Hilbert space of $n \times m$ matrices with components in $\left(B^{k}\right)_{\Lambda}$. In the periodic case $\left(\Lambda=\mathbb{Z}^{k}\right)$ we may identify $\left(B^{k}\right)_{\Lambda}$ with $L_{2}\left(\mathbb{T}^{k}\right)$.

A subset $S$ of $\mathbb{R}^{k}$ is called a halfspace if it has the following properties:
(i) $\mathbb{R}^{k}=S \cup(-S)$;
(ii) $S \cap(-S)=\{0\}$;
(iii) if $x, y \in S$ then $x+y \in S$;
(iv) if $x \in S$ and $\alpha$ is a nonnegative real number, then $\alpha x \in S$.

Note that conditions (iii) and (iv) mean that $S$ is a cone, and conditions (ii), (iii), and (iv) together mean that $S$ is a pointed cone. A standard example of a halfspace is given by

$$
\begin{gathered}
E_{k}=\left\{\left(x_{1}, \cdots, x_{k}\right)^{T} \in \mathbb{R}^{k} \backslash\{0\}:\right. \\
\left.x_{1}=x_{2}=\cdots=x_{j-1}=0, x_{j} \neq 0 \Rightarrow x_{j}>0\right\} \cup\{0\} .
\end{gathered}
$$

(The vectors in $\mathbb{R}^{k}$ are understood as column vectors.) One can show using basic results on linearly ordered real vector spaces (see [5], Section IV. 5 in [9]) that a set $S \subseteq \mathbb{R}^{k}$ is a halfspace if and only if there exists a real invertible $k \times k$ matrix $A$ such that

$$
\begin{equation*}
S=A E_{k} \stackrel{\text { def }}{=}\left\{A x: x \in E_{k}\right\} \tag{2.6}
\end{equation*}
$$

See Section 2 of [26] for more details.

In this section we consider the following positive extension problem. Let $S$ be a halfspace in $\mathbb{R}^{k}$ and $\Lambda$ an additive subgroup of $\mathbb{R}^{k}$. Further, fix

$$
p \in \mathbb{N}, \quad q \in\{1, \ldots, p\}, \quad \text { and } \quad n_{0}, n_{1}, \ldots, n_{p} \in \mathbb{N}
$$

Given are:

$$
f_{i j}=f_{j i}^{*} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times n_{j}} \quad \text { for } \quad|j-i|<q
$$

and

$$
f_{i j}^{-}=\left(f_{j i}^{+}\right)^{*} \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{i} \times n_{j}}, \quad \text { for } \quad j-i=q
$$

Find, if possible,

$$
f_{i j}^{+}=\left(f_{j i}^{-}\right)^{*} \in\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{i} \times n_{j}}, \quad \text { for } \quad j-i=q
$$

and

$$
f_{i j}=f_{j i}^{*} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times n_{j}} \quad \text { for } \quad|j-i|>q
$$

so that

$$
\begin{equation*}
\left(f_{i j}\right)_{i, j=0}^{p}>0, \tag{2.7}
\end{equation*}
$$

where

$$
f_{i j}=f_{j i}^{*}:=f_{i j}^{-}+f_{i j}^{+} \quad \text { for } \quad j-i=q
$$

The inequality in (2.7) is interpreted as follows:

$$
\left(f_{i j}\right)_{i, j=0}^{p}(t) \geq \epsilon I, \text { for every } t \in \mathbb{R}^{k}
$$

where $\epsilon>0$ is independent of $t$. An analogous interpretation is given to inequalities of the form $G>0$ in the sequel. The necessary and sufficient condition for existence of a solution (2.7) will be the positive definiteness of an appropriately defined Toeplitz operator. In addition, in that case we shall construct a solution to the problem that has a certain maximality property, and we shall construct a linear fractional description for the set of solutions as well.

Before we can state and prove our theorems we do have to introduce some additional notation and derive some auxiliary results. We start by quoting a result from [26] (Corollary 5.3) that will be used extensively in the sequel.
Theorem 2.1. Let $G \in\left(A P W^{k}\right)^{n \times n}$ and assume that the matrix $G(t)$ is positive definite for every $t \in \mathbb{R}^{k}$, and $\operatorname{det} G(t) \geq \epsilon$ for every $t \in \mathbb{R}^{k}$, where $\epsilon>0$ is independent of $t$. Let also $\Lambda^{\prime}$ be the minimal additive subgroup of $\mathbb{R}^{k}$ which contains $\sigma(G)$. Then $G(t)$ admits canonical factorizations of the forms

$$
\begin{equation*}
G(t)=A_{+}(t) \cdot\left(A_{+}(t)\right)^{*}=\left(\widetilde{A}_{+}(t)\right)^{*} \cdot \widetilde{A}_{+}(t) \tag{2.8}
\end{equation*}
$$

where $A_{+}^{ \pm 1}, \widetilde{A}_{+}^{ \pm 1} \in\left(A P W^{k}\right)_{S \cap \Lambda^{\prime}}^{n \times n}$. The factors $A_{+}^{ \pm 1}, \widetilde{A}_{+}^{ \pm 1}$ are defined up to a right/left constant unitary multiple.

Let

$$
\begin{aligned}
& \mathcal{M}=\left(A P W^{k}\right)_{\Lambda}^{N \times N}, \quad N=n_{0}+\cdots+n_{p} \\
& \mathcal{M}_{1}=\mathcal{M}_{4}^{*}=\left\{\left(f_{i j}\right)_{i, j=0}^{p} \in \mathcal{M}: f_{i j}=0,\right. j-i<q \\
&\text { and } \left.f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{i} \times n_{j}}, j-i=q\right\},
\end{aligned}
$$

$$
\begin{gathered}
\mathcal{M}_{2}^{0}=\left(\mathcal{M}_{3}^{0}\right)^{*}=\left\{\left(f_{i j}\right)_{i, j=0}^{p} \in \mathcal{M}: f_{i j}=0, j-i<0 ; f_{i i} \in\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{i} \times n_{i}},\right. \\
\\
\left.i=0, \ldots, p ; f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{i} \times n_{j}}, j-i=q ; \text { and } f_{i j}=0, j-i>q\right\},
\end{gathered}
$$

and

$$
\mathcal{M}_{d}=\mathbb{C}^{n_{0} \times n_{0}} \oplus \cdots \oplus \mathbb{C}^{n_{p} \times n_{p}}
$$

We identify here $\mathbb{C}^{n_{j} \times n_{j}}$ with $\left(A P W^{k}\right)_{\{0\}}^{n_{j} \times n_{j}}$. Then

$$
\begin{gather*}
e:=\bigoplus_{i=0}^{p} I_{n_{i} \times n_{i}} \in \mathcal{M}_{d}, \\
\mathcal{M}=\mathcal{M}_{1}+\mathcal{M}_{2}^{0}+\mathcal{M}_{d}+\mathcal{M}_{3}^{0}+\mathcal{M}_{4}, \tag{2.9}
\end{gather*}
$$

and the following multiplication table holds:

| $\cdot$ | $\mathcal{M}_{1}$ | $\mathcal{M}_{2}^{0}$ | $\mathcal{M}_{d}$ | $\mathcal{M}_{3}^{0}$ | $\mathcal{M}_{4}$ |
| :---: | :--- | :--- | :--- | :--- | :--- |
| $\mathcal{M}_{1}$ | $\mathcal{M}_{1}$ | $\mathcal{M}_{1}$ | $\mathcal{M}_{1}$ | $\mathcal{M}_{+}^{0}$ | $\mathcal{M}$ |
| $\mathcal{M}_{2}^{0}$ | $\mathcal{M}_{1}$ | $\mathcal{M}_{+}^{0}$ | $\mathcal{M}_{2}^{0}$ | $\mathcal{M}_{c}$ | $\mathcal{M}_{-}^{0}$ |
| $\mathcal{M}_{d}$ | $\mathcal{M}_{1}$ | $\mathcal{M}_{2}^{0}$ | $\mathcal{M}_{d}$ | $\mathcal{M}_{3}^{0}$ | $\mathcal{M}_{4}$ |
| $\mathcal{M}_{3}^{0}$ | $\mathcal{M}_{+}^{0}$ | $\mathcal{M}_{c}$ | $\mathcal{M}_{3}^{0}$ | $\mathcal{M}_{-}^{0}$ | $\mathcal{M}_{4}$ |
| $\mathcal{M}_{4}$ | $\mathcal{M}$ | $\mathcal{M}_{-}^{0}$ | $\mathcal{M}_{4}$ | $\mathcal{M}_{4}$ | $\mathcal{M}_{4}$ |

where

$$
\begin{aligned}
\mathcal{M}_{+}^{0} & =\left(\mathcal{M}_{-}^{0}\right)^{*}=\mathcal{M}_{2}^{0} \dot{+} \mathcal{M}_{1} \\
\mathcal{M}_{c} & =\mathcal{M}_{2}^{0}+\mathcal{M}_{d} \dot{+} \mathcal{M}_{3}^{0} \\
\mathcal{M}_{+} & =\mathcal{M}_{-}^{*}=\mathcal{M}_{+}^{0} \dot{+} \mathcal{M}_{d}
\end{aligned}
$$

For example, $\mathcal{M}_{2}^{0} \mathcal{M}_{4} \subseteq \mathcal{M}_{-}^{0}$. Note that $\mathcal{M}_{-}^{0}$ and $\mathcal{M}_{+}^{0}$ are independent of $q$. We shall use the letter $e$ to also refer to identity matrices of smaller size as well, e.g., $e=\bigoplus_{i=0}^{p-1} I_{n_{i}}$. The appropriate interpretation should be clear from the context. Notice that we may restate the positive extension problem formulated in this section as follows: Given $k=k^{*} \in \mathcal{M}_{c}$, find, if possible, $m_{1} \in \mathcal{M}_{1}$ so that

$$
k+m_{1}+m_{1}^{*}>0 .
$$

The element $k+m_{1}+m_{1}^{*}>0$ is called a positive extension of $k$, i.e., $f$ is a positive extension of $k$ if and only if $f>0$ and $f-k \in \mathcal{M}_{1} \dot{+} \mathcal{M}_{4}$. It should be noted that the problem is now completely formulated as in the general band method (see Chapter XXXIV of [10] or [29]), except that we have to show that the notion of positivity $f>0$ coincides with the notion of positivity as required in the general band method. This follows easily as a consequence of the next lemma.

Lemma 2.2. Let $f \in \mathcal{M}$ be so that $f>0$. Then there exist $f_{+} \in \mathcal{M}_{+}^{0}$ and $f_{d} \in \mathcal{M}_{d}$ so that

$$
\begin{equation*}
f=\left(e+f_{+}\right)^{*} f_{d}\left(e+f_{+}\right), \tag{2.11}
\end{equation*}
$$

$f_{d}>0, f_{d} \in \mathcal{M}_{d}, e+f_{+}$is invertible in $\mathcal{M}$ and $\left(e+f_{+}\right)^{-1}-e \in \mathcal{M}_{+}^{0}$. Moreover, this factorization is unique. In addition, $f$ also allows the factorization

$$
\begin{equation*}
f=\left(e+f_{-}\right)^{*} \widetilde{f}_{d}\left(e+f_{-}\right) \tag{2.12}
\end{equation*}
$$

with $\tilde{f}_{d}>0, \widetilde{f}_{d} \in \mathcal{M}_{d}$ and $f_{-},\left(e+f_{-}\right)^{-1}-e \in \mathcal{M}_{-}^{0}$. Again, factorization (2.12) is unique.

The factorizations (2.11) and (2.12) are referred to as the right and the left canonical factorizations of $f$, respectively, with respect to the decomposition $\mathcal{M}=$ $\mathcal{M}_{-}^{0} \dot{+} \mathcal{M}_{d} \dot{+} \mathcal{M}_{+}^{0}$.
Proof of Lemma 2.2. We do this by induction on the number of row and column blocks $p$. Note that the statement of this lemma also makes sense when $p=0$ (in contrast with other parts of this section). In this case we interpret

$$
\mathcal{M}_{-}^{0}=\left(A P W^{k}\right)_{\Lambda \cap(-S \backslash\{0\})}^{n_{0} \times n_{0}}, \quad \mathcal{M}_{d}=\mathbb{C}^{n_{0} \times n_{0}}, \quad \mathcal{M}_{+}^{0}=\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{0} \times n_{0}}
$$

It will be most convenient for the proof to include the case $p=0$.
When $p=0$, the result of Lemma 2.2 follows from Theorem 2.1. Now assume the result has been established up to $p-1$. Letting $f=\left(f_{i j}\right)_{i, j=0}^{p}$ and writing

$$
\left(f_{i j}\right)_{i, j=0}^{p-1}=\left(e+f_{+}^{(p-1)}\right)^{*} f_{d}^{(p-1)}\left(e+f_{+}^{(p-1)}\right)
$$

we get that the right canonical factors of $f$ in (2.11) are given by

$$
\begin{aligned}
e+f_{+} & :=\left[\begin{array}{cc}
e+f_{+}^{(p-1)} & 0 \\
0 & e+g_{+}
\end{array}\right]\left[\begin{array}{ll}
e & h \\
0 & e
\end{array}\right] \\
f_{d} & :=\left[\begin{array}{cc}
f_{d}^{(p-1)} & 0 \\
0 & g_{d}
\end{array}\right] .
\end{aligned}
$$

Here $h=\left[\left(f_{i j}\right)_{i, j=0}^{p-1}\right]^{-1}\left(f_{i p}\right)_{i=0}^{p-1}$, and

$$
f_{p p}-\left(f_{p i}\right)_{i=0}^{p-1}\left[\left(f_{i j}\right)_{i, j=0}^{p-1}\right]^{-1}\left(f_{i p}\right)_{i=0}^{p-1}=\left(e+g_{+}\right)^{*} g_{d}\left(e+g_{+}\right)
$$

is a right canonical factorization with respect to the decomposition

$$
\left(A P W^{k}\right)_{\Lambda \cap(-S \backslash\{0\})}^{n_{p} \times n_{p}} \dot{+} \mathbb{C}^{n_{p} \times n_{p}} \dot{+}\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{p} \times n_{p}} .
$$

Notice that

$$
f_{p p}-\left(f_{p i}\right)_{i=0}^{p-1}\left[\left(f_{i j}\right)_{i, j=0}^{p-1}\right]^{-1}\left(f_{i p}\right)_{i=0}^{p-1}>0 .
$$

The uniqueness of the factorization follows using standard arguments (see, e.g., [26]). The proof of (2.12) is analogous.

