# Fiedler companion linearizations for rectangular matrix polynomials ${ }^{\boldsymbol{\alpha}}$ 

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#### Abstract

The development of new classes of linearizations of square matrix polynomials that generalize the classical first and second Frobenius companion forms has attracted much attention in the last decade. Research in this area has two main goals: finding linearizations that retain whatever structure the original polynomial might possess, and improving properties that are essential for accurate numerical computation, such as eigenvalue condition numbers and backward errors. However, all recent progress on linearizations has been restricted to square matrix polynomials. Since rectangular polynomials arise in many applications, it is natural to investigate if the new classes of linearizations can be extended to rectangular polynomials. In this paper, the family of Fiedler linearizations is extended from square to rectangular matrix polynomials, and it is shown that minimal indices and bases of polynomials can be recovered from those of any linearization in this class via the same simple procedures developed previously for square polynomials. Fiedler linearizations are one of the most important classes of linearizations introduced in recent years, but their generalization to rectangular polynomials is nontrivial, and requires a completely different approach to the one used in the square case. To the best of our knowledge, this is the first class of new linearizations that has been generalized to rectangular polynomials.


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## 1. Introduction

We consider in this paper $m \times n$ matrix polynomials with degree $k \geq 2$ of the form

$$
\begin{equation*}
P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}, \quad A_{0}, \ldots, A_{k} \in \mathbb{F}^{m \times n}, \quad A_{k} \neq 0 \tag{1}
\end{equation*}
$$

where $\mathbb{F}$ is an arbitrary field and $\lambda$ is a scalar variable in $\mathbb{F}$. Our main focus is on rectangular matrix polynomials, i.e., with $m \neq n$, although new results for square polynomials are also presented. A matrix polynomial $P(\lambda)$ is said to be singular either if it is rectangular, or it is square and $\operatorname{det} P(\lambda)$ is identically zero, i.e., if all the coefficients of $\operatorname{det} P(\lambda)$ are zero; otherwise $P(\lambda)$ is regular.

Matrix polynomials arise in many applications like systems of differential-algebraic equations, vibration analysis of structural systems, acoustics, fluid-structure interaction problems, computer graphics, signal processing, control theory, and linear system theory [4, 18, 24, 25, 29, 30, 31, 32]. Rectangular matrix polynomials appear mainly in control theory and linear system theory. The magnitudes that are

[^0]usually relevant in the applications of regular matrix polynomials are their finite and infinite eigenvalues and the corresponding eigenvectors [18], while in applications of singular polynomials their minimal indices and bases also play an important role [15, 24].

A standard way of dealing, both theoretically and numerically, with a matrix polynomial $P(\lambda)$ is to convert it into an equivalent matrix pencil. This process is known as linearization [18], and is explained in Section 2. The classical approach uses the first and second Frobenius companion forms (4) and (5) as linearizations. However, these companion forms usually do not share any algebraic structure that $P(\lambda)$ might have, and their use in numerical computations, via well-established algorithms for pencils [3, 7, 8, 19, 33], may destroy important qualitative features of the eigenvalues/eigenvectors and minimal indices/bases as a consequence of rounding errors. In addition, the condition numbers of the eigenvalues in the Frobenius companion linearizations may be much larger than in $P(\lambda)$, and small eigenvalue backward errors in the linearization do not guarantee small backward errors in the polynomial [21, 22].

These difficulties have motivated intense activity in the last decade towards the development of new classes of linearizations. At first, only linearizations for regular matrix polynomials were considered $[1,2$, $23,27,28]$, while more recently square singular polynomials have also received attention $[10,11,12,34]$. However, all this recent progress on linearizations has been restricted to square matrix polynomials. The main goal of this paper is to extend one of the most relevant new classes of linearizations from square to rectangular matrix polynomials. This is the family of Fiedler pencils, which was originally introduced by Fiedler for scalar polynomials in [14], generalized to regular matrix polynomials over $\mathbb{C}$ in [2], and then extended and further analyzed in [11] for both regular and singular square matrix polynomials over an arbitrary field $\mathbb{F}$.

Fiedler pencils of square matrix polynomials $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ enjoy a number of important properties that make them attractive candidates for generalization to rectangular polynomials. They are strong linearizations for any square polynomial, regular or singular, over an arbitrary field, and the coefficients of these pencils are simply constructed as block partitioned matrices whose blocks are either $0, \pm I$, or $\pm A_{i}, i=0,1, \ldots, k[11]$. This means that they are all companion forms in the sense of $[12$, Definition 1.1]. Fiedler pencils allow us to very easily recover not only the eigenvalues, but also the eigenvectors, minimal indices, and minimal bases of $P(\lambda)$ from the corresponding magnitudes of the pencil [11]. These pencils can also be generalized to preserve structures of polynomials that are important in applications, like symmetry and palindromicity $[2,12,34]$. No other class of linearizations introduced in recent years simultaneously satisfy all these properties. In fact, for other important classes of new linearizations [27], it is very easy to find pencils that cannot be extended to rectangular matrix polynomials as a consequence of obvious size constraints.

We remark that the extension of Fiedler pencils from square to rectangular matrix polynomials is not direct, since the original definition cannot be applied to rectangular polynomials. This issue is discussed in Section 3.2. Therefore we follow an approach completely different than the one considered in $[2,11,14]$ for square polynomials. This approach in based on the construction presented in Algorithm 2, from which the main Definition 3.8 is established. With this definition in hand, and after considerable technical effort, we prove in Theorem 4.5 that Fiedler pencils of rectangular matrix polynomials are always strong linearizations over arbitrary fields, again using new techniques. Finally, simple recovery procedures for minimal indices and bases are presented in Corollaries 5.4 and 5.7. These recovery rules are essentially the same as the ones derived for square polynomials in [11]. Although the new proofs and definitions may seem complicated, we emphasize that the key idea is very simple: we perform the same operations that we would do in the square case, but proving that the rectangular matrices that appear are always conformable for multiplication. This requires a substantial amount of care. Another essential difference between Fiedler pencils for rectangular and square polynomials is that in the rectangular case there are Fiedler pencils of several different sizes; indeed the two Frobenius companion forms are always the Fiedler pencils with largest and smallest sizes, while the other Fiedler pencils have intermediate sizes. This always makes one of the Frobenius companion forms a privileged choice to work with rectangular matrix polynomials, although the low band-width structure of some other Fiedler pencils might make them useful in certain situations.

The paper is organized as follows. In Section 2 we introduce the basic definitions and notation used throughout the paper. In Section 3 we recall first the notion of Fiedler pencils for square matrix polynomials, then present an algorithm to construct these pencils in a manner that readily generalizes to
rectangular matrix polynomials. It is thus by means of this algorithm that we are able to extend the notion of Fiedler pencils to rectangular polynomials. In the last part of Section 3, we establish the relationship between the reversal of a polynomial and the reversal of any of its Fiedler pencils (Theorem 3.14). This relationship is needed to prove that Fiedler pencils of rectangular polynomials are always strong linearizations in Section 4. Section 5 establishes very simple formulae for the recovery of the minimal indices and bases of a matrix polynomial from the minimal indices and bases of any of its Fiedler pencils. Finally, Section 6 gives some conclusions and describes possible future work motivated by the results in this paper.

## 2. Basic notation and definitions

We present in this section some basic concepts related to rectangular matrix polynomials. The reader can find more information in [10, Section 2] and [11, Section 2], where these concepts were presented in more detail for square polynomials. In the rest of the paper we adopt the following notation: $0_{d}$ and $I_{d}$ are used to denote the $d \times d$ zero and identity matrices, respectively. If there is no risk of confusion, then the sizes are not indicated and we simply write 0 or $I$. Two $m \times n$ matrix polynomials $P(\lambda)$ and $Q(\lambda)$ are strictly equivalent if there exist two constant nonsingular matrices $E$ and $F$ such that $P(\lambda)=E Q(\lambda) F$. We emphasize that any equation in this paper involving expressions in $\lambda$ is to be understood as a formal algebraic identity, and not just as an equality of functions on the field $\mathbb{F}$. For finite fields $\mathbb{F}$ this distinction is important, and we will always intend the stronger meaning of a formal algebraic identity.

Let $\mathbb{F}(\lambda)$ denote the field of rational functions with coefficients in $\mathbb{F}$, so that $\mathbb{F}(\lambda)^{n \times 1}$ is the vector space of column $n$-tuples with entries in $\mathbb{F}(\lambda)$. The normal rank of a matrix polynomial $P(\lambda)$, denoted nrank $P(\lambda)$, is the rank of $P(\lambda)$ considered as a matrix with entries in $\mathbb{F}(\lambda)$, or equivalently, the size of the largest non-identically zero minor of $P(\lambda)$ [16]. A finite eigenvalue of $P(\lambda)$ is an element $\lambda_{0} \in \mathbb{F}$ such that

$$
\operatorname{rank} P\left(\lambda_{0}\right)<\operatorname{nrank} P(\lambda)
$$

We say that $P(\lambda)$ with degree $k$ has an infinite eigenvalue if the reversal polynomial

$$
\begin{equation*}
\operatorname{rev} P(\lambda):=\lambda^{k} P(1 / \lambda)=\sum_{i=0}^{k} \lambda^{i} A_{k-i} \tag{2}
\end{equation*}
$$

has zero as an eigenvalue.
An $m \times n$ singular matrix polynomial $P(\lambda)$ may have right (column) and/or left (row) null vectors, that is, vectors $x(\lambda) \in \mathbb{F}(\lambda)^{n \times 1}$ and $y(\lambda)^{T} \in \mathbb{F}(\lambda)^{1 \times m}$ such that $P(\lambda) x(\lambda) \equiv 0$ and $y(\lambda)^{T} P(\lambda) \equiv 0$, respectively, where $y(\lambda)^{T}$ denotes the transpose of $y(\lambda)$. This leads to the following definition.

Definition 2.1. The right and left nullspaces of the $m \times n$ matrix polynomial $P(\lambda)$, denoted by $\mathcal{N}_{r}(P)$ and $\mathcal{N}_{\ell}(P)$, respectively, are the following subspaces:

$$
\begin{aligned}
\mathcal{N}_{r}(P) & :=\left\{x(\lambda) \in \mathbb{F}(\lambda)^{n \times 1}: P(\lambda) x(\lambda) \equiv 0\right\} \\
\mathcal{N}_{\ell}(P) & :=\left\{y(\lambda)^{T} \in \mathbb{F}(\lambda)^{1 \times m}: y(\lambda)^{T} P(\lambda) \equiv 0^{T}\right\}
\end{aligned}
$$

Note that the identities $\operatorname{nrank} P(\lambda)=n-\operatorname{dim} \mathcal{N}_{r}(P)=m-\operatorname{dim} \mathcal{N}_{\ell}(P)$ hold.
It is well known that the elementary divisors of $P(\lambda)$ corresponding to its finite eigenvalues, as well as the dimensions of $\mathcal{N}_{r}(P)$ and $\mathcal{N}_{\ell}(P)$, are invariant under unimodular equivalence [16], i.e., under pre- and post-multiplication of $P(\lambda)$ by unimodular matrices (square matrix polynomials with nonzero constant determinant). The elementary divisors of $P(\lambda)$ corresponding to the infinite eigenvalue are defined as the elementary divisors corresponding to the zero eigenvalue of the reversal polynomial [20, Definition 1] and may change under unimodular equivalence.

Next we define linearizations and strong linearizations of matrix polynomials.
Definition 2.2. A matrix pencil $L(\lambda)=\lambda X+Y$ is a linearization of an $m \times n$ matrix polynomial $P(\lambda)$, if for some $s \geq 0$ there exist two unimodular matrices $U(\lambda)$ and $V(\lambda)$ such that

$$
U(\lambda) L(\lambda) V(\lambda)=\left[\begin{array}{c|c}
I_{s} & 0  \tag{3}\\
\hline 0 & P(\lambda)
\end{array}\right]
$$

i.e., if $L(\lambda)$ is unimodularly equivalent to $\operatorname{diag}\left[I_{s}, P(\lambda)\right]$. A linearization $L(\lambda)$ is called a strong linearization if $\operatorname{rev} L(\lambda)$ is also a linearization of $\operatorname{rev} P(\lambda)$.

The definition of linearization was introduced in [18], while the notion of strong linearization was introduced in [17] and later named in [26]. In [17, 18, 26] only regular (square) matrix polynomials were considered. These definitions were extended to any matrix polynomial in [9], that is, including rectangular and square (regular or singular) polynomials. The original definition in [18, p. 12] for $n \times n$ regular polynomials considers linearizations with sizes $(n+s) \times(n+s)$ and $s \geq 0$ arbitrary. However, for $n \times n$ matrix polynomials with degree $k$, the definition of linearization presented in most references fixes the size of the linearizations to be $n k \times n k$, which corresponds to $s=(k-1) n$ in Definition 2.2. Perhaps the reason for this commonly encountered size restriction lies in the fact that all linearizations of a matrix polynomial with nonsingular leading coefficient have sizes at least $n k \times n k$ and that, moreover, all strong linearizations of regular matrix polynomials have size exactly $n k \times n k$ [9]. However, if $P(\lambda)$ is an $n \times n$ singular polynomial with degree $k$, then there are strong linearizations with size strictly less than $n k \times n k$ [9] that have interest in applications [6]. For these and other reasons, the size of the matrix pencil $L(\lambda)$ in Definition 2.2 is not fixed. In fact, when $P(\lambda)$ is rectangular there always exist strong linearizations for $P(\lambda)$ with different sizes. This is illustrated by the two most common linearizations used in practice, i.e., the first and second Frobenius companion forms, which for the $n \times n$ polynomial $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ are

$$
C_{1}(\lambda):=\lambda\left[\begin{array}{llll}
A_{k} & & &  \tag{4}\\
& I_{n} & & \\
& & \ddots & \\
& & & I_{n}
\end{array}\right]+\left[\begin{array}{cccc}
A_{k-1} & A_{k-2} & \cdots & A_{0} \\
-I_{n} & 0 & \cdots & 0 \\
\vdots & \ddots & \ddots & \vdots \\
0 & \cdots & -I_{n} & 0
\end{array}\right]
$$

and

$$
C_{2}(\lambda):=\lambda\left[\begin{array}{llll}
A_{k} & & &  \tag{5}\\
& I_{n} & & \\
& & \ddots & \\
& & & I_{n}
\end{array}\right]+\left[\begin{array}{cccc}
A_{k-1} & -I_{n} & \cdots & 0 \\
A_{k-2} & 0 & \ddots & \vdots \\
\vdots & \vdots & \ddots & -I_{n} \\
A_{0} & 0 & \cdots & 0
\end{array}\right]
$$

and both have size $n k \times n k$. However, if $P(\lambda)$ is rectangular with size $m \times n$, then the identity matrices in $C_{2}(\lambda)$ must have size $m \times m$. So $C_{1}(\lambda)$ has size $(m+(k-1) n) \times k n$, and $C_{2}(\lambda)$ has size $k m \times((k-1) m+n)$. These sizes are different when $m \neq n$.

It is well known that strong linearizations are relevant in the study of both regular and singular square matrix polynomials, because they are the only matrix pencils preserving the dimension of the left and right nullspaces as well as the finite and infinite elementary divisors of $P(\lambda)$ [10, Lemma 2.3]. For the rectangular case, the same result is true because the arguments used to prove this fact do not depend on $P(\lambda)$ being square or rectangular (see the proof of Lemma 2.3 in [10]). Thus for rectangular matrix polynomials we have the following analogue of Lemma 2.3 in [10].

Lemma 2.3. Let $P(\lambda)$ be an $m \times n$ matrix polynomial and let $L(\lambda)$ be an $(m+s) \times(n+s)$ matrix pencil for some $s \geq 0$, and consider the following conditions on $L(\lambda)$ and $P(\lambda)$ :
(a) $\operatorname{dim} \mathcal{N}_{r}(L)=\operatorname{dim} \mathcal{N}_{r}(P)$,
(b) $L(\lambda)$ and $P(\lambda)$ have exactly the same finite elementary divisors,
(c) $L(\lambda)$ and $P(\lambda)$ have exactly the same infinite elementary divisors.

Then $L(\lambda)$ is

- a linearization of $P(\lambda)$ if and only if conditions (a) and (b) hold,
- a strong linearization of $P(\lambda)$ if and only if conditions (a), (b) and (c) hold.

Note that condition (a) in Lemma 2.3 is equivalent to $\operatorname{dim} \mathcal{N}_{\ell}(L)=\operatorname{dim} \mathcal{N}_{\ell}(P)$.
A vector polynomial is a vector whose entries are polynomials in the variable $\lambda$. For any subspace of $\mathbb{F}(\lambda)^{n \times 1}$, it is always possible to find a basis consisting entirely of vector polynomials. The degree of a vector polynomial is the greatest degree of its components, and the order of a polynomial basis is defined as the sum of the degrees of its vectors [15, p. 494]. Then the following definition makes sense.

Definition 2.4. [15] Let $\mathcal{V}$ be a subspace of $\mathbb{F}(\lambda)^{n \times 1}$. A minimal basis of $\mathcal{V}$ is any polynomial basis of $\mathcal{V}$ with least order among all polynomial bases of $\mathcal{V}$.

It can be shown [15] that for any given subspace $\mathcal{V}$ of $\mathbb{F}(\lambda)^{n \times 1}$, the ordered list of degrees of the vector polynomials in any minimal basis of $\mathcal{V}$ is always the same. These degrees are then called the minimal indices of $\mathcal{V}$. Given a matrix polynomial $P(\lambda)$, the minimal indices and bases of the subspace $\mathcal{N}_{r}(P)$ are called the right minimal indices and bases of $P(\lambda)$, while the minimal indices and bases of $\mathcal{N}_{\ell}(P)$ are called the left minimal indices and bases of $P(\lambda)$. These magnitudes have important applications in Linear System Theory [24].