We shall formulate our necessary and sufficient condition in terms of a Hilbert space operator. We thus start with introducing the appropriate Hilbert spaces. Let $\mathcal{B}=\left(B^{k}\right)^{N \times N}$, and let

$$
\mathcal{B}_{1}, \mathcal{B}_{2}^{0}, \mathcal{B}_{2}, \mathcal{B}_{d}, \mathcal{B}_{3}^{0}, \mathcal{B}_{3}, \mathcal{B}_{4}, \mathcal{B}_{-}^{0}, \mathcal{B}_{+}^{0}, \mathcal{B}_{-}, \mathcal{B}_{+}, \mathcal{B}_{d}
$$

be the closures of

$$
\begin{gathered}
\mathcal{M}_{1}, \mathcal{M}_{2}^{0}, \mathcal{M}_{2}=\mathcal{M}_{2}^{0} \dot{+} \mathcal{M}_{d}, \mathcal{M}_{d}, \mathcal{M}_{3}^{0}, \mathcal{M}_{3}=\mathcal{M}_{3}^{0} \dot{+} \mathcal{M}_{d} \\
\mathcal{M}_{4}, \mathcal{M}_{-}^{0}, \mathcal{M}_{+}^{0}, \mathcal{M}_{-}, \mathcal{M}_{+}, \mathcal{M}_{d}
\end{gathered}
$$

in $\mathcal{B}$, respectively. Alternatively, one may for example introduce $\mathcal{B}_{1}$ in the same way $\mathcal{M}_{1}$ was defined with $A P W^{k}$ replaced by $B^{k}$. Notice that we have the following orthogonal sum decomposition

$$
\mathcal{B}=\mathcal{B}_{1} \oplus \mathcal{B}_{2}^{0} \oplus \mathcal{B}_{d} \oplus \mathcal{B}_{3}^{0} \oplus \mathcal{B}_{4}
$$

In addition, we introduce $P_{1}, P_{2}^{0}$, etc., as the orthogonal projections of $\mathcal{B}$ onto $\mathcal{B}_{1}$, $\mathcal{B}_{2}^{0}$, etc. It is not hard to see that if one of these projections is applied to a member of $\mathcal{M}$ the result ends up in the appropriate subspace of $\mathcal{M}$. So, for instance, if $m \in \mathcal{M}$, then $P_{c} m \in \mathcal{M}_{c}$.

We will now state our main result. For $f \in\left(A P W^{k}\right)^{m \times m}$ we denote

$$
\|f\|_{\infty}=\sup _{x \in \mathbb{R}^{k}}\|f(x)\|,
$$

where $\|X\|$ is the operator norm of an operator or matrix $X$.
Theorem 2.3. Let $k=k^{*} \in \mathcal{M}_{c}$ be given. Then the following statements are equivalent:
(i) $k$ has a positive extension;
(ii) the operator $T: \mathcal{B}_{2} \rightarrow \mathcal{B}_{2}$ defined by $T(g)=P_{2}(k g)$, is positive definite;
(iii) the operator $\widetilde{T}: \mathcal{B}_{3} \rightarrow \mathcal{B}_{3}$ defined by $\widetilde{T}(g)=P_{3}(\mathrm{~kg})$, is positive definite.

In that case, let

$$
x=T^{-1}(e), \quad y=\widetilde{T}^{-1}(e) .
$$

Then $x^{-1} \in \mathcal{M}_{+}, y^{-1} \in \mathcal{M}_{-}, P_{d} x>0, P_{d} y>0$, and

$$
f_{0}:=x^{*-1}\left(P_{d} x\right) x^{-1}=y^{*-1}\left(P_{d} y\right) y^{-1}
$$

is a positive extension of $k$. In fact, $f_{0}$ is the unique positive extension of $k$ with $f_{0}^{-1} \in \mathcal{M}_{c}$. Moreover, if we let $u=x\left(P_{d} x\right)^{-\frac{1}{2}}$ and $v=y\left(P_{d} y\right)^{-\frac{1}{2}}$, then $f$ is a positive extension of $k$ if and only if

$$
f=(u+v g)^{*-1}\left(e-g^{*} g\right)(u+v g)^{-1}
$$

for some (unique!) $g \in \mathcal{M}_{1}$ with $\|g\|_{\infty}<1$. Also, $f$ is a positive extension of $k$ if and only if

$$
f=(v+u h)^{*-1}\left(e-h^{*} h\right)(v+u h)
$$

for some (unique!) $h \in \mathcal{M}_{4}$ with $\|h\|_{\infty}<1$.
Lastly, if $f$ is a positive extension of $k$ with canonical factorizations

$$
f=\left(e+f_{+}\right)^{*} f_{d}\left(e+f_{+}\right)=\left(e+f_{-}\right)^{*} \widetilde{f}_{d}\left(e+f_{-}\right)
$$

then

$$
\begin{equation*}
f_{d} \leq\left(P_{d} x\right)^{-1} \text { and } \widetilde{f}_{d} \leq\left(P_{d} y\right)^{-1} \tag{2.13}
\end{equation*}
$$

i.e., the matrices $\left(P_{d} x\right)^{-1}-f_{d}$ and $\left(P_{d} y\right)^{-1}-\widetilde{f}_{d}$ are positive semidefinite, and equality holds in at least one of the two inequalities in (2.13) if and only if $f=f_{0}$ (and thus both inequalities are equalities).

The "maximizing" property of $f_{0}$ exhibited in (2.13) may also be presented as follows. Let $f \in\left(A P W^{k}\right)^{N \times N}$ be positive definite. From Proposition 2.3 in [26], it follows that $\log (\operatorname{det} f)$ belongs to $\left(A P W^{k}\right)$. The number

$$
\begin{equation*}
\Delta(f):=M\{\log (\operatorname{det} f)\} \tag{2.14}
\end{equation*}
$$

will be referred to as the entropy of $f$.

Proposition 2.4. Let $k=k^{*} \in \mathcal{M}_{c}$ be given and suppose that $k$ has a positive extension. Let $f_{0}$ be defined as in Theorem 2.3. If $f$ is a positive extension of $k$ then

$$
\Delta(f) \leq \Delta\left(f_{0}\right)
$$

and equality holds if and only if $f=f_{0}$.
The main effort in proving Theorem 2.3 lies in showing that $x^{-1} \in \mathcal{M}_{+}$and $y^{-1} \in \mathcal{M}_{-}$provided $T>0$ and $\widetilde{T}>0$. The remainder of the proof is a standard application of the band method.

For the reader's convenience, we state a result (that includes a particular case of Theorem 2.1 of [24] and remarks thereafter) that will be used in the proof of Theorem 2.3.

For $g \in\left(A P W^{k}\right)_{\Lambda^{\prime}}^{m \times n}$, where $\Lambda^{\prime}$ is an additive subgroup of $\mathbb{R}^{k}$, and for $p \in \mathbb{Z}$ define the generalized Hankel operator

$$
\mathbf{H}(g)_{\Lambda^{\prime}}:\left(B^{k}\right)_{S \cap \Lambda^{\prime}}^{n \times p} \rightarrow\left(B^{k}\right)_{(-S) \cap \Lambda^{\prime}}^{m \times p}
$$

by

$$
\begin{equation*}
\mathbf{H}(g)_{\Lambda^{\prime}} h=P_{-S}(g h), \quad h \in\left(B^{k}\right)_{S \cap \Lambda^{\prime}}^{n \times p} . \tag{2.15}
\end{equation*}
$$

We suppress the dependence of $\mathbf{H}(g)_{\Lambda^{\prime}}$ on $S$ in our notation. It is not hard to see that the norm of $\mathbf{H}(g)_{\Lambda^{\prime}}$ is independent of the choice of the positive integer $p$.
Theorem 2.5. [24] Let $f \in\left(A P W^{k}\right)_{(-S) \cap \Lambda^{\prime}}^{m \times n}$ be given, where $\Lambda^{\prime}$ is an additive subgroup of $\mathbb{R}^{k}$. Then the following two statements are equivalent:
(i) $f$ has a strictly contractive extension $h \in\left(A P W^{k}\right)_{\Lambda^{\prime}}^{m \times n}$, i.e.,

$$
\|h\|_{\infty}:=\sup _{t \in \mathbb{R}^{k}}\|h(t)\|<1, \quad \text { and } \quad h_{\lambda}=f_{\lambda}, \quad \text { for } \quad \lambda \in(-S) \cap \Lambda^{\prime} .
$$

(ii) The generalized Hankel operator $\mathbf{H}(f)_{\Lambda^{\prime}}$ is a strict contraction.

When one (and thus both) of (i)-(ii) is satisfied, put

$$
\begin{aligned}
\widehat{\alpha}(t) & =\left[I-\mathbf{H}(f)_{\Lambda^{\prime}}\left(\mathbf{H}(f)_{\Lambda^{\prime}}\right)^{*}\right]^{-1}(e), \\
\widehat{\beta}(t) & =\mathbf{H}(f)_{\Lambda^{\prime}}\left[I-\left(\mathbf{H}(f)_{\Lambda^{\prime}}\right)^{*} \mathbf{H}(f)_{\Lambda^{\prime}}\right]^{-1}(e), \\
\widehat{\gamma}(t) & =\left(\mathbf{H}(f)_{\Lambda^{\prime}}\right)^{*}\left[I-\mathbf{H}(f)_{\Lambda^{\prime}}\left(\mathbf{H}(f)_{\Lambda^{\prime}}\right)^{*}\right]^{-1}(e), \\
\widehat{\delta}(t) & =\left[I-\left(\mathbf{H}(f)_{\Lambda^{\prime}}\right)^{*} \mathbf{H}(f)_{\Lambda^{\prime}}\right]^{-1}(e),
\end{aligned}
$$

where $e$ stands for matrix function on $\mathbb{R}^{k}$ that takes the constant identity matrix value. Then $\widehat{\alpha}$ is invertible in $\left(A P W^{k}\right)_{\Lambda^{\prime}}^{m \times m}$ and $M\{\widehat{\alpha}\}$ is positive definite. Similarly, $\widehat{\delta}$ is invertible in $\left(A P W^{k}\right)_{\Lambda^{\prime}}^{n \times n}$ and $M\{\widehat{\delta}\}>0$. Also, $\widehat{\alpha}^{-1} \in\left(A P W^{k}\right)_{(-S) \cap \Lambda^{\prime}}^{m \times m}$ and $\widehat{\delta}^{-1} \in\left(A P W^{k}\right)_{S \cap \Lambda^{\prime}}^{n \times n}$.

Further, let

$$
\begin{array}{ll}
\alpha(t)=\widehat{\alpha}(t) M\{\widehat{\alpha}\}^{-\frac{1}{2}}, & \beta(t)=\widehat{\beta}(t) M\{\widehat{\delta}\}^{-\frac{1}{2}} \\
\gamma(t)=\widehat{\gamma}(t) M\{\widehat{\alpha}\}^{-\frac{1}{2}}, & \delta(t)=\widehat{\delta}(t) M\{\widehat{\delta}\}^{-\frac{1}{2}} \tag{2.17}
\end{array}
$$

Then the function

$$
\begin{equation*}
h_{0}(t)=\beta(t) \delta(t)^{-1}=\left[\alpha(t)^{*}\right]^{-1} \gamma(t)^{*}, t \in \mathbb{R}^{k} \tag{2.18}
\end{equation*}
$$

is a strictly contractive extension in $\left(A P W^{k}\right)_{\Lambda^{\prime}}^{m \times n}$ of $f$.
Proposition 2.6. Let $T: \mathcal{B}_{2} \rightarrow \mathcal{B}_{2}$ be as in Theorem 2.3 and suppose that $T>0$. Then $x:=T^{-1}(e)$ is an invertible element of $\mathcal{M}_{+}$.
Proof. Since $x \in \mathcal{B}_{2}$ it has the form $x=\left(x_{i j}\right)_{i, j=0}^{p}$ with $x_{i j}=0$ for $j-i<0$; $x_{i i} \in\left(B^{k}\right)_{\Lambda \cap S}^{n_{i} \times n_{j}}$ for $i=0, \ldots, p ; x_{i j} \in\left(B^{k}\right)_{\Lambda}^{n_{i} \times n_{j}}$ for $0<j-i<q ; x_{i j} \in\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{i} \times n_{j}}$ for $j-i=q$; and $x_{i j}=0$ for $j-i>q$. It therefore suffices to show that $x_{i i}^{ \pm 1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{i} \times n_{i}}$ and $x_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times n_{j}}$ for $j-i>0$.

Denote $k=\left(k_{i j}\right)_{i, j=0}^{p}$. First observe that by applying $T$ to elements $g=$ $\left(g_{i j}\right)_{i, j=0}^{p} \in \mathcal{B}_{2}$ with $g_{i j}=0$ for $(i, j) \neq(0,0)$ we get that the operator $T_{00}$ : $\left(B^{k}\right)_{\Lambda \cap S}^{n_{0} \times n_{0}} \rightarrow\left(B^{k}\right)_{\Lambda \cap S}^{n_{0} \times n_{0}}$ defined by

$$
T_{00}\left(g_{00}\right)=P_{S}\left(k_{00} g_{00}\right)
$$

is positive definite. Here $P_{S}$ is the orthogonal projection of $\left(B^{k}\right)_{\Lambda}^{m \times m}$ onto $\left(B^{k}\right)_{\Lambda \cap S}^{m \times m}$. It follows from Section 5 in [26] that $k_{00}>0$, and therefore has a right canonical factorization (Theorem 2.1)

$$
k_{00}=\left(e+k_{00,+}\right)^{*} k_{00, d}\left(e+k_{00,+}\right) .
$$

One easily checks that the equation $x=T^{-1}(e)$, or equivalently $P_{2}(k x)=e$, implies that $x_{00}=\left(e+k_{00,+}\right)^{-1} k_{00, d}^{-1}$, and therefore $x_{00}^{ \pm 1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{0} \times n_{0}}$.

Next (assuming $q \geq 2$, otherwise this step is not necessary), observe that by applying $T$ to elements $g=\left(g_{i j}\right)_{i, j=0}^{p} \in \mathcal{B}_{2}$ with $g_{i j}=0$ for $j \neq 1$ we obtain positive definiteness of the operator

$$
\left[\begin{array}{cc}
M_{00} & M_{10}^{*} \\
M_{10} & T_{11}
\end{array}\right]
$$

on $\left(B^{k}\right)_{\Lambda}^{n_{0} \times n_{1}} \oplus\left(B^{k}\right)_{\Lambda \cap S}^{n_{1} \times n_{1}}$. Here

$$
M_{00}\left(g_{0}\right)=k_{00} g_{0}, M_{10}\left(g_{0}\right)=P_{S}\left(k_{10} g_{0}\right), g_{0} \in\left(B^{k}\right)_{\Lambda}^{n_{0} \times n_{1}}
$$

and

$$
T_{11}\left(g_{1}\right)=P_{S}\left(k_{11} g_{1}\right), g_{1} \in\left(B^{k}\right)_{\Lambda \cap S}^{n_{1} \times n_{1}}
$$

We had already established that $k_{00}>0$ and thus $M_{00}>0$. Consequently, we get that $S_{11}:=T_{11}-M_{10} M_{00}^{-1} M_{10}^{*}>0$. Observe that $S_{11}$ is the operator

$$
S_{11}\left(g_{1}\right)=P_{S}\left(\left(k_{11}-k_{10} k_{00}^{-1} k_{10}^{*}\right) g_{1}\right)
$$

and consequently (using again Section 5 in [26]) we obtain

$$
k_{11}-k_{10} k_{00}^{-1} k_{10}^{*}>0
$$

Performing a right canonical factorization

$$
k_{11}-k_{10} k_{00}^{-1} k_{10}^{*}=\left(e+s_{11,+}\right)^{*} s_{11, d}\left(e+s_{11,+}\right),
$$

it is straightforward to check that

$$
x_{11}=\left(e+s_{11,+}\right)^{-1} s_{11, d}^{-1}, \quad \text { and } \quad x_{01}=-k_{00}^{-1} k_{01} x_{11} .
$$

(Note: $k_{10}^{*}=k_{01}$.) Thus $x_{01} \in\left(A P W^{k}\right)_{\Lambda}^{n_{0} \times n_{1}}$ and $x_{11}^{ \pm 1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{1} \times n_{1}}$. Repeating this argument we obtain that $x_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times n_{j}}$ for $0 \leq i<j \leq q-1$, and $x_{i i}^{ \pm 1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{i} \times n_{i}}$ for $i=0, \ldots, q-1$.

Let us now look at the $q^{\text {th }}$ column of $x$. For this, observe that by applying $T$ to elements $g=\left(g_{i j}\right)_{i, j=0}^{p} \in \mathcal{B}_{2}$ with $g_{i j}=0$ for $j \neq q$ we get that

$$
\left[\begin{array}{ccc}
\widetilde{T}_{00} & A & H \\
A^{*} & M_{2} & B^{*} \\
H^{*} & B & T_{q q}
\end{array}\right]>0
$$

where

$$
\begin{aligned}
& \widetilde{T}_{00}: \quad\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}} \rightarrow\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}}, \quad \widetilde{T}_{00}\left(g_{0 q}\right)=P_{-S}\left(k_{00} g_{0 q}\right), \\
& A: \quad\left(B^{k}\right)_{\Lambda}^{Q \times n_{q}} \rightarrow\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}}, \quad A\left[\begin{array}{c}
g_{1 q} \\
\vdots \\
g_{q-1, q}
\end{array}\right]=P_{-S}\left(\sum_{r=1}^{q-1} k_{0 r} g_{r q}\right), \\
& H: \quad\left(B^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}} \rightarrow\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}}, \quad H g_{q q}=P_{-S}\left(k_{0 q} g_{q q}\right), \\
& M_{2}: \quad\left(B^{k}\right)_{\Lambda}^{Q \times Q} \rightarrow\left(B^{k}\right)_{\Lambda}^{Q \times Q}, \quad M_{2}\left[\begin{array}{c}
g_{1 q} \\
\vdots \\
g_{q-1, q}
\end{array}\right]=\left[\left(k_{i j}\right)_{i, j=1}^{q-1}\right]\left[\begin{array}{c}
g_{1 q} \\
\vdots \\
g_{q-1, q}
\end{array}\right], \\
& B: \quad\left(B^{k}\right)_{\Lambda}^{Q \times n_{q}} \rightarrow\left(B^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}}, \quad B\left[\begin{array}{c}
g_{1 q} \\
\vdots \\
g_{q-1, q}
\end{array}\right]=P_{S}\left(\sum_{r=1}^{q-1} k_{q r} g_{r q}\right), \\
& T_{q q}: \quad\left(B^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}} \rightarrow\left(B^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}}, \quad T_{q q} g_{q q}=P_{S}\left(k_{q q} g_{q q}\right),
\end{aligned}
$$

and where $Q=n_{1}+\cdots+n_{q-1}$. Thus $M_{2}>0$ and

$$
\left[\begin{array}{cc}
\widetilde{T}_{00}-A M_{2}^{-1} A^{*} & H-A M_{2}^{-1} B^{*}  \tag{2.19}\\
H^{*}-B M_{2}^{-1} A^{*} & T_{q q}-B M_{2}^{-1} B^{*}
\end{array}\right]>0
$$

Notice that $\widetilde{T}_{00}-A M_{2}^{-1} A^{*}$ is the Toeplitz operator

$$
g_{0 q} \mapsto P_{-S}\left(\left(k_{00}-\left(k_{0 i}\right)_{i=0}^{q-1}\left[\left(k_{i j}\right)_{i, j=1}^{q-1}\right]^{-1}\left(k_{i 0}\right)_{i=0}^{q-1}\right) g_{0 q}\right),
$$

and thus its symbol

$$
h_{1}=k_{00}-\left(k_{0 i}\right)_{i=0}^{q-1}\left[\left(k_{i j}\right)_{i, j=1}^{q-1}\right]^{-1}\left(k_{i 0}\right)_{i=0}^{q-1}
$$

is positive definite: $h_{1}>0$. Likewise

$$
h_{2}=k_{q q}-\left(k_{q i}\right)_{i=1}^{q-1}\left[\left(k_{i j}\right)_{i, j=1}^{q-1}\right]\left(k_{i q}\right)_{i=1}^{q-1}>0
$$

Write now the canonical factorizations (which exist by Theorem 2.1)

$$
h_{1}=\left(e+h_{1,-}\right)^{*} h_{1, d}\left(e+h_{1,-}\right), \quad h_{2}=\left(e+h_{2,+}\right)^{*} h_{2, d}\left(e+h_{2,+}\right)
$$

where

$$
h_{1,-},\left(e+h_{1,-}\right)^{-1}-e \in\left(A P W^{k}\right)_{\Lambda \cap(-S \backslash\{0\})}^{n_{0} \times n_{0}}, \quad h_{1, d} \in \mathbb{C}^{n_{0} \times n_{0}}, \quad h_{1, d}>0,
$$

and

$$
h_{2,+},\left(e+h_{2,+}\right)^{-1}-e \in\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{q} \times n_{q}}, \quad h_{2, d} \in \mathbb{C}^{n_{1} \times n_{1}}, \quad h_{2, d}>0 .
$$

Introducing the operators

$$
\begin{equation*}
H_{1}:\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}} \rightarrow\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}}, \quad H_{2}:\left(B^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}} \rightarrow\left(B^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}} \tag{2.20}
\end{equation*}
$$

by

$$
H_{1} g=h_{1, d}^{\frac{1}{2}}\left(e+h_{1,-}\right) g, \quad H_{2} g=h_{2, d}^{\frac{1}{2}}\left(e+h_{2,+}\right) g
$$

we obtain the equalities

$$
\widetilde{T}_{00}-A M_{2}^{-1} A^{*}=H_{1}^{*} H_{1}, \quad T_{q q}-B M_{2}^{-1} B^{*}=H_{2}^{*} H_{2}
$$

Observe that

$$
H_{3}:=H_{1}^{*-1}\left(H-A M_{2}^{-1} B^{*}\right) H_{2}^{-1}:\left(B^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}} \rightarrow\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}}
$$

is the operator defined by

$$
H_{3} g=P_{-S}\left(\left(e+h_{1,-}\right)^{*-1} h_{1, d}^{-\frac{1}{2}} \phi\left(e+h_{2,+}\right)^{-1} h_{2, d}^{-\frac{1}{2}} g\right), g \in\left(B^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}},
$$

where

$$
\phi=k_{0 q}-\left(k_{0 i}\right)_{i=1}^{q-1}\left[\left(k_{i j}\right)_{i, j=1}^{q-1}\right]^{-1}\left(k_{i q}\right)_{i=0}^{q-1}
$$

By (2.19) we get that $\left\|H_{3}\right\|<1$. Applying now Theorem 2.5, we get that

$$
\begin{aligned}
\beta & :=H_{3}\left(I-H_{3}^{*} H_{3}\right)^{-1}(e) \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}}, \\
\delta & :=\left(I-H_{3}^{*} H_{3}\right)^{-1}(e) \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}},
\end{aligned}
$$

$\delta^{-1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}}$ and $M\{\delta\}>0$. It is now straightforward to check the equalities

$$
\begin{equation*}
x_{0 q}=\left(e+h_{1,-}\right)^{-1} h_{1, d}^{-\frac{1}{2}} \beta h_{2, d}^{-\frac{1}{2}}, \quad x_{q q}=\left(e+h_{2,+}\right)^{-1} h_{2, d}^{-\frac{1}{2}} \delta h_{2, d}^{-\frac{1}{2}}, \tag{2.21}
\end{equation*}
$$

and

$$
\begin{align*}
\left(x_{i q}\right)_{i=1}^{q-1} & =-M_{2}^{-1}\left(A^{*} x_{0 q}+B^{*} x_{q q}\right)  \tag{2.22}\\
& =-\left[\left(k_{i j}\right)_{i, j=1}^{q-1}\right]^{-1}\left(\left[\left(k_{0 i}\right)_{i=1}^{q-1}\right]^{*} x_{0 q}+\left[\left(k_{q i}\right)_{i=1}^{q-1}\right]^{*} x_{q q}\right) .
\end{align*}
$$

Indeed, to verify (2.21) and (2.22), notice that

$$
\left[\begin{array}{ccc}
\widetilde{T}_{00} & A & H  \tag{2.23}\\
A^{*} & M_{2} & B^{*} \\
H^{*} & B & T_{q q}
\end{array}\right]\left[\begin{array}{c}
x_{0 q} \\
x_{1 q} \\
\vdots \\
x_{q q}
\end{array}\right]=\left[\begin{array}{c}
0 \\
0 \\
\vdots \\
e
\end{array}\right]
$$

and substitute formulas (2.21) and (2.22) for $x_{j q}$ in (2.23). Now clearly,

$$
x_{0 q} \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{0} \times n_{q}}, x_{q q}^{ \pm 1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{q} \times n_{q}}
$$

$$
M\left\{x_{q q}\right\}>0,\left(x_{i q}\right)_{i=1}^{q-1} \in\left(A P W^{k}\right)_{\Lambda}^{Q \times n_{q}} .
$$

This yields that the $q^{t h}$ column of $x$ is of the required form. For columns $q+1$, $\ldots, p$ of $x$ one uses the same reasoning with $\left(k_{i j}\right)_{i, j=0}^{q}$ replaced by $\left(k_{i j}\right)_{i, j=s}^{q+s}, s=$ $1, \ldots, p-q$, respectively. This concludes the proof.

Analogously to the proof of Proposition 2.6, one may prove the following.
Proposition 2.7. Let $\widetilde{T}: \mathcal{B}_{3} \rightarrow \mathcal{B}_{3}$ be as in Theorem 2.3 and suppose $\widetilde{T}>0$. Then $y:=\widetilde{T}^{-1}(e)$ is an invertible element of $\mathcal{M}_{-}$.

Now we may prove our main result.
Proof of Theorem 2.3. First assume (i), i.e., $k$ has a positive extension $f$. Then the multiplication operator $g \mapsto f g$ is a positive definite operator on $\mathcal{B}$. But then its restrictions to $\mathcal{B}_{2}$ and $\mathcal{B}_{3}$ are also positive definite. Using the multiplication table (2.10), one obtains easily that for $g \in \mathcal{M}_{2}$ we have that $T(g)=P_{2}(k g)=P_{2}(f g)$. But then the same holds for $g \in \mathcal{B}_{2}$. Consequently, it follows that $T>0$. Likewise, one shows that $\widetilde{T}>0$.

Assume now that (ii) holds. By Proposition 2.6 we have that $x:=T^{-1}(e) \in \mathcal{M}_{2}$, $f x \in e+\mathcal{M}_{4}+\mathcal{M}_{3}^{0}+\mathcal{M}_{1}$, and $x^{-1} \in \mathcal{M}_{+}$. In addition, it is easy to see that $P_{d} x>0$. By Theorem 1.1 in Chapter XXXIV of [10] we obtain that $f$ has a positive extension $f_{0}$ (in fact, a so-called "band extension"), which is given by

$$
f_{0}=x^{*-1}\left(P_{d} x\right) x^{-1}
$$

But then (i) holds, and consequently (iii) holds as well. It follows by Proposition 2.7 and Theorems 1.2 and 1.3 in Chapter XXXIV of [10] that $f_{0}$ may also be found via

$$
f_{0}=y^{*-1}\left(P_{d} y\right) y^{-1}
$$

This shows the first part of the theorem.
Next, observe that in the terminology of Chapter XXXIV of [10], $\mathcal{M}$ is "an algebra with band structure (2.9) in the unital $C^{*}$-algebra $\mathcal{B}$ ". Furthermore, note that Axiom (A) in Chapter XXXIV of [10] is satisfied; in other words, if $F \in \mathcal{M}_{+}$ is such that sup $\|F(t)\|<1$, then $(I-F)^{-1} \in \mathcal{M}_{+}$. Indeed, by Proposition 2.3 $t \in \mathbb{R}^{k}$
in [26] (with $f=I-F$ and $\Psi(z)=z^{-1}$ ) we obtain $(I-F)^{-1} \in(A P W)^{N \times N}$; on the other hand, the series $(I-F)^{-1}=I+F+F^{2}+\cdots$ converges in $\|\cdot\|_{\infty}$, hence we have $(I-F)^{-1} \in \mathcal{M}_{+}$as required. Since Axiom (A) in Chapter XXXIV of [10] holds, we may apply Theorem 2.1 of Chapter XXXIV in [10], yielding the linear fractional description of the set of all positive extensions.

Next, we need to check that $\mathcal{M}$ with decomposition (2.9) satisfies Axioms (C1) and (C2) in Chapter XXXIV, Section 4 of [10]. Indeed, if $f \in \mathcal{M}$ is positive then $P_{d} f$ is a positive semi-definite matrix, and $P_{d} f$ equals zero if and only if $f=0$. Thus we may apply Theorems 1.3 and 4.2 in Chapter XXXIV in [10] to obtain the last statement in the theorem.

Proof of Proposition 2.4. Observe that if $f$ has the right spectral factorization (2.11), then

$$
\log (\operatorname{det} f)=\log \left(\operatorname{det}\left(e+f_{+}\right)^{*}\right)+\log \operatorname{det} f_{d}+\log \left(\operatorname{det}\left(e+f_{+}\right)\right)
$$

Notice: $\log \left(\operatorname{det}\left(e+f_{+}\right)\right) \in \mathcal{M}_{+}$(see Proposition 3.2 in [24]). Since $\left(e+f_{+}\right)^{ \pm 1} \in$ $e+\mathcal{M}_{+}^{0}$ we have that $M\left\{\log \left(\operatorname{det}\left(e+f_{+}\right)\right)\right\}=0$. Thus

$$
\Delta(f)=\log \operatorname{det} f_{d}
$$

But then it follows from (2.13) that

$$
\Delta(f)=\log \operatorname{det} f_{d} \leq \log \operatorname{det}\left(P_{d} x\right)^{-1}=\Delta\left(f_{0}\right)
$$

with equality if and only if $f=f_{0}$ (since for matrices $A \neq B$ with $A \geq B>0$ we have $\operatorname{det} A>\operatorname{det} B$; in the terminology of Section XXXIV. 4 in [10], the function "log det" is strictly $\mathcal{B}$-monotone). The reasoning may also be applied to the other inequality in (2.13).