The left (right) minimal indices of a matrix pencil can be read off from the sizes of the left (right) singular blocks of the Kronecker canonical form of the pencil [16, Chap. XII]. Consequently, the minimal indices of a pencil can be stably computed via the GUPTRI form [33, 7, 8, 13]. Therefore it is natural to look for relationships between the minimal indices of a singular matrix polynomial $P(\lambda)$ and the minimal indices of a given linearization of $P(\lambda)$, since this would lead to a numerical method for computing the minimal indices of $P(\lambda)$. In the case of square singular matrix polynomials, such relationships were found in [10] for the pencils introduced in [27], in [11] for Fiedler pencils, and in [5] for generalized Fiedler pencils. In the case of Fiedler pencils of rectangular polynomials, we will develop analogous relationships in Section 5.

## 3. Fiedler pencils: definition and structural properties

In this section we first recall the notion of Fiedler pencils for square matrix polynomials, introduced in [2] and named later in [11]. In Section 3.1 we will present Algorithm 1 to construct these pencils. In Section 3.2 we extend the notion of Fiedler pencils to rectangular $m \times n$ matrix polynomials by means of Algorithm 2, which is a slight modification of Algorithm 1. This will motivate the main definition in this paper, Definition 3.8, which includes the one for the square case by just considering $n=m$. Also in Section 3.2 we will present some structural properties of Fiedler pencils that will be used later. Finally in Section 3.3 we will show the connection between the reversal of a Fiedler pencil and the reversal of the polynomial.

To introduce the Fiedler pencils of an $n \times n$ matrix polynomial $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$, we need the following block-partitioned matrices:

$$
M_{k}:=\left[\begin{array}{ll}
A_{k} &  \tag{6}\\
& I_{(k-1) n}
\end{array}\right], \quad M_{0}:=\left[\begin{array}{ll}
I_{(k-1) n} & \\
& -A_{0}
\end{array}\right]
$$

and

$$
M_{i}:=\left[\begin{array}{cccc}
I_{(k-i-1) n} & & &  \tag{7}\\
& -A_{i} & I_{n} & \\
& I_{n} & 0 & \\
& & & I_{(i-1) n}
\end{array}\right], \quad i=1, \ldots, k-1 .
$$

Notice that

$$
\begin{equation*}
M_{i} M_{j}=M_{j} M_{i} \quad \text { for }|i-j| \neq 1 \tag{8}
\end{equation*}
$$

Now, we introduce Fiedler pencils in the same way as in [11].
Definition 3.1 (Fiedler Pencils for square matrix polynomials). Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $n \times n$ matrix polynomial and let $M_{i}, i=0,1, \ldots, k$, be the matrices defined in (6) and (7). Given any bijection $\sigma:\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$, the Fiedler pencil of $P(\lambda)$ associated with $\sigma$ is the nk $\times$ nk matrix pencil

$$
\begin{equation*}
F_{\sigma}(\lambda):=\lambda M_{k}-M_{\sigma^{-1}(1)} \cdots M_{\sigma^{-1}(k)} . \tag{9}
\end{equation*}
$$

Note that $\sigma(i)$ describes the position of the factor $M_{i}$ in the product $M_{\sigma^{-1}(1)} \cdots M_{\sigma^{-1}(k)}$ defining the zero-degree term in (9): i.e., $\sigma(i)=j$ means that $M_{i}$ is the $j$ th factor in the product. For brevity, we denote this product by

$$
\begin{equation*}
M_{\sigma}:=M_{\sigma^{-1}(1)} \cdots M_{\sigma^{-1}(k)} \tag{10}
\end{equation*}
$$

so that $F_{\sigma}(\lambda):=\lambda M_{k}-M_{\sigma}$.

As in [11], sometimes we will write the bijection $\sigma$ using the array notation $\sigma=(\sigma(0), \sigma(1), \ldots, \sigma(k-1))$. Unless otherwise stated, the matrices $M_{i}, i=0, \ldots, k, M_{\sigma}$, and the Fiedler pencil $F_{\sigma}(\lambda)$ refer to the matrix polynomial $P(\lambda)$ in (1). When necessary, we will explicitly indicate the dependence on a certain polynomial $Q(\lambda)$ by writing $M_{i}(Q), M_{\sigma}(Q)$ and $F_{\sigma}(Q)$. In this situation, the dependence on $\lambda$ is dropped in the Fiedler pencil $F_{\sigma}(Q)$ for simplicity. Since matrix polynomials will always be denoted by capital letters, there is no risk of confusion between $F_{\sigma}(\lambda)$ and $F_{\sigma}(Q)$.

The set of Fiedler pencils includes the first and second companion forms [18, 11]. More precisely, the first companion form corresponds to the bijection $\sigma_{1}=(k, k-1, \ldots, 2,1)$ and the second to the bijection $\sigma_{2}=(1,2, \ldots, k-1, k)$. Other relevant Fiedler pencils are the pentadiagonal Fiedler pencils that are described in detail in [11, Example 3.2].

It is shown in [11] that the relative positions of the matrices $M_{i}$ and $M_{i+1}$, for $i=0,1, \ldots, k-2$, in the product $M_{\sigma}$ determine most of the relevant properties of the Fiedler pencil $F_{\sigma}(\lambda)$. This motivates Definition 3.2, that was introduced in [11, Definition 3.3].

Definition 3.2. Let $\sigma:\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$ be a bijection.
(a) For $i=0, \ldots, k-2$, we say that $\sigma$ has a consecution at $i$ if $\sigma(i)<\sigma(i+1)$ and that $\sigma$ has an inversion at $i$ if $\sigma(i)>\sigma(i+1)$.
(b) We denote by $\mathfrak{c}(\sigma)$ the total number of consecutions in $\sigma$, and by $\mathfrak{i}(\sigma)$ the total number of inversions in $\sigma$.
(c) For $i \leq j$, we denote by $\mathfrak{c}(\sigma(i: j))$ the total number of consecutions that $\sigma$ has at $i, i+1, \ldots, j$, and by $\mathfrak{i}(\sigma(i: j))$ the total number of inversions that $\sigma$ has at $i, i+1, \ldots, j$. Observe that $\mathfrak{c}(\sigma)=\mathfrak{c}(\sigma(0$ : $k-2)$ ) and $\mathfrak{i}(\sigma)=\mathfrak{i}(\sigma(0: k-2))$.
(d) The consecution-inversion structure sequence of $\sigma$, denoted by $\operatorname{CISS}(\sigma)$, is the tuple $\left(c_{1}, i_{1}, c_{2}, i_{2}, \ldots\right.$, $\left.c_{\ell}, i_{\ell}\right)$, where $\sigma$ has $c_{1}$ consecutive consecutions at $0,1, \ldots, c_{1}-1 ; i_{1}$ consecutive inversions at $c_{1}, c_{1}+1, \ldots, c_{1}+i_{1}-1$ and so on, up to $i_{\ell}$ inversions at $k-1-i_{\ell}, \ldots, k-2$.

We want to point out that, though the notions introduced in Definition 3.2 depend only on the bijection $\sigma$ and not on the Fiedler pencil $F_{\sigma}(\lambda)$, they are closely related to the definition of $F_{\sigma}(\lambda)$, as it is shown in the following remark.

Remark 3.3. The following simple observations on Definition 3.2 will be used freely.

1. $\sigma$ has a consecution at $i$ if and only if $M_{i}$ is to the left of $M_{i+1}$ in $M_{\sigma}$, while $\sigma$ has an inversion at $i$ if and only if $M_{i}$ is to the right of $M_{i+1}$ in $M_{\sigma}$.
2. Either $c_{1}$ or $i_{\ell}$ in $\operatorname{CISS}(\sigma)$ may be zero (in the first case $\sigma$ has an inversion at 0 , in the second it has a consecution at $k-2$ ), but $i_{1}, c_{2}, i_{2}, \ldots, i_{\ell-1}, c_{\ell}$ are all strictly positive. These conditions uniquely determine $\operatorname{CISS}(\sigma)$ and, in particular, the parameter $\ell$.
3. $\mathfrak{c}(\sigma)=\sum_{j=1}^{\ell} c_{j}, \mathfrak{i}(\sigma)=\sum_{j=1}^{\ell} i_{j}$, and $\mathfrak{c}(\sigma)+\mathfrak{i}(\sigma)=k-1$.

The reader may find in [11, Example 3.5] explicit examples of CISS $(\sigma)$ for some relevant Fiedler pencils.

### 3.1. A multiplication free algorithm to construct Fiedler pencils of square matrix polynomials

We focus only on how to construct the zero-degree term $M_{\sigma}$ in the Fiedler pencil (9), since the firstdegree term is already known. The obvious option is to perform directly the multiplication of all factors, but this is not convenient if the degree is large. ${ }^{1}$ Theorem 3.4 shows how to construct Fiedler pencils without performing any arithmetic operation. Throughout this paper, we will use MATLAB notation for submatrices on block indices, that is, if $A$ is a matrix partioned into blocks, then $A(i: j,:)$ indicates the submatrix of $A$ consisting of block rows $i$ through $j$ and $A(:, k: l)$ indicates the submatrix of $A$ consisting of block columns $k$ through $l$.

[^1]Theorem 3.4. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $n \times n$ matrix polynomial with degree $k \geq 2$, let $\sigma$ : $\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$ be a bijection, and let $M_{\sigma}$ be the zero-degree term of the Fiedler pencil of $P(\lambda)$ associated with $\sigma$. Consider the matrices $W_{0}, W_{1}, \ldots, W_{k-2}$ constructed by Algorithm 1 below partitioned, respectively, into $2 \times 2,3 \times 3, \ldots, k \times k$ blocks of size $n \times n$. Then Algorithm 1 constructs $M_{\sigma}$, more precisely, $M_{\sigma}=W_{k-2}$.
Algorithm 1. Given $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ with size $n \times n$ and a bijection $\sigma$, the following algorithm constructs $M_{\sigma}$.
if $\sigma$ has a consecution at 0 then

$$
W_{0}=\left[\begin{array}{cc}
-A_{1} & I_{n} \\
-A_{0} & 0
\end{array}\right]
$$

else

$$
W_{0}=\left[\begin{array}{cc}
-A_{1} & -A_{0} \\
I_{n} & 0
\end{array}\right]
$$

endif
for $i=1: k-2$
if $\sigma$ has a consecution at $i$ then

$$
W_{i}=\left[\begin{array}{ccc}
-A_{i+1} & I_{n} & 0 \\
W_{i-1}(:, 1) & 0 & W_{i-1}(:, 2: i+1)
\end{array}\right]
$$

else

$$
W_{i}=\left[\begin{array}{cc}
-A_{i+1} & W_{i-1}(1,:) \\
I_{n} & 0 \\
0 & W_{i-1}(2: i+1,:)
\end{array}\right]
$$

endif
endfor
$M_{\sigma}=W_{k-2}$
Proof. The proof proceeds by induction on the degree $k$. The result is obvious for $k=2$, because in this case there are only two possible options for $M_{\sigma}$, namely, $M_{\sigma}=M_{0} M_{1}$ if $\sigma$ has a consecution at 0 or $M_{\sigma}=M_{1} M_{0}$ if $\sigma$ has an inversion at 0 . A direct computation shows that

$$
M_{0} M_{1}=\left[\begin{array}{cc}
-A_{1} & I_{n}  \tag{11}\\
-A_{0} & 0
\end{array}\right] \quad \text { and } \quad M_{1} M_{0}=\left[\begin{array}{cc}
-A_{1} & -A_{0} \\
I_{n} & 0
\end{array}\right] \quad \text { for } k=2,
$$

and the result follows.
Assume now that the result is valid for matrix polynomials of degree $k-1 \geq 2$, and let us prove it for the polynomial $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ and the bijection $\sigma:\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$. Note first that the matrices $M_{i}(P)$ defined in (6-7) for $P(\lambda)$ satisfy

$$
\begin{equation*}
M_{i}(P)=\operatorname{diag}\left(I_{n}, M_{i}(Q)\right), \quad \text { for } i=0, \ldots, k-2, \tag{12}
\end{equation*}
$$

where $M_{i}(Q)$ are the $n(k-1) \times n(k-1)$ matrices corresponding to the polynomial $Q(\lambda)=\sum_{i=0}^{k-1} \lambda^{i} A_{i}$. We need to distinguish two cases in the proof.

Case 1. If $\sigma$ has a consecution at $k-2$, then the commutativity relations (8) of the $M_{i}$ 's matrices allow us to write

$$
M_{\sigma}(P)=M_{i_{0}}(P) \cdots M_{i_{k-2}}(P) M_{k-1}(P)
$$

where $\left(i_{0}, i_{1}, \ldots, i_{k-2}\right)$ is a permutation of $(0,1, \ldots, k-2)$. By using (12), we can write

$$
\begin{equation*}
M_{\sigma}(P)=\operatorname{diag}\left(I_{n}, M_{\widetilde{\sigma}}(Q)\right) M_{k-1}(P) \tag{13}
\end{equation*}
$$

where $\widetilde{\sigma}:\{0,1, \ldots, k-2\} \rightarrow\{1, \ldots, k-1\}$ is a bijection such that, for $i=0, \ldots, k-3, \widetilde{\sigma}$ has a consecution (resp. inversion) at $i$ if and only if $\sigma$ has a consecution (resp. inversion) at $i$. Therefore, by the induction hypothesis, $M_{\widetilde{\sigma}}(Q)=W_{k-3}$. Finally, we perform the simple block product in (13) as follows

$$
\begin{aligned}
M_{\sigma}(P) & =\left[\begin{array}{ccc}
I_{n} & 0_{n} & 0 \\
0 & W_{k-3}(:, 1) & W_{k-3}(:, 2: k-1)
\end{array}\right]\left[\begin{array}{ccc}
-A_{k-1} & I_{n} & \\
I_{n} & 0_{n} & \\
& & I_{(k-2) n}
\end{array}\right] \\
& =\left[\begin{array}{ccc}
-A_{k-1} & I_{n} & 0 \\
W_{k-3}(:, 1) & 0 & W_{k-3}(:, 2: k-1)
\end{array}\right]
\end{aligned}
$$

which is precisely the matrix $W_{k-2}$ constructed by Algorithm 1 when $\sigma$ has a consecution at $k-2$.
Case 2. If $\sigma$ has an inversion at $k-2$ the proof is similar, but with $M_{k-1}(P)$ placed on the left, i.e.,

$$
M_{\sigma}(P)=M_{k-1}(P) M_{i_{0}}(P) \cdots M_{i_{k-2}}(P)=M_{k-1}(P) \operatorname{diag}\left(I_{n}, M_{\widetilde{\sigma}}(Q)\right)
$$

### 3.2. Fiedler pencils of rectangular matrix polynomials

The extension of equation (9) to a rectangular $m \times n$ matrix polynomial $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ presents difficulties, because it is not clear how to define the sizes of the identity blocks in the main block diagonal of the factors $M_{i}$. A tentative approach is simply to choose the sizes of the diagonal identities in both (6) and (7) such that all the factors in (10) are conformal for multiplication (notice that the non-diagonal identities in the central $2 \times 2$ block submatrix of (7) are determined by the size of $A_{i} \in \mathbb{F}^{m \times n}$ ). This can be done, but it is not immediate and is cumbersome, because the presence of the block $-A_{i}$ in the matrix $M_{i}$ imposes restrictions on the sizes of the diagonal identities of the factors to both the left and the right of $M_{i}$ in the product defining $M_{\sigma}$. Hence, proceeding in this way, the sizes of the matrices $M_{i}$ should be carefully determined and, even more, these sizes would depend, for each Fiedler pencil, on the position of the $M_{i}$ factor in the product defining $M_{\sigma}$. These questions are better explained with an example.

Example 3.5. Let $P(\lambda)=A_{0}+\lambda A_{1}+\lambda^{2} A_{2}+\lambda^{3} A_{3}$, with $A_{i} \in \mathbb{F}^{m \times n}$, be a matrix polynomial with degree 3 and $\sigma_{1}=(1,3,2)$ and $\sigma_{2}=(2,3,1)$ be bijections of $\{0,1,2\}$ on $\{1,2,3\}$. Let us try to give a meaning to the symbolic expressions

$$
F_{\sigma_{1}}(\lambda)=\lambda M_{3}-M_{0} M_{2} M_{1} \quad \text { and } \quad F_{\sigma_{2}}(\lambda)=\lambda M_{3}-M_{2}^{\prime} M_{0}^{\prime} M_{1}^{\prime}
$$

that is, let us try to define the Fiedler pencils for $P(\lambda)$ associated with the bijections $\sigma_{1}$ and $\sigma_{2}$. When $P(\lambda)$ is square $(n=m)$ the commutativity relations (8) imply that $F_{\sigma_{1}}(\lambda)=F_{\sigma_{2}}(\lambda)$. Assume now that $m \neq n$, then the factors in the zero degree term of $F_{\sigma_{1}}(\lambda)$ are conformal for multiplication if and only if they are

$$
M_{0}=\left[\begin{array}{ccc}
I_{m} & 0 & 0 \\
0 & I_{n} & 0 \\
0 & 0 & -A_{0}
\end{array}\right], \quad M_{1}=\left[\begin{array}{ccc}
I_{n} & 0 & 0 \\
0 & -A_{1} & I_{m} \\
0 & I_{n} & 0
\end{array}\right], \quad M_{2}=\left[\begin{array}{ccc}
-A_{2} & I_{m} & 0 \\
I_{n} & 0 & 0 \\
0 & 0 & I_{n}
\end{array}\right]
$$

and the factors in the zero degree term of $F_{\sigma_{2}}(\lambda)$ are conformal for multiplication if and only if they are

$$
M_{0}^{\prime}=\left[\begin{array}{ccc}
I_{n} & 0 & 0 \\
0 & I_{m} & 0 \\
0 & 0 & -A_{0}
\end{array}\right], \quad M_{1}^{\prime}=\left[\begin{array}{ccc}
I_{n} & 0 & 0 \\
0 & -A_{1} & I_{m} \\
0 & I_{n} & 0
\end{array}\right], \quad M_{2}^{\prime}=\left[\begin{array}{ccc}
-A_{2} & I_{m} & 0 \\
I_{n} & 0 & 0 \\
0 & 0 & I_{m}
\end{array}\right]
$$

Note that the size of $M_{2}$ is different than the size of $M_{2}^{\prime}$. However, the reader is invited to check that $F_{\sigma_{1}}(\lambda)=F_{\sigma_{2}}(\lambda)$. This example shows that defining Fiedler pencils for rectangular polynomials in a similar way as in the square case would force the sizes of the $M_{i}$ matrices to depend on the specific bijection $\sigma$. It is easy to devise examples of rectangular matrix polynomials of degree higher than 3 where the sizes are different for more than one factor $M_{i}$.

A first option to extend Fiedler pencils from square to rectangular polynomials that is not affected by the difficulties illustrated in Example 3.5 would be the following. Use, in the square case, the commutativity relations (8) to order the factors $M_{i}$ in $M_{\sigma}(10)$ in a certain canonical order that is exactly the same for all Fiedler pencils with the same $\operatorname{CISS}(\sigma)$ (that are, in fact, the same pencil by Theorem $3.4)$. One possible order may be found in [2, eq. (2.9)]. Then, use this order and force the conformability of all $M_{i}$ factors for multiplication, by choosing properly the sizes of their identity blocks, to extend the Fiedler pencil to rectangular matrix polynomials. Again, this can be done, but it requires to prove, for each different CISS $(\sigma)$, that the sizes of the $M_{i}$ factors can always be properly chosen and to determine these sizes. This is not obvious and is tedious. In addition, the reader may easily check that the sizes of the $M_{i}$ factors may be different for different $\operatorname{CISS}(\sigma)$, what is still unpleasant.

Another option to extend Fiedler pencils from square to rectangular polynomials bypassing all difficulties mentioned above is to avoid the use of the factors $M_{i}$. To this purpose, we might start by performing symbolically in the square case the product defining $M_{\sigma}$ in (10), in order to obtain an explicit expression of the block-entries of $M_{\sigma}$ in terms of the coefficients $A_{i}$ of the polynomial $P(\lambda)$. This can be done by using $\operatorname{CISS}(\sigma)$, although is complicated and requires a cumbersome notation. Once this explicit expression is known, we would replace the square $n \times n$ blocks $A_{i}, i=0,1, \ldots, k-1$, by rectangular $m \times n$ blocks $A_{i}$, and we would check that the sizes of the block rows and block columns fit properly with assigning either a size $n \times n$ or $m \times m$ to every identity block that appears in $M_{\sigma}$. This requires again a tedious proof. Therefore, we will follow a simpler approach based on adapting Algorithm 1 to rectangular matrix polynomials. This approach is developed in Theorem 3.6 and Definition 3.8 and is, in fact, equivalent to the process described above of obtaining an explicit expression of the block-entries of $M_{\sigma}$ in terms of the coefficients $A_{i}$. Note that in Algorithm 2 we use again MATLAB notation for submatrices on block indices.

Theorem 3.6. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$ and let $\sigma$ : $\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$ be a bijection. If in Algorithm 2 below each matrix $W_{i}$, for $i=1,2, \ldots, k-$ 2, is partitioned into blocks in such a way that the blocks of $W_{i-1}$ are blocks of $W_{i}$, then Algorithm 2 constructs a sequence $\left\{W_{0}, W_{1}, \ldots, W_{k-2}\right\}$ of matrices partitioned in $2 \times 2,3 \times 3, \ldots, k \times k$ blocks, respectively.
Algorithm 2. Given $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ with size $m \times n$ and a bijection $\sigma$, the following algorithm constructs a sequence of matrices $\left\{W_{0}, W_{1}, \ldots, W_{k-2}\right\}$.
if $\sigma$ has a consecution at 0 then

$$
W_{0}=\left[\begin{array}{cc}
-A_{1} & I_{m} \\
-A_{0} & 0
\end{array}\right]
$$

else

$$
W_{0}=\left[\begin{array}{cc}
-A_{1} & -A_{0} \\
I_{n} & 0
\end{array}\right]
$$

endif
for $i=1: k-2$
if $\sigma$ has a consecution at $i$ then

$$
W_{i}=\left[\begin{array}{ccc}
-A_{i+1} & I_{m} & 0 \\
W_{i-1}(:, 1) & 0 & W_{i-1}(:, 2: i+1)
\end{array}\right]
$$

else

$$
W_{i}=\left[\begin{array}{cc}
-A_{i+1} & W_{i-1}(1,:) \\
I_{n} & 0 \\
0 & W_{i-1}(2: i+1,:)
\end{array}\right]
$$

endif
endfor
In addition, the matrices $\left\{W_{0}, W_{1}, \ldots, W_{k-2}\right\}$ satisfy the following properties:
(a) The size of $W_{i}$ is

$$
(m+m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0: i))) \times(n+m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0: i)))
$$

(b) The first diagonal block of $W_{i}$ is $-A_{i+1}$ and, so, has size $m \times n$. The rest of diagonal blocks of $W_{i}$ are square zero matrices, more precisely

$$
W_{i}(i+2-j, i+2-j)=\left\{\begin{array}{ll}
0_{m} & \text { if } \sigma \text { has a consecution at } j \\
0_{n} & \text { if } \sigma \text { has an inversion at } j
\end{array} \quad, \quad \text { for } j=0,1, \ldots, i\right.
$$

Proof. The proof is elementary. We simply sketch the main points. First, notice that the matrix $W_{0}$ is well defined in Algorithm 2. Therefore, $W_{1}$ is also well defined both when $\sigma$ has a consecution at 1 and has an inversion at 1 , because in both cases $W_{1}(1,1)=-A_{2}, W_{0}(:, 1)$ has $n$ columns, and $W_{0}(1,:)$ has $m$ rows. The same argument can be applied inductively to show that also $W_{2}, \ldots, W_{k-2}$ are well defined.

The fact that $W_{i}$ is partitioned into $(i+2) \times(i+2)$ blocks is true by definition for $W_{0}$, and for the rest of matrices in the sequence it follows from the fact that one block row and one block column are added in each step of the "for" loop of Algorithm 2. Part (a) is again true for $W_{0}$, and for obtaining the result for the rest of matrices in the sequence note that: (1) if $\sigma$ has a consecution at $i$, then $W_{i}$ has $m$ rows and $m$ columns more than $W_{i-1} ;(2)$ if $\sigma$ has an inversion at $i$, then $W_{i}$ has $n$ rows and $n$ columns more than $W_{i-1}$. Finally, let us prove part (b). The result is true for $W_{0}$. For the rest of matrices in the sequence assume that it is true for $W_{i-1}$ and let us prove it for $W_{i}$. Note that by construction $W_{i}(1,1)=-A_{i+1}$ and

$$
W_{i}(2,2)= \begin{cases}0_{m} & \text { if } \sigma \text { has a consecution at } i \\ 0_{n} & \text { if } \sigma \text { has an inversion at } i\end{cases}
$$

which is Part (b) for $j=i$. Observe also that

$$
W_{i}(3: i+2,3: i+2)=W_{i-1}(2: i+1,2: i+1)
$$

which implies $W_{i}(i+2-j, i+2-j)=W_{i-1}((i-1)+2-j,(i-1)+2-j)$ for $j=0,1, \ldots, i-1$. This proves the result since we are assuming that the result is true for $W_{i-1}$.
Remark 3.7. In part (b) of Theorem 3.6 we assume, as in the rest of the paper, that the block indices of $W_{i}$ run from 1 to $i+2$. Thus, the diagonal blocks of $W_{i}$ are $W_{i}(1,1), \ldots, W_{i}(i+2, i+2)$. If we let the block indices of $W_{i}$ run from $k-i-1$ to $k$, the result in part (b) is expressed as

$$
W_{i}(k-j, k-j)=\left\{\begin{array}{ll}
0_{m} & \text { if } \sigma \text { has a consecution at } j \\
0_{n} & \text { if } \sigma \text { has an inversion at } j
\end{array}, \quad \text { for } j=0,1, \ldots, i\right.
$$

which shows that the sizes of these blocks only depend on $j$ and not on $i$, as long as $0 \leq j \leq i$.
Observe that Algorithm 2 differs from Algorithm 1 only in the sizes of the identity blocks, that are chosen to fit the size $m \times n$ of the polynomial $P(\lambda)$. This fact and Theorem 3.4 motivate Definition 3.8, which is the main definition in this paper.

Definition 3.8 (Fiedler Pencils for rectangular matrix polynomials). Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, let $\sigma:\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$ be a bijection, and denote by $M_{\sigma}$ the last matrix of the sequence constructed by Algorithm 2, that is,

$$
M_{\sigma}:=W_{k-2} .
$$

The Fiedler pencil of $P(\lambda)$ associated with $\sigma$ is the $(m+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)) \times(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma))$ matrix pencil

$$
F_{\sigma}(\lambda):=\lambda\left[\begin{array}{ll}
A_{k} &  \tag{14}\\
& I_{m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)}
\end{array}\right]-M_{\sigma} .
$$

Remark 3.9. Some remarks on Definition 3.8 may be useful for the reader.

1. The leading coefficient $\left[{ }_{A_{k}}{ }_{I}\right]$ of the Fiedler pencil $F_{\sigma}(\lambda)$ introduced in Definition 3.8 has the same structure as the matrix $M_{k}$ in (6), but the size of the block diagonal identity is different when $m \neq n$.
2. If $m \neq n$, then there are Fiedler pencils with different sizes, because the sum $\mathfrak{c}(\sigma)+\mathfrak{i}(\sigma)=k-1$ is fixed for all $\sigma$ and, so, different pairs of $(\mathfrak{c}(\sigma), \mathfrak{i}(\sigma))$ produce different sizes of $F_{\sigma}(\lambda)$. For instance, if $m>n$, then the Fiedler pencil with smallest size corresponds to $\mathfrak{c}(\sigma)=0$, i.e., to the first companion form, and the one with largest size corresponds to $\mathfrak{i}(\sigma)=0$, i.e., to the second companion form [11]. If $n>m$, then the opposite situation holds.
3. We use in Theorem 3.6 and Definition 3.8 a bijection $\sigma$ for the only purpose of keeping a parallelism with the standard definition of Fiedler pencils for square polynomials. However, a bijection is not really needed since we do not use factors $M_{i}$ in our definition. Observe that Algorithm 2 only needs a sequence of decisions that we have identified with $\sigma$ having a consecution or inversion.
4. A comparison between Algorithms 1 and 2 indicates that, for the same bijection $\sigma$, Fiedler pencils of square and rectangular matrix polynomials look symbolically the same except for the sizes of the identity blocks. This means that a fundamental consequence of Theorems 3.4 and 3.6 is that the zero degree term $M_{\sigma}$ of every Fiedler pencil of a square $n \times n$ matrix polynomial $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ satisfies that the identity blocks of $M_{\sigma}$ are placed in block-entries such that if the polynomial becomes rectangular, i.e., with size $m \times n$, then the square identity blocks may be transformed consistently into $I_{m}$ or $I_{n}$ matrices. Let us see an specific example. Consider $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ with degree $k=6$ and size $n \times n$, and the bijection $\tau=(1,2,5,3,6,4)$. In this case

$$
M_{\tau}=M_{0} M_{1} M_{3} M_{5} M_{2} M_{4}=\left[\begin{array}{cccccc}
-A_{5} & -A_{4} & I_{n} & 0 & 0 & 0  \tag{15}\\
I_{n} & 0 & 0 & 0 & 0 & 0 \\
0 & -A_{3} & 0 & -A_{2} & I_{n} & 0 \\
0 & I_{n} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & -A_{1} & 0 & I_{n} \\
0 & 0 & 0 & -A_{0} & 0 & 0
\end{array}\right]
$$

which can be constructed by direct multiplication of the factors or via Algorithm 1, since $\tau$ has consecutions at $0,1,3$ and inversions at 2,4 . If the size becomes $m \times n$, then Algorithm 2 produces

$$
M_{\tau}=\left[\begin{array}{cccccc}
-A_{5} & -A_{4} & I_{m} & 0 & 0 & 0  \tag{16}\\
I_{n} & 0 & 0 & 0 & 0 & 0 \\
0 & -A_{3} & 0 & -A_{2} & I_{m} & 0 \\
0 & I_{n} & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & -A_{1} & 0 & I_{m} \\
0 & 0 & 0 & -A_{0} & 0 & 0
\end{array}\right]
$$

which is nothing else that (15) but modifying the sizes of some identity blocks according to the size $m \times n$ of the coefficients $-A_{i}$.

Theorem 3.10 is a direct consequence of Theorem 3.6 and establishes that the zero-degree term $M_{\sigma}$ of any Fiedler pencil of $P(\lambda)$ has as non-zero blocks $-A_{0},-A_{1}, \ldots,-A_{k-1}$ and $(k-1)$ identities of size $n \times n$ or $m \times m$. This property is very well known in the case of the first and second companion forms, that are particular cases of Fiedler pencils. Theorem 3.10 also includes additional information on the structure of $M_{\sigma}$ that will be used later.

Theorem 3.10. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$ and let $\sigma:\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$ be a bijection. Suppose that $F_{\sigma}(\lambda)=\lambda\left[{ }^{A_{k}}{ }_{I}\right]-M_{\sigma}$ is the Fiedler pencil of $P(\lambda)$ associated with $\sigma$ and consider $M_{\sigma}$ partioned into $k \times k$ blocks according to Algorithm 2. Then:
(a) $M_{\sigma}$ has $k$ blocks equal to $-A_{0},-A_{1}, \ldots,-A_{k-1}$.
(b) $M_{\sigma}$ has $k-1$ identity blocks: $\mathfrak{c}(\sigma)$ blocks equal to $I_{m}$ and $\mathfrak{i}(\sigma)$ blocks equal to $I_{n}$.
(c) The rest of the blocks of $M_{\sigma}$ are equal to 0 matrices of size $n \times n, m \times m, n \times m$, or $m \times n$.
(d) The $k-1$ identity blocks in part (b) satisfy:

1. None of them is in the main block diagonal of $M_{\sigma}$.
2. Two of these blocks are never in the same block row (or in the same block column) of $M_{\sigma}$.
3. If an identity block is in the $(i, j)$ block-entry of $M_{\sigma}$, then one and only one of the following two properties holds: (a) the rest of the blocks in the $i$ th block row of $M_{\sigma}$ are equal to 0 and at least one of the matrices $-A_{0},-A_{1}, \ldots,-A_{k-1}$ is in the $j$ th block column of $M_{\sigma}$; (b) the rest of the blocks in the $j$ th block column of $M_{\sigma}$ are equal to 0 and at least one of the matrices $-A_{0},-A_{1}, \ldots,-A_{k-1}$ is in the $i$ th block row of $M_{\sigma}$.
4. If $\left\{i_{1}, i_{2}, \ldots, i_{t}\right\}$ (resp. $\left\{j_{t+1}, j_{t+2}, \ldots, j_{k-1}\right\}$ ) are the block indices of the block rows (resp. block columns) of $M_{\sigma}$ containing one identity block and having the remaining blocks equal to zero, then the (unordered) set $\left\{i_{1}, i_{2}, \ldots, i_{t}, j_{t+1}, j_{t+2}, \ldots, j_{k-1}\right\}$ is equal to $\{2,3, \ldots, k\}$.

Proof. Parts (a), (b), and (c) are obvious from Algorithm 2. Part (d)-1 follows from Theorem 3.6-(b). The proofs of Parts (d)-2, (d)-3, and (d)-4 proceed by induction on the matrices $W_{0}, \ldots, W_{k-2}\left(=M_{\sigma}\right)$ constructed by Algorithm 2. A direct inspection shows that Parts (d)-2, (d)-3, and (d)-4 hold for $W_{0}$ with $k=2$. Let us assume that they hold for $W_{k-3}$ with $k-1$ instead of $k$. Next partition $W_{k-3}$ as follows

$$
W_{k-3}=\left[\begin{array}{cc}
-A_{k-2} & Z_{12} \\
Z_{21} & Z_{22}
\end{array}\right] .
$$

Then Algorithm 2 gives for $W_{k-2}=M_{\sigma}$

$$
W_{k-2}=\left[\begin{array}{ccc}
-A_{k-1} & I_{m} & 0  \tag{17}\\
-A_{k-2} & 0 & Z_{12} \\
Z_{21} & 0 & Z_{22}
\end{array}\right] \quad \text { or } \quad W_{k-2}=\left[\begin{array}{ccc}
-A_{k-1} & -A_{k-2} & Z_{12} \\
I_{n} & 0 & 0 \\
0 & Z_{21} & Z_{22}
\end{array}\right]
$$

and observe that the main block diagonal of $Z_{22}$ is on the main block diagonal of $W_{k-2}$. The structure of $M_{\sigma}$ in (17) and the fact that $W_{k-3}$ satisfies Parts (d)-2 and (d)-3 make evident that $M_{\sigma}$ also satisfies Parts (d)-2 and (d)-3. The block indices of the identity blocks of $W_{k-3}$ in Part (d)-4 are $\{2,3, \ldots, k-1\}$ and observe that the induction hypothesis implies that if an identity block is a block-entry of $Z_{12}$ (resp. $Z_{21}$ ) then the corresponding block column (resp. block row) in $Z_{22}$ is zero. This fact and the structure of $M_{\sigma}$ in (17) imply that the block indices in Part (d)-(4) of the identity blocks of $W_{k-3}$ as block entries of $M_{\sigma}$ are $\{3,4, \ldots, k\}$. Finally, note that the identity block that is added to construct $W_{k-2}$ from $W_{k-3}$ has always index 2 in the set of indices in Part (d)-4.