It should be noted that similar results may be obtained when the $q^{t h}$ diagonal of $k$ does not solely consist of elements in $\left(A P W^{k}\right)_{\Lambda \cap(-S)}$ but is per block in any of

$$
\begin{gathered}
\left(A P W^{k}\right)_{\Lambda \cap(-S)}, \quad\left(A P W^{k}\right)_{\Lambda \cap(-S \backslash\{0\})}, \\
\left(A P W^{k}\right)_{\Lambda}, \quad\{0\},\left(A P W^{k}\right)_{\Lambda \cap(\nu-S)},
\end{gathered}
$$

or

$$
\left(A P W^{k}\right)_{\Lambda \cap(\nu-(S \backslash\{0\}))},
$$

where $\nu \in \Lambda$. Each time a mixture of choices is made, the spaces $\mathcal{M}_{1}, \mathcal{M}_{2}^{0}$, $\mathcal{M}_{d}, \mathcal{M}_{3}^{0}, \mathcal{M}_{4}$ should be appropriately defined and will result in different types of Toeplitz operators, each time yielding an analogue of Theorem 2.3. The proofs of such variations require the same type of simple modifications as the ones outlined in [23, Section 10] and [25, Section 4.5]. We omit further details.

## 3. The contractive st-block problem

Fix $s, t \in \mathbb{N}, \max \{t-s, 0\} \leq q \leq t+1$, a halfspace $S$ of $\mathbb{R}^{k}$ and a subgroup $\Lambda$ of $\mathbb{R}^{k}$. Our problem is the following. Let $f=\left(f_{i j}\right)_{i=1, j=1}^{s, t}$ be given so that $f_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times m_{j}}, j-i<q$, and $f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{i} \times m_{j}}, j-i=q$. Find, if possible,

$$
\widetilde{f}=\left(\widetilde{f}_{i j}\right)_{i=1, j=1}^{s, t} \in\left(A P W^{k}\right)_{\Lambda}^{\left(n_{1}+\cdots+n_{s}\right) \times\left(m_{1}+\cdots+m_{t}\right)}
$$

so that

$$
\begin{gathered}
\tilde{f}_{i j}=f_{i j} \quad \text { for } \quad j-i<q \\
\widetilde{f}_{i j}-f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{i} \times m_{j}} \text { for } j-i=q
\end{gathered}
$$

and

$$
\|\widetilde{f}\|_{\infty}<1
$$

Such an $\tilde{f}$ is called a strictly contractive extension of $f$.
In order to state the results we introduce the following spaces. Let

$$
\begin{aligned}
& \mathcal{A}_{1}=\left\{f=\left(f_{i j}\right)_{i, j=1}^{s}: f_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times n_{j}} \quad \text { for } \quad i<j\right. \\
& \left.\qquad f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{i} \times n_{j}} \quad \text { for } \quad i=j ; \quad f_{i j}=0 \text { for } i>j\right\},
\end{aligned}
$$

$$
\begin{aligned}
& \mathcal{A}_{2}=\left\{f=\left(f_{i j}\right)_{i=1, j=1}^{s, t}: f_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times m_{j}} \text { for } j-i<q ;\right. \\
& \left.f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{i} \times m_{j}} \text { for } j-i=q ; \quad f_{i j}=0 \quad \text { for } j-i>q\right\}, \\
& \mathcal{A}_{3}=\left\{f=\left(f_{i j}\right)_{i, j=1}^{t}: f_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{m_{i} \times m_{j}} \quad \text { for } i<j ;\right. \\
& \left.f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{m_{i} \times m_{j}} \quad \text { for } i=j ; \quad f_{i j}=0 \text { for } i>j\right\}, \\
& \mathcal{A}_{4}=\left\{f=\left(f_{i j}\right)_{i=1, j=1}^{s, t}: f_{i j}=0 \text { for } j-i<q ;\right. \\
& \left.f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{i} \times m_{j}} \text { for } j-i=q ; \quad f_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times m_{j}} \text { for } j-i>q\right\} .
\end{aligned}
$$

Further, let

$$
\left(\mathcal{A}_{1}^{*}\right)_{d}=\mathbb{C}^{n_{1} \times n_{1}} \oplus \cdots \oplus \mathbb{C}^{n_{s} \times n_{s}}, \quad\left(\mathcal{A}_{3}\right)_{d}=\mathbb{C}^{m_{1} \times m_{1}} \oplus \cdots \oplus \mathbb{C}^{m_{t} \times m_{t}}
$$

We let $\mathcal{C}_{1}, \mathcal{C}_{2}$, and $\mathcal{C}_{3}$ denote the closures of $\mathcal{A}_{1}, \mathcal{A}_{2}$, and $\mathcal{A}_{3}$ in $\left(B^{k}\right)_{\Lambda}^{\left(\Sigma n_{i}\right) \times\left(\Sigma n_{i}\right)}$, $\left(B^{k}\right)_{\Lambda}^{\left(\Sigma n_{i}\right) \times\left(\Sigma m_{i}\right)}$, and $\left(B^{k}\right)_{\Lambda}^{\left(\Sigma m_{i}\right) \times\left(\Sigma m_{i}\right)}$, respectively. Also, $\mathcal{C}_{1}^{*}$ and $\mathcal{C}_{2}^{*}$ stand for the closures of $\mathcal{A}_{1}^{*}$ and $\mathcal{A}_{2}^{*}$ in $\left(B^{k}\right)_{\Lambda}^{\left(\Sigma n_{i}\right) \times\left(\Sigma n_{i}\right)}$ and $\left(B^{k}\right)_{\Lambda}^{\left(\Sigma m_{i}\right) \times\left(\Sigma n_{i}\right)}$, respectively. For a closed subspace $\mathcal{C}$ of $\left(B^{k}\right)^{N \times M}$ the symbol $P_{\mathcal{C}}$ stands for the orthogonal projection of $\left(B^{k}\right)^{N \times M}$ onto $\mathcal{C}$.

Theorem 3.1. Let $f \in \mathcal{A}_{2}$ be given. Then the following statements are equivalent:
(i) $f$ has a strictly contractive extension.
(ii) The operator $\underset{\sim}{H}: \mathcal{C}_{3} \rightarrow \mathcal{C}_{2}$ defined by $\underset{\sim}{H}(g)=P_{\mathcal{C}_{2}}(f g)$, is a strict contraction.
(iii) The operator $\widetilde{H}: \mathcal{C}_{2}^{*} \rightarrow \mathcal{C}_{1}^{*}$ defined by $\widetilde{H}(g)=P_{\mathcal{C}_{1}^{*}}(f g)$, is a strict contraction.

In that case, let

$$
\widehat{\alpha}=\left(I-\widetilde{H} \widetilde{H}^{*}\right)^{-1}(e), \quad \widehat{\delta}=\left(I-H^{*} H\right)^{-1}(e)
$$

Then $\widehat{\alpha}^{-1} \in \mathcal{A}_{1}^{*}, \widehat{\delta}^{-1} \in \mathcal{A}_{3}$, and $P_{\left(\mathcal{A}_{1}^{*}\right)_{d}}(\widehat{\alpha})>0$ and $P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta})>0$. Further, put

$$
\alpha=\widehat{\alpha}\left[P_{\left(\mathcal{A}_{1}^{*}\right)_{d}}(\widehat{\alpha})\right]^{-\frac{1}{2}}, \quad \delta=\widehat{\delta}\left[P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta})\right]^{-\frac{1}{2}}
$$

$\beta=H(\delta)$ and $\gamma=\widetilde{H}^{*}(\alpha)$. Then

$$
h_{0}:=\beta \delta^{-1}=\alpha^{*-1} \gamma^{*}
$$

is a strictly contractive extension of $f$. In fact, $h_{0}$ is the unique strictly contractive extension with $h_{0}\left(e-h_{0}^{*} h_{0}\right)^{-1} \in \mathcal{A}_{2}$. Moreover, $h$ is a strictly contractive extension of $f$ if and only if

$$
h=(\alpha g+\beta)(\gamma g+\delta)^{-1}
$$

for some (unique!) $g \in \mathcal{A}_{4}$ with $\|g\|_{\infty}<1$. Also, $h$ is a strictly contractive extension of $f$ if and only if

$$
h=\left(\alpha^{*}+g \beta^{*}\right)^{-1}\left(\gamma^{*}+g \delta^{*}\right),
$$

for some (unique!) $g \in \mathcal{A}_{4}^{*}$ with $\|g\|_{\infty}<1$.

In addition, let $h$ be a strictly contractive extension of $f$, and perform factorizations

$$
\begin{equation*}
I-h^{*} h=\left(e+h_{+}\right)^{*} h_{1, d}\left(e+h_{+}\right), \tag{3.1}
\end{equation*}
$$

and

$$
\begin{equation*}
I-h h^{*}=\left(e+h_{-}\right)^{*} h_{2, d}\left(e+h_{-}\right), \tag{3.2}
\end{equation*}
$$

with $h_{+},\left(e+h_{+}\right)^{-1} \in \mathcal{A}_{3} ; P_{\left(\mathcal{A}_{3}\right)_{d}}\left(h_{+}\right)=0 ; h_{1, d} \in\left(\mathcal{A}_{3}\right)_{d} ; h_{-},\left(e+h_{-}\right)^{-1} \in \mathcal{A}_{1}^{*}$; $P_{\left(\mathcal{A}_{1}^{*}\right)_{d}}\left(h_{-}\right)=0$, and $h_{2, d} \in\left(\mathcal{A}_{1}\right)_{d}$. Then

$$
h_{1, d} \leq P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta}), h_{2, d} \leq P_{\left(\mathcal{A}_{1}\right)_{d}}(\widehat{\alpha})
$$

and equality occurs in one of the inequalities if and only if $h=h_{0}$ (and thus both inequalities are equalities).

Finally,

$$
\begin{equation*}
\left\|h_{0}\right\|_{B^{k}} \leq \frac{\|f\|_{B^{k}}}{\sqrt{1-\|H\|^{2}}}=\frac{\|f\|_{B^{k}}}{\sqrt{1-\|\tilde{H}\|^{2}}} \tag{3.3}
\end{equation*}
$$

where

$$
\begin{equation*}
\|g\|_{B^{k}}=\left[\operatorname{trace} M\left\{g^{*} g\right\}\right]^{1 / 2}, \quad g \in\left(B^{k}\right)^{\left(\sum n_{j}\right) \times\left(\sum m_{j}\right)} \tag{3.4}
\end{equation*}
$$

is the Besikovitch norm.
Proof. It is straightforward that (i) $\Rightarrow$ (ii) and (i) $\Rightarrow$ (iii) hold. Indeed, if $f_{\text {ext }}$ is a strictly contractive extension of $f$, then the map

$$
g \mapsto f_{e x t} g
$$

is a strictly contractive operator acting $\left(B^{k}\right)_{\Lambda}^{\left(\Sigma m_{i}\right) \times \ell} \rightarrow\left(B^{k}\right)_{\Lambda}^{\left(\Sigma n_{i}\right) \times \ell}$, where $\ell \in \mathbb{N}$ (in particular, we shall use it for the cases when $\ell=\Sigma m_{i}$ or $\ell=\Sigma n_{i}$ ). But then so are any restrictions of this multiplication operator. In particular, it follows that $\|H\|<1$ and $\|\widetilde{H}\|<1$.

Let us show (ii) $\Rightarrow$ (i). So assume that $\|H\|<1$. Let $Q=n_{1}+\ldots+n_{s}$ and $R=m_{1}+\ldots+m_{t}$. Put

$$
\begin{equation*}
\mathcal{M}=\left(A P W^{k}\right)_{\Lambda}^{(Q+R) \times(Q+R)}=\mathcal{M}_{1} \dot{+} \mathcal{M}_{2}^{0} \dot{+} \mathcal{M}_{d} \dot{+} \mathcal{M}_{3}^{0} \dot{+} \mathcal{M}_{4}, \tag{3.5}
\end{equation*}
$$

where

$$
\begin{aligned}
& \mathcal{M}_{1}=\left[\begin{array}{cc}
0 & \mathcal{A}_{4} \\
0 & 0
\end{array}\right], \quad \mathcal{M}_{2}^{0}=\left[\begin{array}{cc}
\mathcal{A}_{1}^{0} & \mathcal{A}_{2} \\
0 & \mathcal{A}_{3}^{0}
\end{array}\right], \\
& \mathcal{M}_{d}=\left[\begin{array}{cc}
\left(\mathcal{A}_{1}^{*}\right)_{d} & 0 \\
0 & \left(\mathcal{A}_{3}\right)_{d}
\end{array}\right], \quad \mathcal{M}_{3}^{0}=\left(\mathcal{M}_{2}^{0}\right)^{*}, \quad \mathcal{M}_{4}=\left(\mathcal{M}_{1}\right)^{*},
\end{aligned}
$$

with

$$
\mathcal{A}_{r}^{0}=\left\{f=\left(f_{j k}\right)_{j, k} \in \mathcal{A}_{r}: M\left\{f_{i i}\right\}=0 \text { for all } i\right\}, \quad r=1,3 .
$$

Furthermore, let

$$
k=\left[\begin{array}{cc}
e & f \\
f^{*} & e
\end{array}\right] .
$$

Now we are in the setting of Theorem 2.3. Moreover, with $k$ as defined above, we find that the operator $T$ in Theorem 2.3(ii) corresponds to

$$
T=\left[\begin{array}{cc}
I & H \\
H^{*} & I
\end{array}\right]
$$

Thus we apply Theorem 2.3(ii) and obtain that $k$ has a positive extension

$$
k_{e x t}=\left[\begin{array}{cc}
e & f_{e x t} \\
f_{e x t}^{*} & e
\end{array}\right]
$$

where $f_{\text {ext }}-f \in \mathcal{A}_{4}$. But then it follows that $f$ has a strictly contractive extension $f_{\text {ext }}$. This shows (ii) $\Rightarrow$ (i).

In the same manner one may show (iii) $\Rightarrow$ (i) using Theorem 2.3.
In particular, the equivalence of (ii) and (iii) shows that $\|H\|=\|\widetilde{H}\|$.
Next we observe that with $T$ and $\widetilde{T}$ as in Theorem 2.3 with the above choice of $k$, we have that

$$
T\left[\begin{array}{cc}
e & -\widehat{\beta} \\
0 & \widehat{\delta}
\end{array}\right]=\left[\begin{array}{ll}
e & 0 \\
0 & e
\end{array}\right], \quad \widetilde{T}\left[\begin{array}{cc}
\widehat{\alpha} & 0 \\
-\widehat{\gamma} & e
\end{array}\right]=\left[\begin{array}{ll}
e & 0 \\
0 & e
\end{array}\right]
$$

where $\widehat{\beta}=H(\widehat{\delta})$ and $\widehat{\gamma}=\widetilde{H}^{*}(\widehat{\alpha})$. Thus when we apply Theorem 2.3 with the above choice of $k$, we find that $x$ and $y$ defined in Theorem 2.3 correspond to

$$
x=\left[\begin{array}{cc}
e & -\widehat{\beta} \\
0 & \widehat{\delta}
\end{array}\right], \quad y=\left[\begin{array}{cc}
\widehat{\alpha} & 0 \\
-\widehat{\gamma} & e
\end{array}\right]
$$

Thus we obtain from the properties of $x$ and $y$ that $\widehat{\alpha}^{-1} \in \mathcal{A}_{1}^{*}, \widehat{\delta}^{-1} \in \mathcal{A}_{3}$, and $P_{\left(\mathcal{A}_{1}^{*}\right)_{d}}(\widehat{\alpha})>0$ and $P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta})>0$. Next, when we introduce $u$ and $v$ as in Theorem 2.3, and let

$$
g=\left[\begin{array}{cc}
0 & \tilde{g} \\
0 & 0
\end{array}\right] \in \mathcal{M}_{1}
$$

then we get that

$$
(u+v g)^{*-1}\left(e-g^{*} g\right)(u+v g)^{-1}=\left[\begin{array}{cc}
e & (\alpha \widetilde{g}+\beta)(\gamma \widetilde{g}+\delta)^{-1} \\
{\left[(\alpha \widetilde{g}+\beta)(\gamma \widetilde{g}+\delta)^{-1}\right]^{*}} & e
\end{array}\right]
$$

This yields the first linear fractional description of the set of all solutions. Likewise, using the linear fractional description $(v+u h)^{*-1}\left(e-h^{*} h\right)(v+u h)^{-1}$ of Theorem 2.3 yields the second linear fractional description.

Note also that it follows from Theorem 2.3 that

$$
f_{0}^{-1}=\left[\begin{array}{cc}
I & h_{0} \\
h_{0}^{*} & I
\end{array}\right]^{-1}=\left[\begin{array}{cc}
* & h_{0}\left(I-h_{0}^{*} h_{0}\right)^{-1} \\
* & *
\end{array}\right] \in \mathcal{M}_{2}^{0}+\mathcal{M}_{d} \dot{+} \mathcal{M}_{3}^{0} .
$$

Thus $h_{0}\left(I-h_{0}^{*} h_{0}\right)^{-1} \in \mathcal{A}_{2}$. In addition, the statement regarding the inequalities

$$
h_{1, d} \leq P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta}), h_{2, d} \leq P_{\left(\mathcal{A}_{1}\right)_{d}}(\widehat{\alpha})
$$

follows directly by applying (2.13).
Finally for the last statement, we use a similar type of argument as in the proof of Theorem 7.1 in [24]. Indeed, without loss of generality we may assume that $\|H\| \neq$ 0 , excluding the trivial case $f=0$. Let $\epsilon=\frac{1}{\|H\|}$. Let now $\widetilde{f}(t)=\beta(t) \delta(t)^{-1}-f(t)$. Using the easily derived inequality

$$
\operatorname{trace} M^{*} M \leq-\log \operatorname{det}\left(I-M^{*} M\right)
$$

which holds for every $M \in \mathbb{C}^{Q \times R}$ with $\|M\|<1$, we have:

$$
\begin{aligned}
\|f+\widetilde{f}\|_{B^{k}}^{2} & =\lim _{T \rightarrow \infty} \frac{1}{(2 T)^{k}} \int_{[-T, T]^{k}} \operatorname{trace}\left[((f+\widetilde{f})(t))^{*}(f+\widetilde{f})(t)\right] d t \\
& \leq-\lim _{T \rightarrow \infty} \frac{1}{(2 T)^{k}} \int_{[-T, T]^{k}} \log \operatorname{det}\left(I-((f+\widetilde{f})(t))^{*}(f+\widetilde{f})(t)\right) d t \\
& =-\lim _{T \rightarrow \infty} \frac{1}{(2 T)^{k}} \int_{[-T, T]^{k}} \log \operatorname{det}\left(I-\left(\beta(t) \delta(t)^{-1}\right)^{*} \beta(t) \delta(t)^{-1}\right) d t
\end{aligned}
$$

Since $\left(\delta(t) P_{\left(\mathcal{A}_{3}\right)_{d}}(\delta)^{-1}\right)^{ \pm 1} \in e+\mathcal{A}_{3}^{0}$ we get the equality

$$
M\left\{\log \operatorname{det}\left(\delta(t) P_{\left(\mathcal{A}_{3}\right)_{d}}(\delta)^{-1}\right)\right\}=0
$$

Also
$I-\left(\beta(t) \delta(t)^{-1}\right)^{*} \beta(t) \delta(t)^{-1}=\delta(t)^{-1 *}\left(\delta(t)^{*} \delta(t)-\beta(t)^{*} \beta(t)\right) \delta(t)^{-1}=\delta(t)^{-1 *} \delta(t)^{-1}$.
Using the last two equalities, we obtain

$$
\begin{equation*}
\|f+\widetilde{f}\|_{B^{k}}^{2} \leq M\left\{\log \operatorname{det}\left(\delta(t) \delta(t)^{*}\right)\right\}=\log \operatorname{det}\left(P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta})\right)=\operatorname{trace}\left(\log P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta})\right) \tag{3.6}
\end{equation*}
$$

Note that
$\left.P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta})=P_{\left(\mathcal{A}_{3}\right)_{d}}\left(\left(I-H^{*} H\right)^{-1}\left(I_{\sum m_{i}}\right)\right)=P_{\left(\mathcal{A}_{3}\right)_{d}}\left(I+H^{*}\left(I-H H^{*}\right)^{-1} H\right)\left(I_{\sum m_{i}}\right)\right)$.
From the inequality $\log (1+r) \leq r$ valid for $r \geq 0$ we get that

$$
\begin{equation*}
\log \operatorname{det} P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta}) \leq \operatorname{trace}\left[P_{\left(\mathcal{A}_{3}\right)_{d}}\left(H^{*}\left(I-H H^{*}\right)^{-1} H\right)\right] \tag{3.7}
\end{equation*}
$$

Since $\|H\|=\epsilon$, it follows from (3.6) and (3.7) that

$$
\|f+\widetilde{f}\|_{B^{k}}^{2} \leq \frac{1}{1-\epsilon^{-2}} \operatorname{trace} M\left\{\left(H^{*} H\right)\left(I_{\sum m_{i}}\right)\right\}=\frac{\epsilon^{2}}{\epsilon^{2}-1}\|f\|_{B^{k}}^{2}
$$

This proves the inequality in (3.3). The equality there follows because $\|H\|=$ $\|\widetilde{H}\|$.

Also in the contractive setting we may state an entropy result along the lines of Proposition 2.4. Its proof is based on the same observations and we will omit it.

Proposition 3.2. Let $f \in \mathcal{A}_{2}$ be given and suppose that $f$ has a strictly contractive extension. Let $h_{0}$ be defined as in Theorem 3.1. If $h$ is a strictly contractive extension of $f$ then

$$
\Delta\left(I-h^{*} h\right) \leq \Delta\left(I-h_{0}^{*} h_{0}\right)
$$

and equality holds if and only if $h=h_{0}$.
Analogously to the observation made at the end of Section 2, the same methods apply for other situations in which the elements of $f$ in the $q^{t h}$ diagonal, per block, are taken from one of the sets

$$
\begin{gathered}
\left(A P W^{k}\right)_{\Lambda \cap(-S)}, \quad\left(A P W^{k}\right)_{\Lambda \cap(-S \backslash\{0\})}, \\
\left(A P W^{k}\right)_{\Lambda},\{0\},\left(A P W^{k}\right)_{\Lambda \cap(-S+\nu)}, \\
\left(A P W^{k}\right)_{\Lambda \cap((-S \backslash\{0\})+\nu)},
\end{gathered}
$$

where $\nu \in \Lambda$. For each selection of one of the above sets for every block on the $q$-th diagonal, a theorem results which is analogous to Theorem 3.1. The corresponding analogs of Proposition 3.2 hold true as well.

We state as an example the variation where the elements of $f$ in the $q$ th diagonal are all taken from $\left(A P W^{k}\right)_{\Lambda \cap(-S \backslash\{0\})}$. This variation will be used in the next section.

Introduce

$$
\begin{aligned}
& \widehat{\mathcal{A}}_{2}=\left\{f=\left(f_{i j}\right)_{i=1, j=1}^{s, t}: f_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times m_{j}} \text { for } j-i<q ;\right. \\
&\left.f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap(-S \backslash\{0\})}^{n_{i} \times m_{j}} \text { for } j-i=q ; \quad f_{i j}=0 \text { for } j-i>q\right\}, \\
& \widehat{\mathcal{A}}_{4}=\left\{f=\left(f_{i j}\right)_{i=1, j=1}^{s, t}: f_{i j}=0 \text { for } j-i<q ;\right. \\
&\left.f_{i j} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{i} \times m_{j}} \text { for } j-i=q ; \quad f_{i j} \in\left(A P W^{k}\right)_{\Lambda}^{n_{i} \times m_{j}} \text { for } j-i>q\right\},
\end{aligned}
$$

and let $\widehat{\mathcal{C}_{2}}$ and $\widehat{\mathcal{C}}_{2}^{*}$ denote the closures in $\left(B^{k}\right)_{\Lambda}^{\left(\Sigma n_{i}\right) \times\left(\Sigma m_{i}\right)}$ and $\left(B^{k}\right)_{\Lambda}^{\left(\Sigma m_{i}\right) \times\left(\Sigma n_{i}\right)}$ of $\widehat{\mathcal{A}}_{2}$ and $\widehat{\mathcal{A}}_{2}^{*}$, respectively. The spaces $\mathcal{A}_{3}, \mathcal{C}_{3}, \mathcal{A}_{1}^{*}$, and $\mathcal{C}_{1}^{*}$ are defined as before.

Theorem 3.3. Let $f \in \widehat{\mathcal{A}}_{2}$ be given. Then the following statements are equivalent:
(i) $f$ has a strictly contractive extension, i.e., there exists an

$$
\tilde{f} \in\left(A P W^{k}\right)_{\Lambda}^{\left(n_{1}+\cdots+n_{s}\right) \times\left(m_{1}+\cdots+m_{t}\right)}
$$

such that $P_{\widehat{\mathcal{A}}_{2}}(\widetilde{f})=f$ and $\|\widetilde{f}\|_{\infty}<1$.
(ii) The operator $H_{0}: \mathcal{C}_{3} \rightarrow \widehat{\mathcal{C}}_{2}$ defined by $H_{0}(g)=P_{\widehat{\mathcal{C}}_{2}}(f g)$, is a strict contraction.
(iii) The operator $\widetilde{H}_{0}: \widehat{\mathcal{C}}_{2}^{*} \rightarrow \mathcal{C}_{1}^{*}$ defined by $\widetilde{H}_{0}(g)=P_{\mathcal{C}_{1}^{*}}(f g)$, is a strict contraction.
In that case, let

$$
\widehat{\alpha}=\left(I-\widetilde{H}_{0} \widetilde{H}_{0}^{*}\right)^{-1}(e), \quad \widehat{\delta}=\left(I-H_{0}^{*} H_{0}\right)^{-1}(e)
$$

Then $\widehat{\alpha}^{-1} \in \mathcal{A}_{1}^{*}, \widehat{\delta}^{-1} \in \mathcal{A}_{3}$, and $P_{\left(\mathcal{A}_{1}^{*}\right)_{d}}(\widehat{\alpha})>0$ and $P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta})>0$. Further, put

$$
\alpha=\widehat{\alpha}\left[P_{\left(\mathcal{A}_{1}^{*}\right)_{d}}(\widehat{\alpha})\right]^{-\frac{1}{2}}, \quad \delta=\widehat{\delta}\left[P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta})\right]^{-\frac{1}{2}}
$$

$\beta=H_{0}(\delta)$ and $\gamma=\widetilde{H}_{0}^{*}(\alpha)$. Then

$$
h_{0}:=\beta \delta^{-1}=\alpha^{*-1} \gamma^{*}
$$

is a strictly contractive extension of $f$. In fact, $h_{0}$ is the unique strictly contractive extension with $h_{0}\left(e-h_{0}^{*} h_{0}\right)^{-1} \in \widehat{\mathcal{A}}_{2}$. Moreover, $h$ is a strictly contractive extension of $f$ if and only if

$$
h=(\alpha g+\beta)(\gamma g+\delta)^{-1}
$$

for some (unique!) $g \in \widehat{\mathcal{A}}_{4}$ with $\|g\|_{\infty}<1$. Also, $h$ is a strictly contractive extension of $f$ if and only if

$$
h=\left(\alpha^{*}+g \beta^{*}\right)^{-1}\left(\gamma^{*}+g \delta^{*}\right),
$$

for some (unique!) $g \in \widehat{\mathcal{A}}_{4}^{*}$ with $\|g\|_{\infty}<1$.

In addition, let $h$ be a strictly contractive extension, and perform factorizations (3.1) and (3.2), with $h_{+},\left(e+h_{+}\right)^{-1} \in \mathcal{A}_{3} ; P_{\left(\mathcal{A}_{3}\right)_{d}}\left(h_{+}\right)=0 ; h_{1, d} \in\left(\mathcal{A}_{3}\right)_{d} ; h_{-},(e+$ $\left.h_{-}\right)^{-1} \in \mathcal{A}_{1}^{*} ; P_{\left(\mathcal{A}_{1}^{*}\right)_{d}}\left(h_{-}\right)=0$, and $h_{2, d} \in\left(\mathcal{A}_{1}\right)_{d}$, as in Theorem 3.1. Then

$$
h_{1, d} \leq P_{\left(\mathcal{A}_{3}\right)_{d}}(\widehat{\delta}), h_{2, d} \leq P_{\left(\mathcal{A}_{1}\right)_{d}}(\widehat{\alpha})
$$

and equality occurs in one of the inequalities if and only if $h=h_{0}$ (and thus both inequalities are equalities). Finally,

$$
\left\|h_{0}\right\|_{B^{k}} \leq \frac{\|f\|_{B^{k}}}{\sqrt{1-\left\|H_{0}\right\|^{2}}}=\frac{\|f\|_{B^{k}}}{\sqrt{1-\left\|\widetilde{H}_{0}\right\|^{2}}}
$$

where $\|\cdot\|_{B^{k}}$ is given by (3.4).
The proof of Theorem 3.3 requires a modification of the proof of Theorem 3.1 (and thus implicitly, of the proof of Theorem 2.1). These modifications are of the same type as the ones in [23, Section 10], [24, Section 4], and [25, Section 4.5]. We omit further details.