### 3.3. The reversal of a Fiedler pencil

The main result in this section is Theorem 3.14, which establishes that for a rectangular matrix polynomial $P(\lambda)$ the reversal of any of its Fiedler pencils is strictly equivalent to a Fiedler pencil of $\operatorname{rev} P(\lambda)$. We think that this result is interesting on its own right and, in addition, it will be used in Section 4 to prove that every Fiedler pencil of a rectangular matrix polynomial $P(\lambda)$ is a strong linearization of $P(\lambda)$. The proof of Theorem 3.14 is long and can be skipped in a first reading. The proof is based on the technical Lemmas 3.11, 3.12, and 3.13 that are presented below.

Lemma 3.11. Let $P(\lambda)$ be an $m \times n$ matrix polynomial with degree $k \geq 2$ and let $F_{\sigma}(P)$ be the Fiedler pencil of $P(\lambda)$ associated with the bijection $\sigma$. Then the Fiedler pencil $F_{\sigma}(-P)$ of $-P(\lambda)$ is strictly equivalent to $F_{\sigma}(P)$.
Proof. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$. Along this proof, we consider that $F_{\sigma}(P)=\lambda \operatorname{diag}\left(A_{k}, I\right)-M_{\sigma}(P)$ and $F_{\sigma}(-P)=\lambda \operatorname{diag}\left(-A_{k}, I\right)-M_{\sigma}(-P)$ are partitioned into $k \times k$ blocks with the sizes of the blocks determined by the way Algorithm 2 constructs $M_{\sigma}(P)$ and $M_{\sigma}(-P)$. In particular, we consider the block $I$ in $\operatorname{diag}\left(A_{k}, I\right)$ and $\operatorname{diag}\left(-A_{k}, I\right)$ as $I=\operatorname{diag}\left(I_{r_{2}}, I_{r_{3}}, \ldots, I_{r_{k}}\right)$, where $r_{i}=m$ or $n$ by Theorem 3.6-Part (b). Note, in the first place, that $F_{\sigma}(-P)$ is strictly equivalent to $-F_{\sigma}(-P)$. On the other hand, according to Algorithm 2 and Theorem 3.10, the only difference between the pencils $-F_{\sigma}(-P)=$ $\lambda \operatorname{diag}\left(A_{k},-I\right)-\left(-M_{\sigma}(-P)\right)$ and $F_{\sigma}(P)=\lambda \operatorname{diag}\left(A_{k}, I\right)-M_{\sigma}(P)$ are the signs of the $k-1$ identity blocks of $M_{\sigma}(P)$ and the signs of the $k-1$ diagonal identity blocks of $\operatorname{diag}\left(A_{k}, I\right)$. Let $\left\{i_{1}, i_{2}, \ldots, i_{t}\right\}$ and $\left\{j_{t+1}, j_{t+2}, \ldots, j_{k-1}\right\}$ be the indices defined in Theorem 3.10-Part (d)-4 and define now the matrices

$$
U:=\operatorname{diag}\left(I_{m}, \eta_{2} I_{r_{2}}, \eta_{3} I_{r_{3}}, \ldots, \eta_{k} I_{r_{k}}\right), \text { where } \eta_{i}= \begin{cases}-1 & \text { if } \eta_{i} \in\left\{i_{1}, i_{2}, \ldots, i_{t}\right\} \\ 1 & \text { otherwise }\end{cases}
$$

and

$$
V:=\operatorname{diag}\left(I_{n}, \alpha_{2} I_{r_{2}}, \alpha_{3} I_{r_{3}}, \ldots, \alpha_{k} I_{r_{k}}\right), \quad \text { where } \alpha_{i}= \begin{cases}-1 & \text { if } \alpha_{i} \in\left\{j_{t+1}, j_{t+2}, \ldots, j_{k-1}\right\} \\ 1 & \text { otherwise }\end{cases}
$$

According to Theorem 3.10-Part (d)-4 and the previous discussion

$$
U F_{\sigma}(P) V=\lambda U \operatorname{diag}\left(A_{k}, I\right) V-U M_{\sigma}(P) V=\lambda \operatorname{diag}\left(A_{k},-I\right)-\left(-M_{\sigma}(-P)\right)=-F_{\sigma}(-P)
$$

which concludes the proof.

Fiedler pencils for rev $P(\lambda)$ can be easily constructed by applying Algorithm 2 to the reversal polynomial. Lemma 3.12 shows us another way to construct Fiedler pencils for rev $P(\lambda)$ that is useful to prove Theorem 3.14. According to Definition 3.8, we only need to pay attention in Lemma 3.12 to the construction of the zero-degree term of the pencil. In addition, for technical reasons that will be clear later, we construct pencils for the polynomial $-\operatorname{rev} P(\lambda)$.
Lemma 3.12 (Construction of Fiedler pencils for $-\operatorname{rev} P(\lambda)$ via block reverse identities). Let $P(\lambda)=$ $\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$ and let $\sigma:\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$ be a bijection. If in Algorithm 3 below each matrix $Y_{i}$, for $i=1,2, \ldots, k-2$, is partitioned into blocks in such a way that the blocks of $Y_{i-1}$ are blocks of $Y_{i}$, then Algorithm 3 constructs a sequence $\left\{Y_{0}, Y_{1}, \ldots, Y_{k-2}\right\}$ of matrices partitioned in $2 \times 2,3 \times 3, \ldots, k \times k$ blocks, respectively.
Algorithm 3. Given $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ with size $m \times n$ and a bijection $\sigma$, the following algorithm constructs a sequence of matrices $\left\{Y_{0}, Y_{1}, \ldots, Y_{k-2}\right\}$.
if $\sigma$ has a consecution at 0 then

$$
Y_{0}=\left[\begin{array}{cc}
0 & A_{k} \\
I_{m} & A_{k-1}
\end{array}\right]
$$

else

$$
Y_{0}=\left[\begin{array}{cc}
0 & I_{n} \\
A_{k} & A_{k-1}
\end{array}\right]
$$

endif
for $i=1: k-2$
if $\sigma$ has a consecution at $i$ then

$$
Y_{i}=\left[\begin{array}{ccc}
Y_{i-1}(:, 1: i) & 0 & Y_{i-1}(:, i+1) \\
0 & I_{m} & A_{k-i-1}
\end{array}\right]
$$

else

$$
Y_{i}=\left[\begin{array}{cc}
Y_{i-1}(1: i,:) & 0 \\
0 & I_{n} \\
Y_{i-1}(i+1,:) & A_{k-i-1}
\end{array}\right]
$$

endif
endfor
In addition, the matrices $\left\{Y_{0}, Y_{1}, \ldots, Y_{k-2}\right\}$ satisfy the following properties:
(a) The size of $Y_{i}$ is

$$
(m+m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0: i))) \times(n+m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0: i)))
$$

(b) The last diagonal block of $Y_{i}$ is $A_{k-i-1}$ and, so, has size $m \times n$. The rest of diagonal blocks of $Y_{i}$ are square zero matrices, more precisely

$$
Y_{i}(j, j)=\left\{\begin{array}{ll}
0_{m} & \text { if } \sigma \text { has a consecution at } j-1 \\
0_{n} & \text { if } \sigma \text { has an inversion at } j-1
\end{array}, \quad \text { for } j=1,2, \ldots, i+1 .\right.
$$

(c) Let $Y_{i}(j, j)=0_{d_{j}}$, for $j=1,2, \ldots, i+1$, and define the $(i+2) \times(i+2)$ block reverse identities

$$
R_{l}^{(i)}:=\left[\begin{array}{llll} 
& & & I_{m} \\
& & I_{d_{i+1}} & \\
& . & & \\
I_{d_{1}} & & & \\
& & & \\
& & & \\
& & & \\
& & & \\
I_{n} & & & \\
& & & \\
I_{d_{i+1}} & & \\
& & &
\end{array}\right]
$$

Then we have

$$
\begin{equation*}
R_{l}^{(i)} Y_{i} R_{r}^{(i)}=W_{i}(-\operatorname{rev} P), \quad \text { for } i=0,1, \ldots, k-2, \tag{18}
\end{equation*}
$$

where $W_{i}(-\operatorname{rev} P)$ are the matrices constructed by Algorithm 2 for the polynomial $-\operatorname{rev} P(\lambda)$ and the bijection $\sigma$. In particular, according to Definition 3.8,

$$
R_{l}^{(k-2)} Y_{k-2} R_{r}^{(k-2)}=M_{\sigma}(-\operatorname{rev} P)
$$

Proof. The proof of the lemma up to part (b) (included) is analogous to the inductive proof of Theorem 3.6 and is omitted. We only indicate that part (b) for the block $Y_{i}(i+1, i+1)$ is a direct consequence of the way Algorithm 3 constructs $Y_{i}$, while the expression of the remaining $Y_{i}(j, j)$ blocks follows from $Y_{i}(1: i, 1: i)=Y_{i-1}(1: i, 1: i)$ via induction. It is important to note that the size of $Y_{i}(j, j)$ only depends on $j$ and not on $i$, whenever $1 \leq j \leq i+1$.

Before proving part (c), it is convenient to pay close attention to the structure of the matrices $R_{l}^{(i)}$ and $R_{r}^{(i)}$. First, note that the upper-right (resp. lower-left) block of $R_{l}^{(i)}$ (resp. $R_{r}^{(i)}$ ) is special because is always equal to $I_{m}$ (resp. $I_{n}$ ) independently of the consecutions/inversions that $\sigma$ may have. The reason of the presence of these special blocks is to make the product in (18) conformable since the last diagonal block of $Y_{i}$ has size $m \times n$. This motivates to define two matrices, $\widehat{R}_{l}^{(i)}$ and $\widehat{R}_{r}^{(i)}$, obtained from $R_{l}^{(i)}$ and $R_{r}^{(i)}$ by removing these special blocks and the corresponding rows/columns, that is,

$$
R_{l}^{(i)}=:\left[\begin{array}{ll} 
& I_{m}  \tag{19}\\
\widehat{R}_{l}^{(i)} &
\end{array}\right] \quad \text { and } \quad R_{r}^{(i)}=:\left[\begin{array}{cc} 
& \widehat{R}_{r}^{(i)} \\
I_{n} &
\end{array}\right] .
$$

Observe that the matrices $\widehat{R}_{l}^{(i)}$ and $\widehat{R}_{r}^{(i)}$ enjoy the following embedding properties

$$
\widehat{R}_{l}^{(i)}=\left[\begin{array}{ll} 
& I_{d_{i+1}}  \tag{20}\\
\widehat{R}_{l}^{(i-1)} &
\end{array}\right] \quad \text { and } \quad \widehat{R}_{r}^{(i)}=\left[\begin{array}{ll} 
& \widehat{R}_{r}^{(i-1)} \\
I_{d_{i+1}} &
\end{array}\right]
$$

that do not hold for the un-hatted matrices $R_{l}^{(i)}$ and $R_{r}^{(i)}$.
We are now in the position of proving (18) by induction on $i$. The definitions of $R_{l}^{(i)}$ and $R_{r}^{(i)}$ guarantee that the three factors in the left-hand side of (18) are conformal for multiplication. The initial step $i=0$ is proved directly, because for $i=0$, we have:

- If $\sigma$ has a consecution at 0 , then

$$
R_{l}^{(0)} Y_{0} R_{r}^{(0)}=\left[\begin{array}{ll} 
& I_{m} \\
I_{m} &
\end{array}\right]\left[\begin{array}{cc}
0 & A_{k} \\
I_{m} & A_{k-1}
\end{array}\right]\left[\begin{array}{cc} 
& I_{m} \\
I_{n} &
\end{array}\right]=\left[\begin{array}{cc}
A_{k-1} & I_{m} \\
A_{k} & 0
\end{array}\right]=W_{0}(-\operatorname{rev} P) .
$$

- If $\sigma$ has an inversion at 0 , then

$$
R_{l}^{(0)} Y_{0} R_{r}^{(0)}=\left[\begin{array}{ll} 
& I_{m} \\
I_{n} &
\end{array}\right]\left[\begin{array}{cc}
0 & I_{n} \\
A_{k} & A_{k-1}
\end{array}\right]\left[\begin{array}{cc} 
& I_{n} \\
I_{n} &
\end{array}\right]=\left[\begin{array}{cc}
A_{k-1} & A_{k} \\
I_{n} & 0
\end{array}\right]=W_{0}(-\operatorname{rev} P)
$$

Assume now that (18) is true for some $i-1$, such that $0 \leq(i-1) \leq k-3$, and we will prove it for $i$. We need to distinguish two cases according to whether $\sigma$ has a consecution or an inversion at $i$.
Case 1: $\sigma$ has a consecution at $i$. In this case $d_{i+1}=m$. Then (19) and (20) imply

$$
\begin{align*}
R_{l}^{(i)} Y_{i} R_{r}^{(i)} & =\left[\begin{array}{ll} 
& I_{m} \\
\widehat{R}_{l}^{(i)} &
\end{array}\right]\left[\begin{array}{ccc}
Y_{i-1}(:, 1: i) & 0 & Y_{i-1}(:, i+1) \\
0 & & I_{m} \\
A_{k-i-1}
\end{array}\right]\left[\begin{array}{cc} 
& \\
& \widehat{R}_{m}
\end{array}\right. \\
& =\left[\begin{array}{cccc}
A_{k-i-1} & I_{m} & \\
I_{n} & \\
\widehat{R}_{l}^{(i)} Y_{i-1}(:, i+1) & 0 & \widehat{R}_{l}^{(i)} Y_{i-1}(:, 1: i) \widehat{R}_{r}^{(i-1)}
\end{array}\right] \tag{21}
\end{align*}
$$

Observe that $d_{i+1}=m$, together with (19) and (20), imply that $R_{l}^{(i-1)}=\widehat{R}_{l}^{(i)}$. Now, we use the induction assumption, that is, that $(18)$ is true for $(i-1)$.

$$
\begin{align*}
W_{i-1}(-\operatorname{rev} P) & =R_{l}^{(i-1)} Y_{i-1} R_{r}^{(i-1)}=\widehat{R}_{l}^{(i)}\left[\begin{array}{lll}
Y_{i-1}(:, 1: i) & Y_{i-1}(:, i+1)
\end{array}\right]\left[\begin{array}{ll} 
& \widehat{R}_{r}^{(i-1)} \\
I_{n} &
\end{array}\right] \\
& =\left[\begin{array}{lll}
\widehat{R}_{l}^{(i)} Y_{i-1}(:, i+1) & \widehat{R}_{l}^{(i)} Y_{i-1}(:, 1: i) \widehat{R}_{r}^{(i-1)}
\end{array}\right] \tag{22}
\end{align*}
$$

We substitute equation (22) in (21) to get

$$
R_{l}^{(i)} Y_{i} R_{r}^{(i)}=\left[\begin{array}{ccc}
A_{k-i-1} & I_{m} & 0 \\
{\left[W_{i-1}(-\operatorname{rev} P)\right](:, 1)} & 0 & {\left[W_{i-1}(-\operatorname{rev} P)\right](:, 2: i+1)}
\end{array}\right]=W_{i}(-\operatorname{rev} P)
$$

where the last step follows from applying Algorithm 2 to $-\operatorname{rev} P(\lambda)$ and $\sigma$. This concludes the proof of Case 1.
Case 2: $\sigma$ has an inversion at $i$. In this case $d_{i+1}=n$. Then (19) and (20) imply

$$
\begin{align*}
R_{l}^{(i)} Y_{i} R_{r}^{(i)} & =\left[\right]\left[\begin{array}{cc}
Y_{i-1}(1: i,:) & 0 \\
0 & I_{n} \\
Y_{i-1}(i+1,:) & A_{k-i-1}
\end{array}\right]\left[\begin{array}{ll} 
& \widehat{R}_{r}^{(i)} \\
I_{n} &
\end{array}\right] \\
& =\left[\begin{array}{cc}
A_{k-i-1} & Y_{i-1}(i+1,:) \widehat{R}_{r}^{(i)} \\
I_{n} & 0 \\
0 & \widehat{R}_{l}^{(i-1)} Y_{i-1}(1: i,:) \widehat{R}_{r}^{(i)}
\end{array}\right] . \tag{23}
\end{align*}
$$

Observe that $d_{i+1}=n$, together with (19) and (20), imply that $R_{r}^{(i-1)}=\widehat{R}_{r}^{(i)}$. Now, we use the induction assumption.

$$
\begin{align*}
W_{i-1}(-\operatorname{rev} P) & =R_{l}^{(i-1)} Y_{i-1} R_{r}^{(i-1)}=\left[\begin{array}{ll} 
& I_{m} \\
\widehat{R}_{l}^{(i-1)} &
\end{array}\right]\left[\begin{array}{c}
Y_{i-1}(1: i,:) \\
Y_{i-1}(i+1,:)
\end{array}\right] \widehat{R}_{r}^{(i)} \\
& =\left[\begin{array}{c}
Y_{i-1}(i+1,:) \widehat{R}_{r}^{(i)} \\
\widehat{R}_{l}^{(i-1)} Y_{i-1}(1: i,:) \widehat{R}_{r}^{(i)}
\end{array}\right] \tag{24}
\end{align*}
$$

We substitute equation (24) in (23) to get

$$
R_{l}^{(i)} Y_{i} R_{r}^{(i)}=\left[\begin{array}{cc}
A_{k-i-1} & {\left[W_{i-1}(-\operatorname{rev} P)\right](1,:)} \\
I_{n} & 0 \\
0 & {\left[W_{i-1}(-\operatorname{rev} P)\right](2: i+1,:)}
\end{array}\right]=W_{i}(-\operatorname{rev} P)
$$

This concludes the proof of Case 2.
Lemma 3.13 shows the result of certain matrix multiplications that are used in the proof of Theorem 3.14 to perform strict equivalences on the reversals of Fiedler pencils when the degree $k$ of the polynomial satisfies $k \geq 3$.

Lemma 3.13. Let $\sigma, \tau:\{0,1, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$ be two bijections such that $\sigma$ has a consecution (resp. inversion) at $i-1$ if and only if $\tau$ has a consecution (resp. inversion) at $k-i-1$ for $i=1, \ldots, k-1$. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 3$, let $\left\{W_{i}\right\}_{i=0}^{k-2}$ be the sequence of block partitioned matrices constructed by Algorithm 2 for $P(\lambda)$ and $\sigma$, let $\left\{Y_{i}\right\}_{i=0}^{k-2}$ be the sequence of block partitioned matrices constructed by Algorithm 3 for $P(\lambda)$ and $\tau$, and let us define $W_{-1}:=-A_{0}$, $Y_{-1}:=A_{k}$. Let us define two sequences, $\left\{\widetilde{I}_{i}\right\}_{i=0}^{k-1}$ and $\left\{\stackrel{\circ}{I}_{i}\right\}_{i=0}^{k-1}$, of partitioned matrices as follows: $\widetilde{I}_{0}$ and $\stackrel{\circ}{I}_{0}$ are $0 \times 0$ empty matrices, and

$$
\widetilde{I}_{i}:=\left[\begin{array}{cccc}
I_{s_{1}} & & & \\
& I_{s_{2}} & & \\
& & \ddots & \\
& & & I_{s_{i}}
\end{array}\right] \quad \text { and } \quad \stackrel{\circ}{I_{i}}:=\left[\begin{array}{cccc}
I_{t_{k-i+1}} & & & \\
& I_{t_{k-i+2}} & & \\
& & \ddots & \\
& & & I_{t_{k}}
\end{array}\right], \quad \text { for } i=1, \ldots, k-1,
$$

where $\left\{s_{j}\right\}_{j=1}^{i}$ are the sizes of the square diagonal blocks $\left\{Y_{i-1}(j, j)\right\}_{j=1}^{i}$ and $\left\{t_{j}\right\}_{j=k-i+1}^{k}$ are the sizes of the square diagonal blocks $\left\{W_{i-1}(j, j)\right\}_{j=2}^{i+1}$. Then the following statements hold.
(a) For each $i=0,1, \ldots, k-1$, the matrices

$$
\widetilde{W}_{i-1}:=\left[\begin{array}{ll}
\widetilde{I}_{k-i-1} & \\
& W_{i-1}
\end{array}\right] \quad \text { and } \quad \widetilde{Y}_{k-i-2}:=\left[\begin{array}{ll}
Y_{k-i-2} & \\
& \stackrel{\circ}{I}_{i}
\end{array}\right]
$$

are partitioned into $k \times k$ blocks and the size of the block $\widetilde{W}_{i-1}(p, q)$ is equal to the size of the block $\widetilde{Y}_{k-i-2}(p, q)$ for all $1 \leq p, q \leq k$. In addition, $\widetilde{W}_{i-1}$ and $\widetilde{Y}_{k-i-2}$ have both size

$$
(m+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)) \times(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)),
$$

that is, the same size as the Fiedler pencil of $P(\lambda)$ associated with $\sigma$.
(b) Define a sequence of matrices $\left\{S_{i}\right\}_{i=1}^{k-1}$ as follows

$$
S_{1}:=\left[\begin{array}{ccc}
\widetilde{I}_{k-2} & & \\
& 0 & I_{n} \\
& I_{m} & A_{1}
\end{array}\right], S_{i}:=\left[\begin{array}{cccccc}
\widetilde{I}_{k-i-1} & & & & & \\
& 0 & I_{n} & & & \\
& I_{m} & A_{i} & & & \\
& & & I_{t_{k-i+2}} & & \\
& & & & \ddots & \\
& & & & & I_{t_{k}}
\end{array}\right], i=2, \ldots, k-1
$$

Then, for each $i=1, \ldots, k-1$, the following statements hold:
(b1) If $\sigma$ has a consecution at $i-1$, then $S_{i}$ has size $(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)) \times(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma))$ and

$$
\widetilde{W}_{i-1} S_{i}=\widetilde{W}_{i-2} \quad \text { and } \quad \widetilde{Y}_{k-i-2} S_{i}=\widetilde{Y}_{k-i-1}
$$

(b2) If $\sigma$ has an inversion at $i-1$, then $S_{i}$ has size $(m+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)) \times(m+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma))$ and

$$
S_{i} \widetilde{W}_{i-1}=\widetilde{W}_{i-2} \quad \text { and } \quad S_{i} \widetilde{Y}_{k-i-2}=\widetilde{Y}_{k-i-1}
$$

Proof. Part (a). $\widetilde{I}_{k-i-1}$ has $(k-i-1) \times(k-i-1)$ blocks, by definition, and, by Theorem 3.6, $W_{i-1}$ has $(i+1) \times(i+1)$ blocks. So $\widetilde{W}_{i-1}$ has $k \times k$ blocks. Analogously, by Lemma 3.12, $Y_{k-i-2}$ has $(k-i) \times(k-i)$ blocks and, by definition, $\stackrel{\circ}{I}_{i}$ has $i \times i$ blocks. So $\widetilde{Y}_{k-i-2}$ has $k \times k$ blocks.

Next we prove that the sizes of the blocks of $\widetilde{W}_{i-1}$ are equal to the sizes of the corresponding blocks of $\widetilde{Y}_{k-i-2}$. Assume first that $1 \leq i \leq k-2$, and recall that $W_{i-1}(1,1)=-A_{i}$ and $Y_{k-i-2}(k-i, k-i)=A_{i+1}$. Then

$$
\begin{align*}
\widetilde{W}_{i-1} & =\left[\begin{array}{ccc}
\widetilde{I}_{k-i-1} & & \\
& -A_{i} & * \\
& * & W_{i-1}(2: i+1,2: i+1)
\end{array}\right],  \tag{25}\\
\widetilde{Y}_{k-i-2} & =\left[\begin{array}{ccc}
Y_{k-i-2}(1: k-i-1,1: k-i-1) & * \\
& * & A_{i+1} \\
& & \circ_{i}
\end{array}\right] . \tag{26}
\end{align*}
$$

By definition, $\widetilde{I}_{k-i-1}$ is partitioned into blocks exactly as $Y_{k-i-2}(1: k-i-1,1: k-i-1)$ and $\stackrel{\circ}{I}_{i}$ is partitioned exactly as $W_{i-1}(2: i+1,2: i+1)$. Therefore, $\widetilde{W}_{i-1}$ and $\widetilde{Y}_{k-i-2}$ have corresponding blocks with equal sizes. For $i=0$ we have $\widetilde{W}_{-1}=\operatorname{diag}\left(\widetilde{I}_{k-1},-A_{0}\right), \widetilde{Y}_{k-2}=Y_{k-2}$, and the definition of $\widetilde{I}_{k-1}$ together with Lemma 3.12-(b) guarantee that the sizes of corresponding blocks are equal. For $i=k-1$ we have $\widetilde{W}_{k-2}=W_{k-2}, \widetilde{Y}_{-1}=\operatorname{diag}\left(A_{k}, \stackrel{\circ}{I}_{k-1}\right)$, and the definition of $\stackrel{\circ}{I}_{k-1}$ together with Theorem 3.6-(b) imply the result.

We consider now the total size of the matrices $\widetilde{W}_{i-1}$ and $\widetilde{Y}_{k-i-2}$. Note first that $\mathfrak{c}(\sigma)=\mathfrak{c}(\tau)$ and $\mathfrak{i}(\sigma)=\mathfrak{i}(\tau)$. For $i=0$, we get from the previous discussion that $\widetilde{W}_{-1}$ and $\widetilde{Y}_{k-2}=Y_{k-2}$ have both the size of $Y_{k-2}$, that is $(m+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)) \times(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma))$ according to Lemma 3.12-(a). For $i=k-1$, we get from the previous discussion that $\widetilde{W}_{k-2}=W_{k-2}$ and $\widetilde{Y}_{-1}$ have both the size of $W_{k-2}$, that is $(m+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)) \times(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma))$ according to Theorem 3.6-(a). For $1 \leq i \leq k-2$, we get again from the previous discussion that $\widetilde{W}_{i-1}$ and $\widetilde{Y}_{k-i-2}$ have both the same size. This size is the sum of the sizes of the three diagonal blocks in (25) (or (26)), which according to Theorem 3.6 and Lemma 3.12 is $(m+r) \times(n+r)$ with

$$
r=m(\mathfrak{c}(\tau(0: k-i-2))+\mathfrak{c}(\sigma(0: i-1)))+n(\mathfrak{i}(\tau(0: k-i-2))+\mathfrak{i}(\sigma(0: i-1)))
$$

From the definition of $\tau$, we see that $r$ is equal to

$$
\begin{aligned}
r & =m(\mathfrak{c}(\sigma(i: k-2))+\mathfrak{c}(\sigma(0: i-1)))+n(\mathfrak{i}(\sigma(i: k-2))+\mathfrak{i}(\sigma(0: i-1))) \\
& =m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma) .
\end{aligned}
$$

This concludes the proof of Part (a).
Part (b). For brevity, we prove only (b1). The proof of (b2) is similar and is omitted. Let us establish the size of $S_{i}$ that is clearly an square matrix for each $i$. So, we only pay attention to the number of rows. Consider first the number of rows of $S_{1}$. Note that if $\sigma$ has a consecution at 0 , then

$$
\widetilde{W}_{0}=\left[\begin{array}{ll}
\widetilde{I}_{k-2} & \\
& W_{0}
\end{array}\right]=\left[\begin{array}{ccc}
\widetilde{I}_{k-2} & & \\
& -A_{1} & I_{m} \\
& -A_{0} & 0
\end{array}\right] .
$$

This makes evident that number of columns of $\widetilde{W}_{0}$ is equal to the number of rows of $S_{1}$ and therefore this number is $(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma))$ by Part (a). Note that we have proved in addition that the partitions of $\widetilde{W}_{0}$ and $S_{1}$ are conformal for the product $\widetilde{W}_{0} S_{1}$, and, by Part (a), the same happens for the product $\widetilde{Y}_{k-3} S_{1}$. Consider next the number of columns of $S_{i}$, for $i=2, \ldots, k-1$. Note that if $\sigma$ has a consecution at $i-1$, then

$$
\widetilde{W}_{i-1}=\left[\begin{array}{lll}
\widetilde{I}_{k-i-1} & \\
& W_{i-1}
\end{array}\right]=\left[\begin{array}{cccc}
\widetilde{I}_{k-i-1} & & & \\
& -A_{i} & I_{m} & 0 \\
& W_{i-2}(:, 1) & 0 & W_{i-2}(:, 2: i)
\end{array}\right] .
$$

By definition, $t_{k-i+2}, \ldots, t_{k}$ are the number of columns of the block columns of $W_{i-1}(:, 3: i+1)$, which have the same number of columns as the block columns of $W_{i-2}(:, 2: i)$. Therefore, the number of columns of $\widetilde{W}_{i-1}$ is equal to the number of rows of $S_{i}$ and this number is $(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma))$ by Part (a). Observe that we have also proved that the partitions of $\widetilde{W}_{i-1}$ and $S_{i}$ are conformal for the product $\widetilde{W}_{i-1} S_{i}$. This implies, by Part (a), that the partitions of $\widetilde{Y}_{k-i-2}$ and $S_{i}$ are conformal for the product $\widetilde{Y}_{k-i-2} S_{i}$.

In the proof of $\widetilde{W}_{i-1} S_{i}=\widetilde{W}_{i-2}$ for $i=1, \ldots, k-1$, we need to deal separately with the case $i=1$, because $\widetilde{W}_{-1}$ has an structure different from $\widetilde{W}_{i}$ for $i>-1$. A direct block multiplication shows that

$$
\widetilde{W}_{0} S_{1}=\left[\begin{array}{ccc}
\widetilde{I}_{k-2} & & \\
& -A_{1} & I_{m} \\
& -A_{0} & 0
\end{array}\right]\left[\begin{array}{ccc}
\widetilde{I}_{k-2} & & \\
& 0 & I_{n} \\
& I_{m} & A_{1}
\end{array}\right]=\left[\begin{array}{ccc}
\widetilde{I}_{k-2} & & \\
& I_{m} & 0 \\
& 0 & -A_{0}
\end{array}\right]=\widetilde{W}_{-1}
$$

where we have used that $\widetilde{I}_{k-1}=\operatorname{diag}\left(\widetilde{I}_{k-2}, I_{m}\right)$, because, according to Lemma 3.12-(b), the sizes of the blocks $\left\{Y_{k-3}(j, j)\right\}_{j=1}^{k-2}$ are equal to the sizes of the blocks $\left\{Y_{k-2}(j, j)\right\}_{j=1}^{k-2}$, and $Y_{k-2}(k-1, k-1)=0_{m}$ because $\sigma$ has a consecution at 0 , that is, $\tau$ has a consecution at $k-2$. Let us consider now $i=2, \ldots, k-1$. Then

$$
\begin{aligned}
\widetilde{W}_{i-1} S_{i} & =\left[\begin{array}{cccc}
\widetilde{I}_{k-i-1} & & & \\
& -A_{i} & I_{m} & 0 \\
& W_{i-2}(:, 1) & 0 & W_{i-2}(:, 2: i)
\end{array}\right]\left[\begin{array}{ccc}
\widetilde{I}_{k-i-1} & & \\
& 0 & I_{n} \\
& & I_{m} \\
& A_{i} & \\
& & \\
& & \\
& & I_{t_{k-i+2}+\cdots+t_{k}}
\end{array}\right] \\
& {\left[\begin{array}{cccc}
\widetilde{I}_{k-i-1} & & & \\
& I_{m} & 0 & 0 \\
& 0 & W_{i-2}(:, 1) & W_{i-2}(:, 2: i)
\end{array}\right]=\left[\begin{array}{cc}
\widetilde{I}_{k-i} & \\
& W_{i-2}
\end{array}\right]=\widetilde{W}_{i-2}, }
\end{aligned}
$$

where we have used that $\widetilde{I}_{k-i}=\operatorname{diag}\left(\widetilde{I}_{k-i-1}, I_{m}\right)$, because, according to Lemma 3.12-(b), the sizes of the blocks $\left\{Y_{k-i-2}(j, j)\right\}_{j=1}^{k-i-1}$ are equal to the sizes of the blocks $\left\{Y_{k-i-1}(j, j)\right\}_{j=1}^{k-i-1}$, and $Y_{k-i-1}(k-i, k-$ $i)=0_{m}$ because $\sigma$ has a consecution at $i-1$, that is, $\tau$ has a consecution at $k-i-1$.

Next, we proceed with the proof of $\widetilde{Y}_{k-i-2} S_{i}=\widetilde{Y}_{k-i-1}$. Here we need to deal separately with the case $i=k-1$ because $\widetilde{Y}_{-1}$ has an structure different from the remaining $\widetilde{Y}_{i}$. We consider first $i=1, \ldots, k-2$. Since $\sigma$ has a consecution at $i-1$, we have that $W_{i-1}(2,2)=0_{m}$ by Theorem 3.6-(b) and the first block of $\stackrel{\circ}{I}_{i}$ is $I_{t k-i+1}=I_{m}$. In addition, note that $\stackrel{\circ}{I}_{i}=\operatorname{diag}\left(I_{m}, \stackrel{\circ}{I}_{i-1}\right)$ because the sizes of the blocks $\left\{W_{i-1}(j, j)\right\}_{j=3}^{i+1}$ are equal to the sizes of the blocks $\left\{W_{i-2}(j, j)\right\}_{j=2}^{i}$ by Theorem 3.6-(b) (recall also

Remark 3.7). Therefore

$$
\begin{aligned}
& \widetilde{Y}_{k-i-2} S_{i}=\left[\begin{array}{llll}
Y_{k-i-2} & & \\
& I_{m} & \\
& & I_{t_{k-i+2}+\cdots+t_{k}}
\end{array}\right]\left[\begin{array}{cccc}
\widetilde{I}_{k-i-1} & & & \\
& 0 & I_{n} & \\
& I_{m} & A_{i} & \\
& & & I_{t_{k-i+2}+\cdots+t_{k}}
\end{array}\right] \\
& =\left[\begin{array}{lllll}
Y_{k-i-2}(:, 1: 1: k-i-1) & Y_{k-i-2}(:, k-i) & & \\
& & I_{m} & \\
& & & \circ_{i-1}
\end{array}\right]\left[\begin{array}{cccc}
\widetilde{I}_{k-i-1} & & & \\
& 0 & I_{n} & \\
& I_{m} & A_{i} & \\
& & & \check{I}_{i-1}
\end{array}\right] \\
& =\left[\begin{array}{cccc}
Y_{k-i-2}(:, 1: k-i-1) & & Y_{k-i-2}(:, k-i) & \\
& I_{m} & A_{i} & \circ_{i-1}
\end{array}\right]=\left[\begin{array}{lll}
Y_{k-i-1} & \\
& \circ_{I} \\
& & \check{I}_{i-1}
\end{array}\right]=\widetilde{Y}_{k-i-1},
\end{aligned}
$$

where we have used Algorithm 3 taking into account that $\tau$ has a consecution at $k-i-1$. We finally cover the case $i=k-1$. Since $\sigma$ has a consecution at $k-2$, an argument similar to the one above shows that $\stackrel{\circ}{I}_{k-1}=\operatorname{diag}\left(I_{m}, \stackrel{\circ}{I}_{k-2}\right)$. Therefore,

$$
\begin{aligned}
& \widetilde{Y}_{-1} S_{k-1}=\left[\begin{array}{ll}
A_{k} & \\
& \stackrel{\circ}{I}_{k-1}
\end{array}\right]\left[\begin{array}{ccc}
0 & I_{n} & \\
I_{m} & A_{k-1} & \\
& & I_{t_{3}+\cdots+t_{k}}
\end{array}\right]=\left[\begin{array}{llll}
A_{k} & & \\
& I_{m} & \\
& & \stackrel{\circ}{I}_{k-2}
\end{array}\right]\left[\begin{array}{ccc}
0 & I_{n} & \\
I_{m} & A_{k-1} & \\
& & \circ^{\circ} \\
& & \\
& & \\
& &
\end{array}\right] \\
& =\left[\begin{array}{ccc}
0 & A_{k} & \\
I_{m} & A_{k-1} & \\
& & \stackrel{\circ}{I}_{k-2}
\end{array}\right]=\left[\begin{array}{ll}
Y_{0} & \\
& \stackrel{\circ}{I}_{k-2}
\end{array}\right]=\widetilde{Y}_{0},
\end{aligned}
$$

where we have used Algorithm 3 taking into account that $\tau$ has a consecution at 0 . This concludes the proof of (b1).