## 4. Model matching

It is well-known that solutions of the standard four block (as well as one-block and two-block) problems lead, when using frequency domain approach, to results concerning model matching, a key problem in control systems. This approach to model matching has been extensively studied in the the engineering literature (see, for example, [4], [7], and references there), especially for one-dimensional systems, and see [27] for some recent results in this direction for multidimensional systems. In this section we provide an interpretation of Theorem 3.3 in the context of model matching. We consider filters acting on square summable sequences indexed by an additive group in $\mathbb{R}^{k}$. The case of the group $\mathbb{Z}$ in $\mathbb{R}$ is the familiar case, treated extensively in the literature (see, e. g., [21]).

Let $\Lambda$ be an additive subgroup of $\mathbb{R}^{k}$. For $\Delta \subseteq \Lambda$ we let $\ell_{2}^{N}(\Delta)$ denote the Hilbert space of sequences $\left(v_{\lambda}\right)_{\lambda \in \Delta}$ where at most countably many $v_{\lambda} \in \mathbb{C}^{N}$ are nonzero and which are square summable in norm, i.e., $\sum_{\lambda \in \Delta}\left\|v_{\lambda}\right\|^{2}<\infty$. By $\ell_{1}^{N \times M}(\Delta)$ we denote the Banach space of sequences $\left(f_{\lambda}\right)_{\lambda \in \Delta}$ where at most countably many $f_{\lambda} \in \mathbb{C}^{N \times M}$ are nonzero and which are summable in norm, i.e., $\sum_{\lambda \in \Delta}\left\|f_{\lambda}\right\|<\infty$.

Fix a halfspace $S$ of $\mathbb{R}^{k}$. With $S$ we associate an ordering $\leq_{S}$ on $\Lambda$ by $q \leq_{S} p$ if and only if $p-q \in S$. We shall use the interval notation with the usual conventions. So, for instance, $S \cap \Lambda=[0, \infty)$. With an element $f \in \ell_{1}^{N \times M}([0, \infty))$, we associate a filter $\Sigma_{f}: \ell_{2}^{M}([0, \infty)) \rightarrow \ell_{2}^{N}([0, \infty))$, defined by

$$
\Sigma_{f}\left(\left(u_{\lambda}\right)_{\lambda \in[0, \infty)}\right)=\left(y_{\lambda}\right)_{\lambda \in[0, \infty)}, y_{\lambda}=\sum_{\alpha \in[0, \lambda]} f_{\alpha} u_{\lambda-\alpha} .
$$

We shall depict the filter as


Figure 1
and call $\left(u_{\lambda}\right)_{\lambda}$ the input and $\left(y_{\lambda}\right)_{\lambda}$ the output of the filter. The concatenation of two filters results in the product filter $\Sigma_{h} \Sigma_{f}$. The difference filter $\Sigma_{f}-\Sigma_{h}$ may be depicted as in Figure 2.


Figure 2
With an element $f=\left(f_{\lambda}\right)_{\lambda \in[0, \infty)} \in \ell_{1}^{N \times M}([0, \infty))$ we may associate a member of $\left(A P W^{k}\right)_{\Lambda \cap S}^{N \times M}$, which with a slight abuse of notation we shall also denote by $f$, and which is defined via

$$
f(t)=\sum_{\lambda \in[0, \infty)} f_{\lambda} e_{\lambda}(t), t \in \mathbb{R}^{k} .
$$

Note that $\Sigma_{h} \Sigma_{f}=\Sigma_{h f}$ and $\Sigma_{h}-\Sigma_{f}=\Sigma_{h-f}$. For a filter $\Sigma_{f}$ we define its norm by

$$
\left\|\Sigma_{f}\right\|=\sup _{u \neq 0} \frac{\left\|\Sigma_{f}(u)\right\|}{\|u\|} .
$$

It is not hard to see that $\left\|\Sigma_{f}\right\|=\|f\|_{\infty}:=\sup _{t \in \mathbb{R}^{k}}\|f(t)\|$.
The model matching problem for linear filters is the following. Given are filters $\Sigma_{f_{1}}, \Sigma_{f_{2}}, \Sigma_{f_{3}}$, find a filter $\Sigma_{h}$ so that the filter $\Sigma_{f_{1}}-\Sigma_{f_{2}} \Sigma_{h} \Sigma_{f_{3}}$ depicted in Figure 3 has minimal possible norm.


Figure 3
Equivalently, given $f_{1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{N \times M}, f_{2} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{N \times P}$ and $f_{3} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{Q \times M}$, find $h \in\left(A P W^{k}\right)_{\Lambda \cap S}^{P \times Q}$ so that $\left\|f_{1}-f_{2} h f_{3}\right\|_{\infty}$ is as small as possible.

We shall assume that $f_{2}$ and $f_{3}$ allow factorizations

$$
f_{2}=H_{2}\left[\begin{array}{cc}
G_{2} & 0  \tag{4.1}\\
0 & 0
\end{array}\right] K_{2}, \quad f_{3}=H_{3}\left[\begin{array}{cc}
G_{3} & 0 \\
0 & 0
\end{array}\right] K_{3},
$$

where $H_{2} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{N \times N}, H_{2}^{-1} \in\left(A P W^{k}\right)_{\Lambda}^{N \times N}, K_{2}^{ \pm 1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{P \times P}, H_{3}^{ \pm 1} \in$ $\left(A P W^{k}\right)_{\Lambda \cap S}^{Q \times Q}, K_{3} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{M \times M}, K_{3}^{-1} \in\left(A P W^{k}\right)_{\Lambda}^{M \times M}$, and $G_{2} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{k_{2} \times k_{2}}$, $G_{2}^{-1} \in\left(A P W^{k}\right)_{\Lambda}^{k_{2} \times k_{2}}, G_{3} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{k_{3} \times k_{3}}, G_{3}^{-1} \in\left(A P W^{k}\right)_{\Lambda}^{k_{3} \times k_{3}}$, for some $k_{2}, k_{3} \in$ $\mathbb{N}$. It follows in particular that $f_{2}(t)$ and $f_{3}(t)$ have constant ranks $k_{2}$ and $k_{3}$, respectively, for every $t \in \mathbb{R}^{k}$.

Under these assumptions we shall provide a solution to the suboptimal problem: Let

$$
\nu>\inf _{h}\left\|f_{1}-f_{2} h f_{3}\right\|_{\infty}
$$

construct one/all $h \in\left(A P W^{k}\right)_{\Lambda \cap S}^{P \times Q}$ such that

$$
\begin{equation*}
\left\|f_{1}-f_{2} h f_{3}\right\|_{\infty}<\nu \tag{4.2}
\end{equation*}
$$

We emphasize that our setting of the suboptimal problem, including factorizations (4.1), involves almost periodic functions with Fourier spectrum in $\Lambda$ only. A suboptimal problem of the above type, but under the more restrictive hypotheses that $f_{2}$ and $f_{3}$ are square size and invertible, was solved in [24]. We solve the above suboptimal problem by reducing it to a contractive four block problem and subsequently applying the results of Section 3. For the reduction we follow closely the ideas of [8].

We first need the following auxiliary result.
Proposition 4.1. Let $f_{2}$ and $f_{3}$ satisfy the condition (4.1). Then there exist

$$
\begin{aligned}
f_{2, i} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{N \times k_{2}}, & f_{2, o} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{k_{2} \times P} \\
f_{3, c i} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{Q \times k_{3}}, & f_{3, c o} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{M \times k_{3}}
\end{aligned}
$$

such that

$$
\begin{equation*}
f_{2}=f_{2, i} f_{2, o}, \quad f_{3}=f_{3, c o} f_{3, c i}, \quad f_{2, i}^{*} f_{2, i} \equiv I, \quad f_{3, c i} f_{3, c i}^{*} \equiv I \tag{4.3}
\end{equation*}
$$

$f_{2, o}$ has a right inverse $f_{2, o}^{\dagger} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{P \times k_{2}}$, and $f_{3, c o}$ has a left inverse $f_{3, c o}^{\dagger} \in$ $\left(A P W^{k}\right)_{\Lambda \cap S}^{k_{3} \times Q}$.

Factorizations (4.3) of $f_{2}$ an $f_{3}$ are to be understood as inner/outer factorizations, where the inner factor is "energy conserving" and the outer factor is "stably invertible".

Proof. We shall prove the statement regarding $f_{2}$. Applying then this result to $f_{3}(-t)^{*}$ will yield the factorization for $f_{3}$.

Let

$$
F=H_{2}\left[\begin{array}{c}
G_{2} \\
0
\end{array}\right]
$$

Then $F^{*} F$ is a positive definite element of $\left(A P W^{k}\right)_{\Lambda}^{k_{2} \times k_{2}}$. By Theorem 2.1 there exists an invertible element $g$ of $\left(A P W^{k}\right)_{\Lambda \cap S}^{k_{2} \times k_{2}}$ so that $F^{*} F=g^{*} g$. Let

$$
f_{2, i}=F g^{-1}, \quad f_{2, o}=\left[\begin{array}{ll}
g & 0
\end{array}\right] K_{2}, \quad f_{2, o}^{\dagger}=K_{2}^{-1}\left[\begin{array}{c}
g^{-1} \\
0
\end{array}\right]
$$

It is straightforward to check that $f_{2, i}, f_{2, o}$, and $f_{2, o}^{\dagger}$ have the required properties.

Theorem 4.2. Let $f_{1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{N \times M}, f_{2} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{N \times P}$ and $f_{3} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{Q \times M}$ be given so that (4.1) is satisfied. Introduce $f_{2, o}, f_{3, c o}, f_{2, i}$ and $f_{3, c i}$ as in Proposition 4.1. Also, let $\nu>0$. Write

$$
\frac{1}{\nu}\left[\begin{array}{c}
f_{2, i}^{*} \\
I-f_{2, i} f_{2, i}^{*}
\end{array}\right] f_{1}\left[I-f_{3, c i}^{*} f_{3, c i} \quad f_{3, c i}^{*}\right]=\left[\begin{array}{ll}
f_{11} & f_{12} \\
f_{21} & f_{22}
\end{array}\right]
$$

and let $s=t=2$ and $q=1$. Then the suboptimal model matching problem (4.2) is solvable if and only if

$$
\left[\begin{array}{cc}
f_{11} & P_{-S \backslash\{0\}}\left(f_{12}\right)  \tag{4.4}\\
f_{21} & f_{22}
\end{array}\right]
$$

has a strictly contractive extension. In that case, for every strictly contractive extension

$$
\left[\begin{array}{ll}
f_{11} & \widetilde{f}_{12}  \tag{4.5}\\
f_{21} & f_{22}
\end{array}\right]
$$

of (4.4) we have that

$$
h:=\nu f_{2, o}^{\dagger}\left(f_{12}-\widetilde{f}_{12}\right) f_{3, c o}^{\dagger}
$$

is a solution to the suboptimal model matching problem (4.2). Conversely, if $h$ is a solution to the suboptimal model matching problem (4.2) then

$$
\left[\begin{array}{cc}
f_{11} & f_{12}-\frac{1}{\nu} f_{2, o} h f_{3, c o} \\
f_{21} & f_{22}
\end{array}\right]
$$

is a strictly contractive extension of (4.4).
Proof. First let (4.5) be a strictly contractive extension of (4.4), i.e.,

$$
\left\|\left[\begin{array}{ll}
f_{11} & \tilde{f}_{12} \\
f_{21} & f_{22}
\end{array}\right]\right\|_{\infty}<1
$$

and $P_{-S \backslash\{0\}}\left(\tilde{f}_{12}\right)=P_{-S \backslash\{0\}}\left(f_{12}\right)$. Letting $h:=\nu f_{2, o}^{\dagger}\left(f_{12}-\widetilde{f}_{12}\right) f_{3, c o}^{\dagger}$, we get that $h \in\left(A P W^{k}\right)_{\Lambda \cap S}^{P \times Q}$, and

$$
\left\|f_{1}-f_{2} h f_{3}\right\|_{\infty}=\|\left[\begin{array}{c}
f_{2, i}^{*} \\
I-f_{2, i} f_{2, i}^{*}
\end{array}\right]\left(f_{1}-f_{2} h f_{3}\right)\left[I-f_{3, c i}^{*} f_{3, c i} \quad f_{3, c i}^{*} \|_{\infty}\right.
$$

the equality follows because $V:=\left[\begin{array}{c}f_{2 i}^{*} \\ I-f_{2 i} f_{2 i}^{*}\end{array}\right]$ is an isometry, i.e., $V^{*} V=I$, and $\left[I-f_{3, c i}^{*} f_{3, c i} \quad f_{3, c i}^{*}\right]$ is a co-isometry. Thus,

$$
\begin{aligned}
\left\|f_{1}-f_{2} h f_{3}\right\|_{\infty} & =\left\|\nu\left[\begin{array}{ll}
f_{11} & f_{12} \\
f_{21} & f_{22}
\end{array}\right]-\left[\begin{array}{cc}
0 & f_{2, o} h f_{3, c o} \\
0 & 0
\end{array}\right]\right\|_{\infty} \\
& =\nu\left\|\left[\begin{array}{ll}
f_{11} & f_{12} \\
f_{21} & f_{22}
\end{array}\right]-\left[\begin{array}{cc}
0 & f_{12}-\widetilde{f}_{12} \\
0 & 0
\end{array}\right]\right\|_{\infty}<\nu
\end{aligned}
$$

For the converse, let $h$ be a solution to the suboptimal model matching problem (4.2), and let $\widetilde{f}_{12}=f_{12}-\frac{1}{\nu} f_{2, o} h f_{3, c o}$. Then $P_{-S \backslash\{0\}}\left(\tilde{f}_{12}\right)=P_{-S \backslash\{0\}}\left(f_{12}\right)$ and

$$
\begin{aligned}
&\left\|\left[\begin{array}{lc}
f_{11} & f_{12}+\left(\widetilde{f}_{12}-f_{12}\right) \\
f_{21} & f_{22}
\end{array}\right]\right\|_{\infty} \\
& \quad=\frac{1}{\nu} \| f_{1}-\left[\begin{array}{ll}
f_{2, i} & I-f_{2, i} f_{2, i}^{*}
\end{array}\right]\left[\begin{array}{cc}
0 & f_{2, o} h f_{3, c o} \\
0 & 0
\end{array}\right] {\left[\begin{array}{c}
I-f_{3, c i}^{*} f_{3, c i} \\
f_{3, c i}
\end{array}\right] \|_{\infty} } \\
&=\frac{1}{\nu}\left\|f_{1}-f_{2} h f_{3}\right\|_{\infty}<1 .
\end{aligned}
$$

By using Theorem 3.3 it is now straightforward to obtain a full solution to the suboptimal model matching problem in terms of Hankel type operators on Besikovitch space. The solution to the classical model matching problem ( $k=1$, $\Lambda=\mathbb{Z}$ ) may be found in [8, Section 8.1] for rational matrix functions.