Now we are in the position of proving the main result in this section.
Theorem 3.14. Let $P(\lambda)$ be an $m \times n$ matrix polynomial with degree $k \geq 2$ and let $F_{\sigma}(P)$ be the Fiedler pencil of $P(\lambda)$ associated with a bijection $\sigma$. Then $\operatorname{rev} F_{\sigma}(P)$ is strictly equivalent to a Fiedler pencil of $\operatorname{rev} P(\lambda)$. More precisely, $\operatorname{rev} F_{\sigma}(P)$ is strictly equivalent to $F_{\tau}(\operatorname{rev} P)$, where $\tau:\{0,1, \ldots, k-1\} \rightarrow$ $\{1, \ldots, k\}$ is any bijection such that $\tau$ has a consecution (resp. inversion) at $k-i-1$ if $\sigma$ has a consecution (resp. inversion) at $i-1$, for $i=1, \ldots, k-1$.
Proof. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$. If the degree is $k=2$, then Algorithm 2 shows that there are only two different Fiedler pencils. These are the two companion forms

$$
C_{1}(P)=\lambda\left[\begin{array}{cc}
A_{2} & 0 \\
0 & I_{n}
\end{array}\right]+\left[\begin{array}{cc}
A_{1} & A_{0} \\
-I_{n} & 0
\end{array}\right] \quad \text { and } \quad C_{2}(P)=\lambda\left[\begin{array}{cc}
A_{2} & 0 \\
0 & I_{m}
\end{array}\right]+\left[\begin{array}{cc}
A_{1} & -I_{m} \\
A_{0} & 0
\end{array}\right]
$$

For $k=2$, direct matrix multiplications show that

$$
\begin{aligned}
& {\left[\begin{array}{cc}
I_{m} & A_{1} \\
0 & -I_{n}
\end{array}\right]\left(\operatorname{rev} C_{1}(P)\right)\left[\begin{array}{cc}
0 & I_{n} \\
I_{n} & 0
\end{array}\right]=C_{1}(\operatorname{rev} P) \quad \text { and }} \\
& {\left[\begin{array}{cc}
0 & I_{m} \\
I_{m} & 0
\end{array}\right]\left(\operatorname{rev} C_{2}(P)\right)\left[\begin{array}{cc}
I_{n} & 0 \\
A_{1} & -I_{m}
\end{array}\right]=C_{2}(\operatorname{rev} P),}
\end{aligned}
$$

which proves the result because the matrices multiplying rev $C_{1}(P)$ and $\operatorname{rev} C_{2}(P)$ are always nonsingular. Observe that for $k=2$ the bijections $\sigma$ and $\tau$ are equal.

For $k \geq 3$, the proof relies in Lemma 3.13, so we use in the rest of the proof exactly the same definitions and notation as in Lemma 3.13. Note that $\widetilde{W}_{k-2}=W_{k-2}=M_{\sigma}$, where $-M_{\sigma}$ is the zero degree term of $F_{\sigma}(P)$ according to Definition 3.8. Then

$$
F_{\sigma}(P)=\lambda\left[\begin{array}{cc}
A_{k} & \\
& I_{m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)}
\end{array}\right]-M_{\sigma}=\lambda \widetilde{Y}_{-1}-\widetilde{W}_{k-2} \quad \text { and } \quad \operatorname{rev} F_{\sigma}(P)=\widetilde{Y}_{-1}-\lambda \widetilde{W}_{k-2}
$$

Next we use the nonsingular matrices $S_{k-1}, S_{k-2}, \ldots, S_{1}$ introduced in Lemma 3.13-(b) and multiply $\operatorname{rev} F_{\sigma}(P)$ first by $S_{k-1}$, second by $S_{k-2}$, and so on until we multiply by $S_{1}$. The multiplications are performed on the right or on the left according to the consecutions or inversions of $\sigma$ as indicated in Lemma 3.13-(b1)-(b2). So, we obtain

$$
\operatorname{rev} F_{\sigma}(P)=\widetilde{Y}_{-1}-\lambda \widetilde{W}_{k-2} \sim_{s} \widetilde{Y}_{0}-\lambda \widetilde{W}_{k-3} \sim_{s} \widetilde{Y}_{1}-\lambda \widetilde{W}_{k-4} \sim_{s} \cdots \sim_{s} \widetilde{Y}_{k-2}-\lambda \widetilde{W}_{-1}
$$

where the symbols $\sim_{s}$ denote that we are performing strict equivalences, because the matrices $S_{i}$ are always nonsingular. From Lemma 3.13 we see that $\widetilde{Y}_{k-2}=Y_{k-2}$ and $\widetilde{W}_{-1}=\operatorname{diag}\left(\widetilde{I}_{k-1},-A_{0}\right)$. Therefore

$$
\operatorname{rev} F_{\sigma}(P) \sim_{s} Y_{k-2}-\lambda \operatorname{diag}\left(\widetilde{I}_{k-1},-A_{0}\right)
$$

We apply now Lemma 3.12-(c) to get

$$
\begin{aligned}
\operatorname{rev} F_{\sigma}(P) & \sim_{s} R_{l}^{(k-2)} Y_{k-2} R_{r}^{(k-2)}-\lambda R_{l}^{(k-2)} \operatorname{diag}\left(\widetilde{I}_{k-1},-A_{0}\right) R_{r}^{(k-2)} \\
& =M_{\tau}(-\operatorname{rev} P)-\lambda \operatorname{diag}\left(-A_{0}, \widetilde{I}_{k-1}\right)=-F_{\tau}(-\operatorname{rev} P)
\end{aligned}
$$

Finally, by Lemma 3.11, $-F_{\tau}(-\operatorname{rev} P)$ is strictly equivalent to $-F_{\tau}(\operatorname{rev} P)$, which in turn is strictly equivalent to $F_{\tau}(\operatorname{rev} P)$. Hence, we conclude that $\operatorname{rev} F_{\sigma}(P)$ is strictly equivalent to $F_{\tau}(\operatorname{rev} P)$.

## 4. Fiedler pencils of rectangular matrix polynomials are strong linearizations

We will prove in this section that all Fiedler pencils $F_{\sigma}(\lambda)$ of a rectangular matrix polynomial $P(\lambda)$ are strong linearizations for $P(\lambda)$. This is proved in Theorem 4.5, which generalizes in a nontrivial way Theorem 4.6 in [11]. The approach we follow is constructive, in the sense that we will show how to construct appropriate unimodular matrices $U(\lambda)$ and $V(\lambda)$ satisfying (3) for every $F_{\sigma}(\lambda)$. The construction of these matrices is accomplished via the construction of sequences of block partitioned matrices in Algorithms 4 and 5, which follow the spirit of Definition 3.8 of Fiedler pencils for rectangular polynomials, and the unimodular transformations generated by these sequences are considered in Lemma 4.4. In this section, we will make a systematic use of the Horner shifts introduced in Definition 4.1.

Definition 4.1. Let $P(\lambda)=A_{0}+\lambda A_{1}+\cdots+\lambda^{k} A_{k}$ be an $m \times n$ matrix polynomial of degree $k$. For $d=$ $0, \ldots, k$, the degree $d$ Horner shift of $P(\lambda)$ is the matrix polynomial $P_{d}(\lambda):=A_{k-d}+\lambda A_{k-d+1}+\cdots+\lambda^{d} A_{k}$.

Observe that the Horner shifts of $P(\lambda)$ satisfy the following recurrence relation

$$
P_{0}(\lambda)=A_{k}, \quad P_{d+1}(\lambda)=\lambda P_{d}(\lambda)+A_{k-d-1}, \text { for } 0 \leq d \leq k-1, \quad \text { and } \quad P_{k}(\lambda)=P(\lambda) .
$$

Lemma 4.2. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$ and let $\sigma$ : $\{0,1, \ldots, k-1\} \rightarrow\{1,2, \ldots, k\}$ be a bijection. Then the following two statements hold.

1. If in Algorithm 4 below each matrix $N_{i}$, for $i=1,2, \ldots, k-2$, is partitioned into blocks in such a way that the blocks of $N_{i-1}$ are blocks of $N_{i}$, then Algorithm 4 constructs a sequence $\left\{N_{0}, N_{1}, \ldots, N_{k-2}\right\}$ of matrices partitioned into $2 \times 2,3 \times 3, \ldots, k \times k$ blocks, respectively.
2. If in Algorithm 5 below each matrix $H_{i}$, for $i=1,2, \ldots, k-2$, is partitioned into blocks in such a way that the blocks of $H_{i-1}$ are blocks of $H_{i}$, then Algorithm 5 constructs a sequence $\left\{H_{0}, H_{1}, \ldots, H_{k-2}\right\}$ of matrices partitioned into $2 \times 2,3 \times 3, \ldots, k \times k$ blocks, respectively.
Algorithm 4. Given $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ with size $m \times n$ and a bijection $\sigma$, the following algorithm constructs a sequence of matrices $\left\{N_{0}, N_{1}, \ldots, N_{k-2}\right\}$. Note that $P_{d}$ denotes the degree $d$ Horner shift of $P(\lambda)$ and that the dependence on $\lambda$ is dropped for simplicity both in $P_{d}$ and in $\left\{N_{i}\right\}_{i=0}^{k-2}$.
if $\sigma$ has a consecution at 0 then
$N_{0}=\left[\begin{array}{cc}I_{m} & 0 \\ \lambda I_{m} & I_{m}\end{array}\right]$
else

$$
N_{0}=\left[\begin{array}{cc}
0 & -I_{n} \\
I_{m} & P_{k-1}
\end{array}\right]
$$

endif
for $i=1: k-2$
if $\sigma$ has a consecution at $i$ then

$$
N_{i}=\left[\begin{array}{cc}
I_{m} & 0 \\
\lambda N_{i-1}(:, 1) & N_{i-1}
\end{array}\right]
$$

else

$$
N_{i}=\left[\begin{array}{ccc}
0 & -I_{n} & 0 \\
N_{i-1}(:, 1) & N_{i-1}(:, 1) P_{k-i-1} & N_{i-1}(:, 2: i+1)
\end{array}\right]
$$

endif
endfor

Algorithm 5. Given $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ with size $m \times n$ and a bijection $\sigma$, the following algorithm constructs a sequence of matrices $\left\{H_{0}, H_{1}, \ldots, H_{k-2}\right\}$. Note that $P_{d}$ denotes the degree d Horner shift of $P(\lambda)$ and that the dependence on $\lambda$ is dropped for simplicity both in $P_{d}$ and in $\left\{H_{i}\right\}_{i=0}^{k-2}$.
if $\sigma$ has a consecution at 0 then

$$
H_{0}=\left[\begin{array}{cc}
0 & I_{n} \\
-I_{m} & P_{k-1}
\end{array}\right]
$$

else

$$
H_{0}=\left[\begin{array}{cc}
I_{n} & \lambda I_{n} \\
0 & I_{n}
\end{array}\right]
$$

endif
for $i=1: k-2$
if $\sigma$ has a consecution at $i$ then

$$
H_{i}=\left[\begin{array}{cc}
0 & H_{i-1}(1,:) \\
-I_{m} & P_{k-i-1} H_{i-1}(1,:) \\
0 & H_{i-1}(2: i+1,:)
\end{array}\right]
$$

else

$$
H_{i}=\left[\begin{array}{cc}
I_{n} & \lambda H_{i-1}(1,:) \\
0 & H_{i-1}
\end{array}\right]
$$

endif
endfor
Moreover, if $\left\{W_{i}\right\}_{i=0}^{k-2}$ is the sequence of block partitioned matrices constructed by Algorithm 2 for $P(\lambda)$ and $\sigma$, then the matrices $\left\{N_{i}\right\}_{i=0}^{k-2}$ and $\left\{H_{i}\right\}_{i=0}^{k-2}$ satisfy the following properties.
(a) For $0 \leq i \leq k-2$ and $1 \leq j \leq i+2$ the number of columns of $N_{i}(:, j)$ is equal to the number of rows of $W_{i}(j,:)$. This means that the matrix product $N_{i}(:, j) W_{i}(j,:)$ is well defined.
(b) For $0 \leq i \leq k-2$ and $1 \leq j \leq i+2$ the number of columns of $W_{i}(:, j)$ is equal to the number of rows of $H_{i}(j,:)$. This means that the matrix product $W_{i}(:, j) H_{i}(j,:)$ is well defined.
(c) The size of $N_{i}$ is $(m+m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0: i))) \times(m+m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0: i)))$.
(d) The size of $H_{i}$ is $(n+m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0: i))) \times(n+m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0: i)))$.
(e) The matrices $N_{i}$ and $H_{i}$ are unimodular. In fact $\operatorname{det}\left(N_{i}\right)= \pm 1$ and $\operatorname{det}\left(H_{i}\right)= \pm 1$.
(f) $N_{i}(i+2,:)$ has $m$ rows and $H_{i}(:, i+2)$ has $n$ columns, that is, the last block row of $N_{i}$ has $m$ rows and the last block column of $H_{i}$ has $n$ columns.

Proof. The proof is trivial by induction. We indicate only the main points. Starting with $N_{0}$ and $H_{0}$, it is obvious to see by induction that $N_{i}(:, 1)$ has $m$ columns and $H_{i}(1,:)$ has $n$ rows for $i=0,1, \ldots, k-2$. Therefore, the sequences $\left\{N_{i}\right\}_{i=0}^{k-2}$ and $\left\{H_{i}\right\}_{i=0}^{k-2}$ are well defined. Since, for $i \geq 1, N_{i}$ and $H_{i}$ are obtained from $N_{i-1}$ and $H_{i-1}$, respectively, by adding one block row and one block column, then we get that $N_{i}$ and $H_{i}$ are partitioned into $(i+2) \times(i+2)$ blocks.

To prove part (a), note that the result is true for $N_{0}$ and $W_{0}$, and make the induction assumption that is true for $N_{i-1}$ and $W_{i-1}$ with $(i-1) \geq 0$. Then the constructions of $N_{i}$ in Algorithm 4 and $W_{i}$ in Algorithm 2 make evident that the result is true for $N_{i}$ and $W_{i}$. Part (b) follows from a similar inductive argument.

To prove parts (c) and (d), we note first that all matrices in the sequences $\left\{N_{i}\right\}_{i=0}^{k-2}$ and $\left\{H_{i}\right\}_{i=0}^{k-2}$ are square, because, by definition, this is true for $N_{0}$ and $W_{0}$ and, for $i \geq 1, N_{i}$ is obtained from $N_{i-1}$ by adding the same number of rows as columns and $H_{i}$ is also obtained from $H_{i-1}$ by adding the same number of rows as columns. Then (c) follows from (a), and (d) from (b), by using Theorem 3.6-(a). Finally Parts (e) and (f) follow again by induction.

Since the matrices $N_{k-2}$ and $H_{k-2}$ in Lemma 4.2 will play a key role in the rest of the paper, next we give them an special name.
Definition 4.3. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, let $\sigma$ : $\{0,1, \ldots, k-1\} \rightarrow\{1,2, \ldots, k\}$ be a bijection, and let $N_{k-2}$ and $H_{k-2}$ be, respectively, the last matrices of the sequences constructed by Algorithms 4 and 5 for $P(\lambda)$ and $\sigma$. Then the left unimodular equivalence matrix associated with $P(\lambda)$ and $\sigma$ is the matrix $U_{\sigma}(\lambda):=N_{k-2}$ and the right unimodular equivalence matrix associated with $P(\lambda)$ and $\sigma$ is the matrix $V_{\sigma}(\lambda):=H_{k-2}$.

Lemma 4.4 studies the unimodular transformations generated by the sequences $\left\{N_{i}\right\}_{i=0}^{k-2}$ and $\left\{H_{i}\right\}_{i=0}^{k-2}$.
Lemma 4.4. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, let $\sigma:\{0, \ldots, k-1\}$ $\rightarrow\{1, \ldots, k\}$ be a bijection, let $\left\{W_{i}\right\}_{i=0}^{k-2},\left\{N_{i}\right\}_{i=0}^{k-2},\left\{H_{i}\right\}_{i=0}^{k-2}$ be the sequences of block partitioned matrices constructed, respectively, by Algorithms 2, 4 and 5 for $P(\lambda)$ and $\sigma$, and let $\alpha_{i}=m \mathfrak{c}(\sigma(0: i))+n \mathfrak{i}(\sigma(0$ : i)) and $\beta_{i}=m \mathfrak{c}(\sigma(i))+n \mathfrak{i}(\sigma(i))$, that is, $\beta_{i}=m$ if $\sigma$ has a consecution at $i$ and $\beta_{i}=n$ if $\sigma$ has an inversion at $i$, for $i=0,1, \ldots, k-2$. Then, the following two identities hold.
(a) For $1 \leq i \leq k-2$,

$$
N_{i}\left(\left[\begin{array}{ll}
\lambda P_{k-i-2} & \\
& \lambda I_{\alpha_{i}}
\end{array}\right]-W_{i}\right) H_{i}=\left[\begin{array}{ccc}
I_{\beta_{i}} & & \\
& N_{i-1}\left(\left[\begin{array}{ll}
\lambda P_{k-i-1} & \\
& \lambda I_{\alpha_{i-1}}
\end{array}\right]-W_{i-1}\right) H_{i-1}
\end{array}\right]
$$

(b) and, for $i=0$,

$$
N_{0}\left(\left[\begin{array}{ll}
\lambda P_{k-2} & \\
& \lambda I_{\alpha_{0}}
\end{array}\right]-W_{0}\right) H_{0}=\left[\begin{array}{ll}
I_{\beta_{0}} & \\
& P
\end{array}\right],
$$

where $P_{d}$ is the degree $d$ Horner shift of $P$ and the dependences on $\lambda$ are dropped for simplicity.
Proof. Observe first that, for $0 \leq i \leq k-2$, parts (a) and (b) of Lemma 4.2 guarantee that the products $N_{i} W_{i} H_{i}$ are well defined and that the block partitions of $N_{i}, W_{i}$, and $H_{i}$ are conformal for this matrix product. Moreover, from Theorem 3.6, the block $W_{i}(1,1)$ has always size $m \times n$ and, so, the size of $W_{i}(2: i+2,2: i+2)$ is $\alpha_{i} \times \alpha_{i}$. Therefore, $\operatorname{diag}\left(\lambda P_{k-i-2}, \lambda I_{\alpha_{i}}\right)$ has the same size as $W_{i}$ and can be partitioned into blocks exactly in the same way as $W_{i}$, where recall that the diagonal blocks $W_{i}(2,2), \ldots, W_{i}(i+2, i+2)$ are square. As a consequence, also the products $N_{i} \operatorname{diag}\left(\lambda P_{k-i-2}, \lambda I_{\alpha_{i}}\right) H_{i}$ are well defined.

The rest of the proof consists in performing carefully block matrix products. We start with the proof of part (a) that has to be split in two cases.
Case Part (a)-(1): $\sigma$ has a consecution at $i$. Then
$N_{i}=\left[\begin{array}{cc}I_{m} & 0 \\ \lambda N_{i-1}(:, 1) & N_{i-1}\end{array}\right], \quad W_{i}=\left[\begin{array}{ccc}-A_{i+1} & I_{m} & 0 \\ W_{i-1}(:, 1) & 0 & W_{i-1}(:, 2: i+1)\end{array}\right], \quad H_{i}=\left[\begin{array}{cc}0 & H_{i-1}(1,:) \\ -I_{m} & P_{k-i-1} H_{i-1}(1,:) \\ 0 & H_{i-1}(2: i+1,:)\end{array}\right]$.
This implies

$$
N_{i}\left[\begin{array}{ll}
\lambda P_{k-i-2} &  \tag{27}\\
& \lambda I_{\alpha_{i}}
\end{array}\right] H_{i}=\left[\begin{array}{cc}
0 & \lambda P_{k-i-2} H_{i-1}(1,:) \\
-\lambda N_{i-1}(:, 1) & \left(S_{i}\right)_{22}
\end{array}\right],
$$

where

$$
\left(S_{i}\right)_{22}=N_{i-1}(:, 1)\left(\lambda^{2} P_{k-i-2}+\lambda P_{k-i-1}\right) H_{i-1}(1,:)+\lambda N_{i-1}(:, 2: i+1) H_{i-1}(2: i+1,:),
$$

and

$$
N_{i} W_{i} H_{i}=\left[\begin{array}{cc}
-I_{m} & \left(-A_{i+1}+P_{k-i-1}\right) H_{i-1}(1,:)  \tag{28}\\
-\lambda N_{i-1}(:, 1) & \left(T_{i}\right)_{22}
\end{array}\right]
$$

where

$$
\begin{aligned}
\left(T_{i}\right)_{22}= & N_{i-1}(:, 1)\left(-\lambda A_{i+1}+\lambda P_{k-i-1}\right) H_{i-1}(1,:) \\
& +N_{i-1}\left(W_{i-1}(:, 1) H_{i-1}(1,:)+W_{i-1}(:, 2: i+1) H_{i-1}(2: i+1,:)\right) .
\end{aligned}
$$

Now, we use (27), (28), and $-A_{i+1}+P_{k-i-1}=\lambda P_{k-i-2}$ to get

$$
N_{i}\left(\left[\begin{array}{cc}
\lambda P_{k-i-2} & \\
& \lambda I_{\alpha_{i}}
\end{array}\right]-W_{i}\right) H_{i}=\left[\begin{array}{cc}
I_{m} & 0 \\
0 & \left(Z_{i}\right)_{22}
\end{array}\right],
$$

where

$$
\begin{aligned}
\left(Z_{i}\right)_{22}= & \lambda N_{i-1}(:, 1) P_{k-i-1} H_{i-1}(1,:)+\lambda N_{i-1}(:, 2: i+1) H_{i-1}(2: i+1,:) \\
& -N_{i-1}\left(W_{i-1}(:, 1) H_{i-1}(1,:)+W_{i-1}(:, 2: i+1) H_{i-1}(2: i+1,:)\right) \\
= & N_{i-1}\left[\begin{array}{ll}
\lambda P_{k-i-1} & \\
& \lambda I_{\alpha_{i-1}}
\end{array}\right] H_{i-1}-N_{i-1} W_{i-1} H_{i-1} \\
= & N_{i-1}\left(\left[\begin{array}{ll}
\lambda P_{k-i-1} & \\
& \lambda I_{\alpha_{i-1}}
\end{array}\right]-W_{i-1}\right) H_{i-1} .
\end{aligned}
$$

Case Part (a)-(2): $\sigma$ has an inversion at $i$. In this case

$$
N_{i}\left[\begin{array}{ll}
\lambda P_{k-i-2} &  \tag{29}\\
& \lambda I_{\alpha_{i}}
\end{array}\right] H_{i}=\left[\begin{array}{cc}
0 & -\lambda H_{i-1}(1,:) \\
\lambda N_{i-1}(:, 1) P_{k-i-2} & \left(S_{i}\right)_{22}
\end{array}\right],
$$

where $\left(S_{i}\right)_{22}$ is the same as in case (a)-(1), and

$$
N_{i} W_{i} H_{i}=\left[\begin{array}{cc}
-I_{n} & -\lambda H_{i-1}(1,:)  \tag{30}\\
\lambda N_{i-1}(:, 1) P_{k-i-2} & \left(\widetilde{T}_{i}\right)_{22}
\end{array}\right],
$$

where now

$$
\begin{aligned}
\left(\widetilde{T}_{i}\right)_{22}= & N_{i-1}(:, 1)\left(-\lambda A_{i+1}+\lambda P_{k-i-1}\right) H_{i-1}(1,:) \\
& +\left(N_{i-1}(:, 1) W_{i-1}(1,:)+N_{i-1}(:, 2: i+1) W_{i-1}(2: i+1,:)\right) H_{i-1} .
\end{aligned}
$$

Subtracting (30) from (29) and reasoning as in case (a)-(1) we get again the desired identity.
The proof of Part (b) is again a direct block matrix product and is omitted. We only remark that one has to consider again separately two cases, that is, $\sigma$ has a consecution at 0 and $\sigma$ has an inversion at 0 , and to use that $\lambda P_{k-2}=P_{k-1}-A_{1}$ and $P(\lambda)=\lambda^{2} P_{k-2}+\lambda A_{1}+A_{0}$.

Now we are in the position of proving the main result in this section.
Theorem 4.5. Let $P(\lambda)$ be an $m \times n$ matrix polynomial with degree larger than or equal to 2 . Then any Fiedler companion pencil $F_{\sigma}(\lambda)$ of $P(\lambda)$ is a strong linearization for $P(\lambda)$.
Proof. We denote as usual $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ and adopt the notation used in Lemma 4.4. Moreover, recall from Definition 4.3 that $U_{\sigma}(\lambda)=N_{k-2}$ and $V_{\sigma}(\lambda)=H_{k-2}$, from Definition 3.8 that $M_{\sigma}=W_{k-2}$, that $P_{0}=A_{k}$, and that $\alpha_{k-2}=m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)$. Therefore, Part (a) in Lemma 4.4 for $i=k-2$ implies

$$
\begin{aligned}
U_{\sigma}(\lambda) F_{\sigma}(\lambda) V_{\sigma}(\lambda) & =N_{k-2}\left(\lambda\left[\begin{array}{ll}
P_{0} & \\
& I_{\alpha_{k-2}}
\end{array}\right]-W_{k-2}\right) H_{k-2} \\
& =\left[\begin{array}{lll}
I_{\beta_{k-2}} & \\
& N_{k-3}\left(\left[\begin{array}{ll}
\lambda P_{1} & \\
& \lambda I_{\alpha_{k-3}}
\end{array}\right]-W_{k-3}\right) H_{k-3}
\end{array}\right] .
\end{aligned}
$$

Now, apply Part (a) in Lemma 4.4 for $i=k-3$ to the lower right block in the previous equation to get

$$
U_{\sigma}(\lambda) F_{\sigma}(\lambda) V_{\sigma}(\lambda)=\left[\begin{array}{lll}
I_{\beta_{k-2}} & & \\
& I_{\beta_{k-3}} & \\
& & N_{k-4}\left(\left[\begin{array}{ll}
\lambda P_{2} & \\
& \lambda I_{\alpha_{k-4}}
\end{array}\right]-W_{k-4}\right) H_{k-4}
\end{array}\right]
$$

Next, we apply Part (a) in Lemma 4.4 consecutively for $i=k-4, k-5, \ldots, 1$ until we get

$$
U_{\sigma}(\lambda) F_{\sigma}(\lambda) V_{\sigma}(\lambda)=\left[\begin{array}{ll}
I_{\beta_{k-2}+\beta_{k-3}+\ldots+\beta_{1}} & \\
& N_{0}\left(\left[\begin{array}{ll}
\lambda P_{k-2} & \\
& \lambda I_{\alpha_{0}}
\end{array}\right]-W_{0}\right) H_{0}
\end{array}\right]
$$

We finally apply Part (b) in Lemma 4.4 and use $\beta_{k-2}+\cdots+\beta_{1}+\beta_{0}=\alpha_{k-2}$ to obtain

$$
U_{\sigma}(\lambda) F_{\sigma}(\lambda) V_{\sigma}(\lambda)=\left[\begin{array}{ll}
I_{m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)} &  \tag{31}\\
& P(\lambda)
\end{array}\right]
$$

which proves that $F_{\sigma}(\lambda)$ is a linearization of $P(\lambda)$, since $U_{\sigma}(\lambda)$ and $V_{\sigma}(\lambda)$ are unimodular.
To prove that $F_{\sigma}(\lambda)$ is a strong linearization of $P(\lambda)$, we invoke Theorem 3.14, which states that $\operatorname{rev} F_{\sigma}(P)$ is strictly equivalent to $F_{\tau}(\operatorname{rev} P)$, where $\tau$ is a bijection defined in the statement of Theorem 3.14 and that has the same total number of consecutions and the same total number of inversions as $\sigma$. We can apply (31) to $F_{\tau}(\operatorname{rev} P)$ and $\operatorname{rev} P$ to prove that $F_{\tau}(\operatorname{rev} P)$ is unimodularly equivalent to $\operatorname{diag}\left(I_{m \mathfrak{c}(\sigma)+n \mathrm{i}(\sigma)}, \operatorname{rev} P\right)$, and, therefore, $\operatorname{rev} F_{\sigma}(P)$ is unimodularly equivalent to $\operatorname{diag}\left(I_{m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)}, \operatorname{rev} P\right)$. This proves that $F_{\sigma}(\lambda)$ is indeed a strong linearization of $P(\lambda)$.

From the proof of Theorem 4.5, we get Corollary 4.6, which will be fundamental in Section 5.
Corollary 4.6. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, let $\sigma$ : $\{0, \ldots, k-1\} \rightarrow\{1, \ldots, k\}$ be a bijection, and let $U_{\sigma}(\lambda)$ and $V_{\sigma}(\lambda)$ be, respectively, the left and right unimodular equivalence matrices associated with $P(\lambda)$ and $\sigma$ introduced in Definition 4.3. Then

$$
U_{\sigma}(\lambda) F_{\sigma}(\lambda) V_{\sigma}(\lambda)=\left[\begin{array}{ll}
I_{m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)} & \\
& P(\lambda)
\end{array}\right] .
$$

## 5. Recovery of minimal indices and bases of rectangular polynomials from Fiedler pencils

In this section we show how to recover in a very simple way the minimal indices and bases of a rectangular matrix polynomial from those of any of its Fiedler pencils. The results and most of the proofs in this section are very similar to the ones presented for singular square matrix polynomials in Section 5 of [11]. Therefore, in order to avoid some repetitions, we omit all proofs that the reader can deduce easily from [11, Section 5] and pay close attention only to those arguments where the fact that the polynomial is rectangular makes a significative difference. Simultaneously, in order to keep the reading of the paper self-contained, we present careful statements of the main results and exact pointers to the results in [11] where the proofs can be found. In this sense, this section has a different character than Sections 3 and 4, where most proofs have been presented with detail since the approaches followed in Sections 3 and 4 are new and very different than those in [11].

The main recovery results in this section are Corollaries 5.4 and 5.7 , which are consequences of Theorems 5.3 and 5.6 , respectively. These results extend to rectangular matrix polynomials what was proved only for square singular polynomials in Corollaries 5.8 and 5.11 , and Theorems 5.7 and 5.9, in [11].

Corollaries 5.4 and 5.7 and Theorems 5.3 and 5.6 rely on a series of previous results. The first one is Lemma 5.1, which is an extension to rectangular matrix polynomials of [11, Lemma 5.1]. The proof is an obvious modification of the one of [11, Lemma 5.1] and is omitted.

Lemma 5.1. Let the pencil $L(\lambda)$ be a linearization of an $m \times n$ matrix polynomial $P(\lambda)$ of degree $k \geq 2$ and let $U(\lambda)$ and $V(\lambda)$ be two unimodular matrix polynomials such that

$$
U(\lambda) L(\lambda) V(\lambda)=\left[\begin{array}{cc}
I_{s} & 0 \\
0 & P(\lambda)
\end{array}\right] .
$$

Let $U^{L}=U^{L}(\lambda)$ and $V^{R}=V^{R}(\lambda)$ be, respectively, the matrices comprising the last $m$ rows of $U(\lambda)$ and the last $n$ columns of $V(\lambda)$. Then
(a) the linear map

$$
\begin{array}{lclc}
\mathcal{L}: & \mathcal{N}_{\ell}(P) & \longrightarrow & \mathcal{N}_{\ell}(L) \\
& w^{T}(\lambda) & \longmapsto & w^{T}(\lambda) \cdot U^{L}
\end{array}
$$

is an isomorphism of $\mathbb{F}(\lambda)$-vector spaces;
(b) the linear map

$$
\begin{array}{cccc}
\mathcal{R}: & \mathcal{N}_{r}(P) & \longrightarrow & \mathcal{N}_{r}(L) \\
& v(\lambda) & \longmapsto & V^{R} \cdot v(\lambda)
\end{array}
$$

is an isomorphism of $\mathbb{F}(\lambda)$-vector spaces.
A consequence of Lemma 5.1 is that the bases of $\mathcal{N}_{r}(P)$ and the ones of $\mathcal{N}_{r}(L)$ are in one-to-one correspondence through the map $\mathcal{R}$. But, for an arbitrary linearization $L(\lambda)$, the map $\mathcal{R}$ does not necessarily establish a one-to-one correspondence between minimal bases. A key point in our developments is that we will prove that for each Fiedler pencil $F_{\sigma}(\lambda)$ of an $m \times n$ matrix polynomial $P(\lambda)$, if $V_{\sigma}(\lambda)$ is the right unimodular equivalence matrix appearing in Corollary 4.6 , then the map $\mathcal{R}_{\sigma}$ associated with $V_{\sigma}(\lambda)$ provides actually a one-to-one correspondence between the right minimal bases of $P(\lambda)$ and those of $F_{\sigma}(\lambda)$. This correspondence is so simple that allows us to obtain very easily the right minimal bases of $P(\lambda)$ from the right minimal bases of $F_{\sigma}(\lambda)$. Analogous results hold for left minimal indices and bases.