Note that Theorem 4.2 together with Theorem 3.3 provide frequency-domain formulas for solutions of the suboptimal problem. These formulas may be not computationally practical. One might anticipate that a conversion to state space formulas, if possible, would give more practical formulas. However, this must remain a subject for future research.

One may wonder whether there exist systems in time domain giving rise to transfer functions that are of the type $f \in\left(A P W^{k}\right)_{\Lambda \cap S}^{N \times M}$ as above. It has been known since the 80s that delay systems give rise to almost periodic transfer functions (see, e.g., [18]). Another source are multidimensional systems; [1] is a general reference. We mention here also the following example (see [16], [17]), which is relevant to the periodic case $\left(\Lambda=\mathbb{Z}^{k}\right)$; for notational simplicity, we let $k=2$ :

$$
\Sigma:\left\{\begin{array}{c}
x\left(n_{1}, n_{2}\right)=A_{1} x\left(n_{1}-1, n_{2}\right)+A_{2} x\left(n_{1}, n_{2}-1\right)+ \\
\\
+B_{1} u\left(n_{1}-1, n_{2}\right)+B_{2} u\left(n_{1}, n_{2}-1\right) \\
y\left(n_{1}, n_{2}\right)=C_{1} x\left(n_{1}-1, n_{2}\right)+C_{2} x\left(n_{1}, n_{2}-1\right)+ \\
\\
+D_{1} u\left(n_{1}-1, n_{2}\right)+D_{2} u\left(n_{1}, n_{2}-1\right) \\
x\left(n_{1}, n_{2}\right)=0 \quad \text { for } \quad\left(n_{1}, n_{2}\right) \in \mathbb{Z}^{2} \quad \text { such that } n_{1} \leq 0, \text { and } n_{1}+n_{2}=0 \\
x\left(n_{1}, n_{2}\right)=0 \quad \text { for } \quad\left(n_{1}, n_{2}\right) \in \mathbb{Z}^{2} \quad \text { such that } n_{1} \geq 1, \text { and } n_{1}+n_{2}=1
\end{array}\right.
$$

Here $A_{1}, A_{2}, B_{1}, B_{2}, C_{1}, C_{2}, D_{1}$, and $D_{2}$ are constant matrices of appropriate sizes. It should be noted that the initial conditions on $x$ are slightly different from those in [16]. The vectors $u(n), x(n)$ and $y(n)$, where $n=\left(n_{1}, n_{2}\right)$, are usually referred to as the input, state and output vectors, respectively. We now apply the $z$-transform, as follows. Let $S$ be the halfspace

$$
S=\left\{\left(v_{1}, v_{2}\right) \in \mathbb{R}^{2}: \text { either } \quad v_{1}+v_{2}>0 \quad \text { or } \quad v_{1} \leq 0 \quad \text { and } \quad v_{1}+v_{2}=0\right\}
$$

Letting

$$
\begin{gathered}
\widehat{x}\left(z_{1}, z_{2}\right)=\sum_{\left(n_{1}, n_{2}\right) \in S \cap \mathbb{Z}^{2}} x\left(n_{1}, n_{2}\right) z_{1}^{n_{1}} z_{2}^{n_{2}}, \\
\widehat{y}\left(z_{1}, z_{2}\right)=\sum_{\left(n_{1}, n_{2}\right) \in S \cap \mathbb{Z}^{2}+(0,1)} y\left(n_{1}, n_{2}\right) z_{1}^{n_{1}} z_{2}^{n_{2}},
\end{gathered}
$$

$$
\widehat{u}\left(z_{1}, z_{2}\right)=\sum_{\left(n_{1}, n_{2}\right) \in S \cap \mathbb{Z}^{2}} u\left(n_{1}, n_{2}\right) z_{1}^{n_{1}} z_{2}^{n_{2}},
$$

and solving for $\widehat{y}$ in terms of $\widehat{u}$, we obtain that

$$
\begin{equation*}
\widehat{y}\left(z_{1}, z_{2}\right)=f\left(z_{1}, z_{2}\right) \widehat{u}\left(z_{1}, z_{2}\right) \tag{4.6}
\end{equation*}
$$

where

$$
f\left(z_{1}, z_{2}\right)=z_{1} D_{1}+z_{2} D_{2}+\left(z_{1} C_{1}+z_{2} C_{2}\right)\left(I-z_{1} A_{1}-z_{2} A_{2}\right)^{-1}\left(z_{1} B_{1}+z_{2} B_{2}\right)
$$

(To make (4.6) precise, one has to assume that $u(0,0)=0$.) The function $f$ is known as the transfer function of the system $\Sigma$. Assuming that $A_{1}$ and $A_{2}$ are such that the spectral radius of $z_{1} A_{1}+z_{2} A_{2}$ is less than one for every $z_{1}, z_{2} \in \mathbb{T}$, upon the substitution $z_{j}=e^{i t_{j}}, j=1,2$, we obtain that $f \in\left(A P W^{2}\right)_{\mathbb{Z}^{2} \cap S}^{N \times M}$.

## 5. Multiblocks having more general patterns: An example

Using the semi-band structure version of the band method developed in [15] instead of the standard band method we used in previous sections, we are able to treat more general positive extension problems. In this section we shall treat one example illustrating the main ideas. Though far more general results can be stated and proved (the theory developed in [14] gives an indication how to produce the most general setup), the notational complexity of doing this is so overwhelming that we restrict ourselves to the following situation only.

Fix a halfspace $S \subset \mathbb{R}^{k}$, an additive subgroup $\Lambda \subseteq \mathbb{R}^{k}$ and a vector $\nu \in S$. Let now

$$
\Delta_{+}=(\nu-S) \cap S \cap \Lambda, \quad \text { and } \quad \Delta=\Delta_{+} \cup\left(-\Delta_{+}\right)
$$

A positive extension problem can be stated as follows: Given

$$
\begin{aligned}
& k_{11} \in\left(A P W^{k}\right)_{\Lambda}^{n_{1} \times n_{1}}, k_{12}^{-}=\left(k_{21}^{+}\right)^{*} \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{1} \times n_{2}}, k_{22} \in\left(A P W^{k}\right)_{\Lambda}^{n_{2} \times n_{2}}, \\
& k_{24}^{-}=\left(k_{42}^{+}\right)^{*} \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{2} \times n_{4}}, k_{33}^{c} \in\left(A P W^{k}\right)_{\Delta}^{n_{3} \times n_{3}}, k_{34}=k_{43}^{*} \in\left(A P W^{k}\right)_{\Lambda}^{n_{3} \times n_{4}}, \\
& k_{44} \in\left(A P W^{k}\right)_{\Lambda}^{n_{4} \times n_{4}},
\end{aligned}
$$

find

$$
k=\left(k_{i j}\right)_{i, j=1}^{4} \in\left(A P W^{k}\right)_{\Lambda}^{\left(n_{1}+\cdots+n_{4}\right) \times\left(n_{1}+\cdots+n_{4}\right)}
$$

so that

$$
k>0, \quad P_{-S}\left(k_{12}\right)=k_{12}^{-}, \quad P_{-S}\left(k_{24}\right)=k_{24}^{-}, \quad P_{\Delta}\left(k_{33}\right)=k_{33}^{c}
$$

We call $k$ a positive extension of

$$
k_{c}=\left[\begin{array}{cccc}
k_{11} & k_{12}^{-} & 0 & 0 \\
k_{21}^{+} & k_{22} & 0 & k_{24}^{-} \\
0 & 0 & k_{33}^{c} & k_{34} \\
0 & k_{42}^{+} & k_{43} & k_{44}
\end{array}\right]
$$

In order to use the framework of [15] we need to introduce several subspaces of $\mathcal{A}:=\left(A P W^{k}\right)_{\Lambda}^{Q \times Q}$, where $Q=n_{1}+\cdots+n_{4}$. In the following formulas we shall omit the subscript $\Lambda$ (for example, $\left(A P W^{k}\right)_{\Delta_{+} \backslash\{0\}}$ is to be understood as $\left.\left(A P W^{k}\right)_{\left(\Delta_{+} \backslash\{0\}\right) \cap \Lambda}\right)$, as well as the sizes of the individual blocks as they are clear from the block's position in a matrix.

Let

$$
\begin{aligned}
& \mathcal{A}_{+}^{0}=\left(\mathcal{A}_{-}^{0}\right)^{*} \\
&=\left[\begin{array}{cccc}
\left(A P W^{k}\right)_{S \backslash\{0\}} & \left(A P W^{k}\right) & \left(A P W^{k}\right) & \left(A P W^{k}\right) \\
0 & \left(A P W^{k}\right)_{S \backslash\{0\}} & \left(A P W^{k}\right) & \left(A P W^{k}\right) \\
0 & 0 & \left(A P W^{k}\right)_{S \backslash\{0\}} & 0 \\
0 & 0 & \left(A P W^{k}\right) & \left(A P W^{k}\right)_{S \backslash\{0\}}
\end{array}\right], \\
& \widetilde{\mathcal{A}}_{d}=\mathcal{A}_{d}=\left[\begin{array}{ccc}
\left(A P W^{k}\right)_{\{0\}} & & 0 \\
& \ddots & \\
0 & & \left(A P W^{k}\right)_{\{0\}}
\end{array}\right],
\end{aligned}
$$

$$
\begin{gathered}
\mathcal{A}_{2}^{0}=\left[\begin{array}{cccc}
\left(A P W^{k}\right)_{S \backslash\{0\}} & \left(A P W^{k}\right)_{-S} & 0 & 0 \\
0 & \left(A P W^{k}\right)_{S \backslash\{0\}} & 0 & \left(A P W^{k}\right)_{-S} \\
0 & 0 & \left(A P W^{k}\right)_{\Delta_{+} \backslash\{0\}} & 0 \\
0 & 0 & \left(A P W^{k}\right) & \left(A P W^{k}\right)_{S \backslash\{0\}}
\end{array}\right], \\
\widetilde{\mathcal{A}}_{4}^{*}=\mathcal{A}_{1}=\left[\begin{array}{cccc}
0 & \left(A P W^{k}\right)_{S \backslash\{0\}} & \left(A P W^{k}\right) & \left(A P W^{k}\right) \\
0 & 0 & \left(A P W^{k}\right) & \left(A P W^{k}\right)_{S \backslash\{0\}} \\
0 & 0 & \left(A P W^{k}\right)_{S \backslash \Delta_{+}} & 0 \\
0 & 0 & 0 & 0
\end{array}\right],
\end{gathered}
$$

$$
\widetilde{\mathcal{A}}_{-}^{0}=\left(\widetilde{\mathcal{A}}_{+}^{0}\right)^{*}
$$

$$
=\left[\begin{array}{cccc}
\left(A P W^{k}\right)_{-S \backslash\{0\}} & 0 & 0 & 0 \\
\left(A P W^{k}\right) & \left(A P W^{k}\right)_{-S \backslash\{0\}} & 0 & 0 \\
\left(A P W^{k}\right) & \left(A P W^{k}\right) & \left(A P W^{k}\right)_{-S \backslash\{0\}} & 0 \\
\left(A P W^{k}\right) & \left(A P W^{k}\right) & \left(A P W^{k}\right) & \left(A P W^{k}\right)_{-S \backslash\{0\}}
\end{array}\right]
$$

$$
\widetilde{\mathcal{A}}_{3}^{0}=\left[\begin{array}{cccc}
\left(A P W^{k}\right)_{-S \backslash\{0\}} & 0 & 0 & 0 \\
\left(A P W^{k}\right)_{S} & \left(A P W^{k}\right)_{-S \backslash\{0\}} & 0 & 0 \\
0 & 0 & \left(A P W^{k}\right)_{-\Delta_{+} \backslash\{0\}} & 0 \\
0 & \left(A P W^{k}\right)_{S} & \left(A P W^{k}\right) & \left(A P W^{k}\right)_{-S \backslash\{0\}}
\end{array}\right],
$$

$$
\begin{gathered}
\mathcal{A}_{2}=\mathcal{A}_{2}^{0} \dot{+} \mathcal{A}_{d}, \widetilde{\mathcal{A}}_{3}=\widetilde{\mathcal{A}}_{3}^{0} \dot{+} \widetilde{\mathcal{A}}_{d} \\
\mathcal{A}_{-}^{*}=\mathcal{A}_{+}=\mathcal{A}_{+}^{0}+\mathcal{A}_{d}, \widetilde{\mathcal{A}}_{-}=\widetilde{\mathcal{A}}_{+}^{*}=\widetilde{\mathcal{A}}_{-}^{0} \dot{+} \widetilde{\mathcal{A}}_{d}
\end{gathered}
$$

It is straightforward to check that the following multiplication tables hold:

| $\cdot$ | $\mathcal{A}_{-}^{0}$ | $\mathcal{A}_{d}$ | $\mathcal{A}_{+}^{0}$ |  | $\cdot$ | $\widetilde{\mathcal{A}}_{-}^{0}$ | $\widetilde{\mathcal{A}}_{d}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $\mathcal{A}_{-}^{0}$ | $\mathcal{A}_{-}^{0}$ | $\mathcal{A}_{-}^{0}$ | $\mathcal{A}$ |  | $\widetilde{\mathcal{A}}_{-}^{0}$ | $\widetilde{\mathcal{A}}_{-}^{0}$ | $\widetilde{\mathcal{A}}_{-}^{0}$ |
| $\mathcal{A}_{d}$ | $\mathcal{A}_{-}^{0}$ | $\mathcal{A}_{d}$ | $\mathcal{A}_{+}^{0}$ |  | $\widetilde{\mathcal{A}}_{d}$ | $\widetilde{\mathcal{A}}_{-}^{0}$ | $\widetilde{\mathcal{A}}_{d}$ |
| $\mathcal{A}_{+}^{0}$ |  |  |  |  |  |  |  |
| $\mathcal{A}_{+}^{0}$ | $\mathcal{A}$ | $\mathcal{A}_{+}^{0}$ | $\mathcal{A}_{+}^{0}$ |  | $\widetilde{\mathcal{A}}_{+}^{0}$ | $\widetilde{\mathcal{A}}$ | $\widetilde{\mathcal{A}}_{+}^{0}$ |
| $\widetilde{\mathcal{A}}_{+}^{0}$ |  |  |  |  |  |  |  |.

Moreover, we have:

$$
\mathcal{A}_{+}=\mathcal{A}_{1} \dot{+} \mathcal{A}_{2}, \quad \widetilde{\mathcal{A}}_{-}=\widetilde{\mathcal{A}}_{4} \dot{+} \tilde{\mathcal{A}}_{3}
$$

and

$$
\begin{gathered}
\mathcal{A}_{1} \mathcal{A}_{+} \subseteq \mathcal{A}_{1}, \quad \mathcal{A}_{2} \mathcal{A}_{d} \subseteq \mathcal{A}_{2}, \quad \mathcal{A}_{1}^{*} \mathcal{A}_{2} \subseteq \mathcal{A}_{-}^{0}, \\
\widetilde{\mathcal{A}}_{4} \widetilde{\mathcal{A}}_{-} \subseteq \widetilde{\mathcal{A}}_{4}, \quad \widetilde{\mathcal{A}}_{3} \widetilde{\mathcal{A}}_{d} \subseteq \widetilde{\mathcal{A}}_{3}, \quad \widetilde{\mathcal{A}}_{4}^{*} \widetilde{\mathcal{A}}_{3} \subseteq \widetilde{\mathcal{A}}_{+}^{0}, \\
\widetilde{\mathcal{A}}_{-} \mathcal{A}_{2} \subseteq \mathcal{A}_{-}^{0}+\mathcal{A}_{2}, \quad \mathcal{A}_{2} \mathcal{A}_{-} \subseteq \mathcal{A}_{-}^{0}+\mathcal{A}_{2}, \quad \mathcal{\mathcal { A }}_{3}^{*} \mathcal{A}_{-} \subseteq \mathcal{A}_{-}^{0}+\mathcal{A}_{2} .
\end{gathered}
$$

We shall let $\mathcal{B}_{1}, \mathcal{B}_{2}^{0}$, etc., denote the closures of $\mathcal{A}_{1}, \mathcal{A}_{2}^{0}$, etc., in $\mathcal{B}=\left(B^{k}\right)_{\Lambda}^{Q \times Q}$. In addition, $P_{1}, P_{2}^{0}$, etc., are the orthogonal projections of $\mathcal{B}$ onto $\mathcal{B}_{1}, \mathcal{B}_{2}^{0}$, etc., respectively. We observe that when $k>0$ is in $\mathcal{A}$, then we have that $k$ allows the factorizations

$$
\begin{equation*}
k=\left(e+k_{+}\right)^{*} k_{d}\left(e+k_{+}\right) \tag{5.2}
\end{equation*}
$$

where $k_{+},\left(e+k_{+}\right)^{-1}-e \in \mathcal{A}_{+}^{0}$ and $k_{d} \in \mathcal{A}_{d}$. Indeed, applying a constant permutation matrix that interchanges the third and fourth blocks, the factorization (5.2) is reduced to that of Lemma 2.2.

We are now ready to state the main result in this section.
Theorem 5.1. Let $k_{c}=k_{c}^{*} \in \mathcal{A}_{2}^{0} \dot{+} \mathcal{A}_{d} \dot{+} \mathcal{A}_{2}^{0 *}=\widetilde{\mathcal{A}}_{3}^{0} \dot{+} \widetilde{\mathcal{A}}_{d} \dot{+} \widetilde{\mathcal{A}}_{3}^{0 *}$ be given. Then the following are equivalent:
(i) $k_{c}$ has a positive extension.
(ii) The operator $T: \mathcal{B}_{2} \rightarrow \mathcal{B}_{2}$ defined by $T(g)=P_{2}\left(k_{c} g\right)$ is positive definite.
(iii) The operator $\widetilde{T}: \widetilde{\mathcal{B}}_{3} \rightarrow \widetilde{\mathcal{B}}_{3}$ defined by $\widetilde{T}(g)=P_{3}\left(k_{c} g\right)$ is positive definite.

In that case, let

$$
x=T^{-1}(e), \quad y=T^{-1}(e)
$$

Then

$$
k_{0}=x^{*-1} P_{d}(x) x^{-1}=y^{*-1} \widetilde{P}_{d}(y) y^{-1}
$$

is a positive extension of $k_{c}$. Moreover, if we let $u=x P_{d}(x)^{-\frac{1}{2}}$ and $v=y \widetilde{P}_{d}(y)^{-\frac{1}{2}}$, then $k$ is a positive extension of $k_{c}$ if and only if

$$
k=(v g+u)^{*-1}\left(e-g^{*} g\right)(v g+u)^{-1}
$$

for some (unique!) $g \in \mathcal{A}_{1}$ with $\|g\|_{\infty}<1$. Lastly, if $k$ is a positive extension of $k_{c}$ with factorization (5.2), then

$$
k_{d} \leq P_{d}(x)^{-1}
$$

and equality holds if and only if $k=k_{0}$.

Proof. As before, the main effort goes into showing that $x^{-1} \in \mathcal{A}_{+}$and $y^{-1} \in \widetilde{\mathcal{A}}_{-}$. We proceed analogously to the proof of Proposition 2.6.

Observing that $T>0$ when applied to $g=\left(g_{i j}\right)_{i, j=1}^{4} \in \mathcal{B}_{2}$ with $g_{i j}=0, j \neq 1$, yields that $k_{11}>0$, and one checks that $x_{11}=\left(e+k_{11,+}\right)^{-1} k_{11, d}^{-1}$, where

$$
k_{11}=\left(e+k_{11,+}\right)^{*} k_{11, d}\left(e+k_{11,+}\right)
$$

is a right canonical factorization.
Next observe that by applying $T$ to elements $g=\left(g_{i j}\right)_{i, j=1}^{4} \in \mathcal{B}_{2}$ with $g_{i j}=0$ for $j \neq 2$, we get that the operator

$$
\left[\begin{array}{cc}
\widetilde{T}_{11} & H_{12}  \tag{5.3}\\
H_{12}^{*} & T_{22}
\end{array}\right]
$$

on $\left(\mathcal{B}^{k}\right)_{\Lambda \cap(-S)}^{n_{1} \times n_{2}} \oplus\left(\mathcal{B}^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}}$ is positive definite. Here

$$
\widetilde{T}_{11}\left(g_{1}\right)=P_{-S}\left(k_{11} g_{1}\right), \quad g_{1} \in\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{1} \times n_{2}}
$$

and

$$
H_{12}\left(g_{2}\right)=P_{-S}\left(k_{12}^{-} g_{2}\right), \quad T_{22}\left(g_{2}\right)=P_{S}\left(k_{22} g_{2}\right), \quad g_{2} \in\left(B^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}}
$$

Since $\widetilde{T}_{11}>0$ and $T_{22}>0$, it follows that their symbols $k_{11}$ and $k_{22}$ are positive definite. Write the canonical factorizations (which exist by Theorem 2.1)

$$
k_{11}=\left(e+k_{1,-}\right)^{*} k_{1, d}\left(e+k_{1,-}\right), \quad k_{22}=\left(e+k_{2,+}\right)^{*} k_{2, d}\left(e+k_{2,+}\right),
$$

where

$$
k_{1,-},\left(e+k_{1,-}\right)^{-1}-e \in\left(A P W^{k}\right)_{\Lambda \cap(-S \backslash\{0\})}^{n_{1} \times n_{1}}, k_{1, d} \in \mathbb{C}^{n_{1} \times n_{1}}, \quad k_{1, d}>0
$$

and

$$
k_{2,+},\left(e+k_{2,+}\right)^{-1}-e \in\left(A P W^{k}\right)_{\Lambda \cap(S \backslash\{0\})}^{n_{2} \times n_{2}}, k_{2, d} \in \mathbb{C}^{n_{2} \times n_{2}}, \quad k_{2, d}>0 .
$$

Introducing the operators

$$
H_{1}:\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{1} \times n_{2}} \rightarrow\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{1} \times n_{2}} \quad \text { and } \quad H_{2}:\left(B^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}} \rightarrow\left(B^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}}
$$

by the formulas

$$
H_{1} g=k_{1, d}^{\frac{1}{2}}\left(e+k_{1,-}\right) g, \quad H_{2} g=k_{2, d}^{\frac{1}{2}}\left(e+k_{2,+}\right) g
$$

we get that

$$
\widetilde{T}_{11}=H_{1}^{*} H_{1}, \quad T_{22}=H_{2}^{*} H_{2} .
$$

Consider the operator

$$
H_{3}:=H_{1}^{*-1} H_{12} H_{2}^{-1}:\left(B^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}} \rightarrow\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{1} \times n_{2}} .
$$

We have

$$
H_{3} g=P_{-S}\left(\left(e+k_{1,-}\right)^{*-1} k_{1, d}^{-\frac{1}{2}} k_{12}^{-}\left(e+k_{2,+}\right)^{-1} k_{2, d}^{-\frac{1}{2}} g\right), \quad g \in\left(B^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}}
$$

Since the operator (5.3) is positive definite, it follows that $\left\|H_{3}\right\|<1$. Applying Theorem 2.5, we get that

$$
\begin{aligned}
\beta & :=H_{3}\left(I-H_{3}^{*} H_{3}\right)^{-1}(e) \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{1} \times n_{2}}, \\
\delta & :=\left(I-H_{3}^{*} H_{3}\right)^{-1}(e) \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}},
\end{aligned}
$$

$\delta^{-1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}}$ and $M\{\delta\}>0$. It is now straightforward to check, analogously to (2.21) and (2.22), that

$$
x_{12}=\left(e+k_{1,-}\right)^{-1} k_{1, d}^{-\frac{1}{2}} \beta k_{2, d}^{-\frac{1}{2}}, \quad x_{22}=\left(e+k_{2,+}\right)^{-1} k_{2, d}^{-\frac{1}{2}} \delta k_{2, d}^{-\frac{1}{2}}
$$

and moreover

$$
x_{12} \in\left(A P W^{k}\right)_{\Lambda \cap(-S)}^{n_{1} \times n_{2}}, \quad x_{22}^{ \pm 1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{2} \times n_{2}}, \quad M\left\{x_{22}\right\}>0
$$

In order to obtain that $x_{33}$ is of the desired form, observe that $T$ applied to $\left(g_{i j}\right)_{i, j=1}^{4} \in \mathcal{B}_{2}$ with $g_{i j}=0, j \neq 3$, yields that

$$
\left[\begin{array}{ll}
T_{33} & M_{43}^{*} \\
M_{43} & M_{44}
\end{array}\right]>0
$$

where

$$
\begin{gathered}
T_{33} g_{3}=P_{\Delta_{+}}\left(k_{33}^{c} g_{3}\right), M_{43} g_{3}=k_{43} g_{3}, g_{3} \in\left(B^{k}\right)_{\Delta_{+}}^{n_{3} \times n_{3}} \\
M_{44} g_{4}=k_{44} g_{4}, g_{4} \in\left(B^{k}\right)^{n_{4} \times n_{3}}
\end{gathered}
$$

Notice that operator $T_{33}-M_{43}^{*} M_{44}^{-1} M_{43}$ is positive definite, and

$$
\left(T_{33}-M_{43}^{*} M_{44}^{-1} M_{43}\right) g_{3}=P_{\Delta_{+}}\left(\left(k_{33}^{c}-k_{34} k_{44}^{-1} k_{43}\right) g\right), \quad g_{3} \in\left(B^{k}\right)_{\Delta_{+}}^{n_{3} \times n_{3}}
$$

It follows from (the proof of) Theorem 4.4.1 in [25] that

$$
x_{33}=\left(T_{33}-M_{43}^{*} M_{44}^{-1} M_{43}\right)^{-1}\left(I_{n_{3}}\right)
$$

has the property that $x_{33} \in\left(A P W^{k}\right)_{\Delta_{+}}^{n_{3} \times n_{3}}, x_{33}^{-1} \in\left(A P W^{k}\right)_{\Lambda \cap S}^{n_{3} \times n_{3}}$ and $M\left\{x_{33}\right\}>0$. In addition,

$$
x_{43}=-k_{44}^{-1} k_{43} x_{33} \in\left(A P W^{k}\right)_{\Lambda}^{n_{4} \times n_{3}}
$$

For the fourth column of $x$ observe that $T$ applied to $\left(g_{i j}\right)_{i, j=1}^{4} \in \mathcal{B}_{2}$ with $g_{i j}=0$, $j \neq 4$, yields that

$$
\left[\begin{array}{cc}
\widetilde{T}_{22} & H_{24} \\
H_{24}^{*} & T_{44}
\end{array}\right]>0
$$

where

$$
\begin{gathered}
\widetilde{T}_{22}\left(g_{2}\right)=P_{-S}\left(k_{22} g_{2}\right), \quad g_{2} \in\left(B^{k}\right)_{\Lambda \cap(-S)}^{n_{2} \times n_{4}}, \\
H_{24}\left(g_{4}\right)=P_{-S}\left(k_{24}^{-} g_{4}\right), \quad T_{44}\left(g_{4}\right)=P_{S}\left(k_{44} g_{4}\right), \quad g_{4} \in\left(B^{k}\right)_{\Lambda \cap S}^{n_{4} \times n_{4}}
\end{gathered}
$$

The same reasoning as for the second column of $x$ yields that $x_{24}$ and $x_{44}$ are of the required form.

We are now in a position to complete the proof of Theorem 5.1, by applying Theorems 1.1, 1.3, 1.4, 1.6, 1.10 in [15] to the current setting. The only additional observation that needs to be made is the following equality (in the notation of Theorem 1.6 of [15]):

$$
\begin{aligned}
& \left(-c^{*} v g+c u\right)(v g+u)^{-1}+(v g+u)^{*-1}\left(-c^{*} v g+c u\right)^{*} \\
& =(v g+u)^{*-1}\left[-g^{*} v^{*} c^{*} v g+g^{*} v^{*} c u-u^{*} c^{*} v g+u^{*} c u\right. \\
& \left.\quad-g^{*} v^{*} c v g+u^{*} c^{*} v g-g^{*} v^{*} c u+u^{*} c^{*} u\right](v g+u)^{-1} \\
& =(v g+u)^{*-1}\left(-g^{*} g+e\right)(v g+u)^{-1}
\end{aligned}
$$

where we used that $c+c^{*}=u^{*-1} u^{-1}=v^{*-1} v^{-1}$.

We may state an analog of Proposition 2.4 in the current setting as well.
Proposition 5.2. Let $k_{c}=k_{c}^{*} \in \mathcal{M}_{c}$ be given and suppose that $k_{c}$ has a positive extension. Let $k_{0}$ be defined as in Theorem 5.1. If $k$ is a positive extension of $k_{c}$ then

$$
\Delta(k) \leq \Delta\left(k_{0}\right)
$$

and equality holds if and only if $k=k_{0}$.
The proof of this proposition is analogous to the proof of Proposition 2.4, and is left to the reader.

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