Lemma 5.1 and Corollary 4.6 indicate that we need to determine the last block column of the matrix $V_{\sigma}(\lambda)=H_{k-2}$ introduced in Definition 4.3, since this last block column has precisely $n$ columns according to Lemma 4.2-(f). For this purpose, we need to define as in [11] some additional magnitudes and matrices, which are based on the consecution-inversion structure sequence of $\sigma$ introduced in Definition 3.2, i.e., $\operatorname{CISS}(\sigma)=\left(c_{1}, i_{1}, \ldots, c_{\ell}, i_{\ell}\right)$. So, we define

$$
\begin{align*}
s_{0} & :=0, \quad s_{j}:=\sum_{p=1}^{j}\left(c_{p}+i_{p}\right) \quad \text { for } j=1, \ldots, \ell,  \tag{32}\\
m_{0} & :=0, \quad m_{j}:=\sum_{p=1}^{j} i_{p} \quad \text { for } j=1, \ldots, \ell \tag{33}
\end{align*}
$$

Recall that $s_{\ell}=k-1$ and $m_{\ell}=\mathfrak{i}(\sigma)$. We also need to define some matrices, denoted $\Lambda_{\sigma, j}(P)$ for $j=1, \ldots, \ell$ and $\widehat{\Lambda}_{\sigma, j}(P)$ for $j=1, \ldots, \ell-1$, associated with the $m \times n$ matrix polynomial $P(\lambda)$ and the bijection $\sigma$. These matrices are defined in terms of the Horner shifts of $P(\lambda)$ as follows:

$$
\Lambda_{\sigma, j}(P):=\left[\begin{array}{c}
\lambda^{i_{j}} I_{n}  \tag{34}\\
\vdots \\
\lambda I_{n} \\
I_{n} \\
P_{k-s_{j-1}-c_{j}} \\
\vdots \\
P_{k-s_{j-1}-2} \\
P_{k-s_{j-1}-1}
\end{array}\right] \quad \text { and } \quad \widehat{\Lambda}_{\sigma, j}(P):=\left[\begin{array}{c}
\lambda^{i_{j}-1} I_{n} \\
\vdots \\
\lambda I_{n} \\
I_{n} \\
P_{k-s_{j-1}-c_{j}} \\
\vdots \\
P_{k-s_{j-1}-2} \\
P_{k-s_{j-1}-1}
\end{array}\right] \quad \text { if } c_{1} \geq 1,
$$

but if $c_{1}=0$, then $\Lambda_{\sigma, 1}(P):=\left[\lambda^{i_{1}} I_{n}, \ldots, \lambda I_{n}, I_{n}\right]^{T}, \widehat{\Lambda}_{\sigma, 1}(P):=\left[\lambda^{i_{1}-1} I_{n}, \ldots, \lambda I_{n}, I_{n}\right]^{T}$, with $\Lambda_{\sigma, j}(P)$, $\widehat{\Lambda}_{\sigma, j}(P)$ as in (34) for $j>1$. Here for simplicity we omit $\lambda$ from the Horner shifts $P_{d}(\lambda)$. With all these
definitions we are in the position of stating and proving Lemma 5.2, which describes explicitly the last block-column of $V_{\sigma}(\lambda)$. Note that Lemma 5.2 generalizes to rectangular polynomials [11, Lemma 5.3]. However, the proof is completely different than that of [11, Lemma 5.3].
Lemma 5.2. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, let $F_{\sigma}(\lambda)$ be the Fiedler pencil of $P(\lambda)$ associated with the bijection $\sigma$, and let $V_{\sigma}(\lambda)$ be the right unimodular equivalence matrix associated with $P(\lambda)$ and $\sigma$ introduced in Definition 4.3. Consider $V_{\sigma}(\lambda)$ partitioned into $k \times k$ blocks according to Algorithm 5. Then the last block-column $V^{R}(\lambda)$ of $V_{\sigma}(\lambda)$, i.e., the last $n$ columns of $V_{\sigma}(\lambda)$, is

$$
\Lambda_{\sigma}^{R}(P):=\left[\begin{array}{c}
\lambda^{m_{\ell-1}} \Lambda_{\sigma, \ell}(P)  \tag{35}\\
\lambda^{m_{\ell-2}} \widehat{\Lambda}_{\sigma, \ell-1}(P) \\
\vdots \\
\lambda^{m_{1}} \widehat{\Lambda}_{\sigma, 2}(P) \\
\widehat{\Lambda}_{\sigma, 1}(P)
\end{array}\right] \quad \text { if } \ell>1
$$

and $V^{R}(\lambda)=\Lambda_{\sigma, 1}(P)=: \Lambda_{\sigma}^{R}(P)$ if $\ell=1$.
Proof. The last block-column of $V_{\sigma}(\lambda)=H_{k-2}$ can be determined by Algorithm 5 just by looking at the last block-column at each step of the algorithm. Set CISS $(\sigma)=\left(c_{1}, i_{1}, \ldots, c_{\ell}, i_{\ell}\right)$. Assume first $c_{1}>0$, which means that $\sigma$ has consecutions at $0,1, \ldots, c_{1}-1$ and, as a consequence, the last block-column of the matrix $H_{c_{1}-1}$ constructed after the steps $0,1, \ldots, c_{1}-1$ of Algorithm 5 is of the form

$$
H_{c_{1}-1}\left(:, c_{1}+1\right)=\left[\begin{array}{c}
I_{n} \\
P_{k-c_{1}} \\
\vdots \\
P_{k-1}
\end{array}\right]
$$

Now, $\sigma$ has inversions at $c_{1}, c_{1}+1, \ldots, c_{1}+i_{1}-1$ and then the last block-column of $H_{c_{1}+i_{1}-1}$ is

$$
H_{c_{1}+i_{1}-1}\left(:, c_{1}+i_{1}+1\right)=\left[\begin{array}{c}
\lambda^{i_{1}} I_{n} \\
\vdots \\
\lambda I_{n} \\
I_{n} \\
P_{k-c_{1}} \\
\vdots \\
P_{k-1}
\end{array}\right]=\Lambda_{\sigma, 1}(P)
$$

The reader may check that if $c_{1}=0$, then $H_{i_{1}-1}\left(:, i_{1}+1\right)=\left[\lambda^{i_{1}} I_{n}, \ldots, \lambda I_{n}, I_{n}\right]^{T}=\Lambda_{\sigma, 1}(P)$. The proof finishes here if $\ell=1$, because in this case $c_{1}+i_{1}=k-1$ and $H_{c_{1}+i_{1}-1}\left(:, c_{1}+i_{1}+1\right)=H_{k-2}(:, k)$ is the last block column of $V_{\sigma}(\lambda)$.

If $\ell>1$, then the following $c_{2}$ consecutions of $\sigma$ at $c_{1}+i_{1}, c_{i}+i_{1}+1, \ldots, c_{1}+i_{1}+c_{2}-1$ give, according to Algorithm 5,

$$
H_{c_{1}+i_{1}+c_{2}-1}\left(:, c_{1}+i_{1}+c_{2}+1\right)=\left[\begin{array}{c}
\lambda^{i_{1}} I_{n} \\
\lambda^{i_{1}} P_{k-s_{1}-c_{2}} \\
\vdots \\
\lambda^{i_{1}} P_{k-s_{1}-1} \\
\lambda^{i_{1}-1} I_{n} \\
\vdots \\
\lambda I_{n} \\
I_{n} \\
P_{k-c_{1}} \\
\vdots \\
P_{k-1}
\end{array}\right]=\left[\begin{array}{c}
\lambda^{i_{1}} I_{n} \\
\lambda^{i_{1}} P_{k-s_{1}-c_{2}} \\
\vdots \\
\lambda^{i_{1}} P_{k-s_{1}-1} \\
\widehat{\Lambda}_{\sigma, 1}(P)
\end{array}\right] .
$$

Notice that this produces the "truncated" block matrix $\widehat{\Lambda}_{\sigma, 1}(P)$ at the bottom of the last block-column of $H_{c_{1}+i_{1}+c_{2}-1}$. The following $i_{2}$ inversions of $\sigma$ produce

$$
H_{s_{2}-1}\left(:, s_{2}+1\right)=\left[\begin{array}{c}
\lambda^{m_{1}} \Lambda_{\sigma, 2}(P) \\
\widehat{\Lambda}_{\sigma, 1}(P)
\end{array}\right] .
$$

The rest of the proof follows by induction with arguments similar to the ones above. Assume that for $j$ such that $2 \leq j<\ell$ we have

$$
H_{s_{j}-1}\left(:, s_{j}+1\right)=\left[\begin{array}{c}
\lambda^{m_{j-1}} \Lambda_{\sigma, j}(P) \\
\lambda^{m_{j-2}} \widehat{\Lambda}_{\sigma, j-1}(P) \\
\vdots \\
\lambda^{m_{1}} \widehat{\Lambda}_{\sigma, 2}(P) \\
\widehat{\Lambda}_{\sigma, 1}(P)
\end{array}\right]
$$

and prove via Algorithm 5 that the corresponding result holds for $j+1$. Note, to finish the proof, that $H_{s_{\ell}-1}\left(:, s_{\ell}+1\right)=H_{k-2}(:, k)$ is precisely the last block column of $V_{\sigma}(\lambda)$.

A fundamental fact in Lemma 5.2 is that $\Lambda_{\sigma}^{R}(P)$ always has exactly one block equal to $I_{n}$ at block index $k-c_{1}$. This is the block that allows us to recover very easily the minimal bases of $P(\lambda)$ from those of $F_{\sigma}(\lambda)$. To this purpose we establish first Theorem 5.3, whose proof is omitted since it is essentially the same as the one of [11, Theorem 5.7].
Theorem 5.3. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, let $F_{\sigma}(\lambda)$ be the Fiedler pencil of $P(\lambda)$ associated with a bijection $\sigma$, let $\mathfrak{i}(\sigma)$ be the total number of inversions of $\sigma$, let $\mathfrak{c}(\sigma)$ be the total number of consecutions of $\sigma$, and let $\Lambda_{\sigma}^{R}(P)$ be the $(n+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)) \times n$ matrix defined in (35). Then the linear map

$$
\begin{aligned}
\mathcal{R}_{\sigma}: \mathcal{N}_{r}(P) & \longrightarrow \mathcal{N}_{r}\left(F_{\sigma}\right) \\
v & \longmapsto \Lambda_{\sigma}^{R}(P) v
\end{aligned}
$$

is an isomorphism of $\mathbb{F}(\lambda)$-vector spaces with uniform degree-shift $\mathfrak{i}(\sigma)$ on the vector polynomials in $\mathcal{N}_{r}(P)$. More precisely, $\mathcal{R}_{\sigma}$ induces a bijection between the subsets of vector polynomials in $\mathcal{N}_{r}(P)$ and $\mathcal{N}_{r}\left(F_{\sigma}\right)$, with the property that

$$
\operatorname{deg} \mathcal{R}_{\sigma}(v)=\mathfrak{i}(\sigma)+\operatorname{deg} v
$$

for every nonzero vector polynomial $v \in \mathcal{N}_{r}(P)$. Furthermore, for any nonzero vector polynomial $v$, $\operatorname{deg} \mathcal{R}_{\sigma}(v)$ is attained only in the topmost $n \times 1$ block of $\mathcal{R}_{\sigma}(v)$.

An immediate consequence of Theorem 5.3 is Corollary 5.4, that establishes a very simple relationship between the right minimal indices and bases of $P(\lambda)$ and $F_{\sigma}(\lambda)$. The proof of Corollary 5.4 is also omitted since it is almost the same as the one of [11, Corollary 5.8].

Corollary 5.4 (recovery of right minimal indices and bases). Let $P(\lambda)$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, and let $F_{\sigma}(\lambda)$ be the Fiedler pencil of $P(\lambda)$ associated with a bijection $\sigma$ having $\operatorname{CISS}(\sigma)=\left(c_{1}, i_{1}, \ldots, c_{\ell}, i_{\ell}\right)$ and total number of consecutions and inversions $\mathfrak{c}(\sigma)$ and $\mathfrak{i}(\sigma)$, respectively. Suppose that each vector $z(\lambda) \in \mathcal{N}_{r}\left(F_{\sigma}\right) \subset \mathbb{F}(\lambda)^{(n+m \mathfrak{c}(\sigma)+n \mathbf{i}(\sigma)) \times 1}$ is partitioned into $k \times 1$ blocks which are conformal for multiplication with the partition of $F_{\sigma}(\lambda)$ given by Algorithm 2.
(a) If $z(\lambda) \in \mathcal{N}_{r}\left(F_{\sigma}\right)$, and $x(\lambda) \in \mathbb{F}(\lambda)^{n \times 1}$ is the $\left(k-c_{1}\right)$ th block of $z(\lambda)$, then $x(\lambda) \in \mathcal{N}_{r}(P)$.
(b) If $\left\{z_{1}(\lambda), \ldots, z_{p}(\lambda)\right\}$ is a right minimal basis of $F_{\sigma}(\lambda)$, and $x_{j}(\lambda)$ is the $\left(k-c_{1}\right)$ th block of $z_{j}(\lambda)$ for each $j=1, \ldots, p$, then $\left\{x_{1}(\lambda), \ldots, x_{p}(\lambda)\right\}$ is a right minimal basis of $P(\lambda)$.
(c) If $0 \leq \varepsilon_{1} \leq \varepsilon_{2} \leq \cdots \leq \varepsilon_{p}$ are the right minimal indices of $P(\lambda)$, then

$$
\varepsilon_{1}+\mathfrak{i}(\sigma) \leq \varepsilon_{2}+\mathfrak{i}(\sigma) \leq \cdots \leq \varepsilon_{p}+\mathfrak{i}(\sigma)
$$

are the right minimal indices of $F_{\sigma}(\lambda)$.

Note that these results hold for the first companion form of $P(\lambda)$ by taking $c_{1}=0$ and $\mathfrak{i}(\sigma)=k-1$, and for the second companion form using $c_{1}=k-1$ and $\mathfrak{i}(\sigma)=0$.

For the recovery of left minimal indices and bases, it is possible to take a similar approach to the one we have used for right minimal indices and bases, that is, according to Lemma 5.1 and Corollary 4.6, the last $m$ rows of $U_{\sigma}(\lambda)$ can be determined via Algorithm 4. However, we follow here a different way, based on the fact that the left minimal indices and bases of a matrix polynomial (and in particular, of a matrix pencil) coincide with the right minimal indices and bases of its transpose, since $y(\lambda)^{T} \in \mathcal{N}_{\ell}(P)$ if and only if $y(\lambda) \in \mathcal{N}_{r}\left(P^{T}\right)$. Then the idea is to relate the right minimal indices and bases of $P(\lambda)^{T}$ with the ones of $F_{\sigma}(\lambda)^{T}$. The advantage of using this approach relies on the fact that $F_{\sigma}(\lambda)^{T}$ is again a Fiedler pencil for $P(\lambda)^{T}$ as the following lemma shows.
Lemma 5.5. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial of degree $k \geq 2$ and $\sigma:\{0, \ldots, k-$ $1\} \rightarrow\{1, \ldots, k\}$ be a bijection. Define the reversal bijection of $\sigma$ as follows: $\operatorname{rev} \sigma(i):=k+1-\sigma(i)$ for $i=0,1, \ldots, k-1$. Then

$$
\left[F_{\sigma}(P)\right]^{T}=F_{\mathrm{rev} \sigma}\left(P^{T}\right)
$$

Proof. This lemma can be easily proved by induction using Algorithm 2. Let $\left\{W_{i}\right\}_{i=0}^{k-2}$ be the sequence of matrices constructed by Algorithm 2 for $P(\lambda)$ and $\sigma$, and let $\left\{W_{i}^{\prime}\right\}_{i=0}^{k-2}$ be the sequence of matrices constructed by Algorithm 2 for $P(\lambda)^{T}=\sum_{i=0}^{k} \lambda^{i} A_{i}^{T}$ and rev $\sigma$. Note also that rev $\sigma$ has a consecution (resp. inversion) at $i$ if and only if $\sigma$ has an inversion (resp. consecution) at $i$. First notice that

$$
W_{0}^{T}=\left[\begin{array}{cc}
-A_{1}^{T} & -A_{0}^{T} \\
I_{m} & 0
\end{array}\right] \quad \text { or } \quad W_{0}^{T}=\left[\begin{array}{cc}
-A_{1}^{T} & I_{n} \\
-A_{0}^{T} & 0
\end{array}\right]
$$

depending on whether $\sigma$ has a consecution or an inversion at 0 . Notice that, in both cases, we get $W_{0}^{T}=W_{0}^{\prime}$. Now, we proceed by induction: assume $W_{i-1}^{T}=W_{i-1}^{\prime}$ for some $0 \leq(i-1)<k-2$, and prove that $W_{i}^{T}=W_{i}^{\prime}$. For this purpose, use Algorithm 2 to see that

$$
W_{i}^{T}=\left[\begin{array}{cc}
-A_{i+1}^{T} & W_{i-1}(:, 1)^{T} \\
I_{m} & 0 \\
0 & W_{i-1}(:, 2: i+1)^{T}
\end{array}\right] \quad \text { or } \quad W_{i}^{T}=\left[\begin{array}{ccc}
-A_{i+1}^{T} & I_{n} & 0 \\
W_{i-1}(1,:)^{T} & 0 & W_{i-1}(2: i+1,:)^{T}
\end{array}\right]
$$

depending on whether $\sigma$ has a consecution or an inversion at $i$. By the induction hypothesis, this is precisely the matrix $W_{i}^{\prime}$. The statement of Lemma 5.5 follows from $W_{k-2}^{T}=W_{k-2}^{\prime}$.

Lemma 5.5 allows us to prove Theorem 5.6 essentially in the same way as Theorem 5.9 in [11] was proved. Therefore, we omit the proof of Theorem 5.6, although we remark that we cannot use here the block-transpose operation $(\cdot)^{\mathcal{B}}$, see [11, Definition 3.6], because the blocks of $\Lambda_{\mathrm{rev} \sigma}^{R}\left(P^{T}\right)$ do not have all the same sizes when $P(\lambda)$ is rectangular. This motivates a minor modification ${ }^{2}$ in the statement of Theorem 5.6 with respect the statement of [11, Theorem 5.9]. Note also that $\mathfrak{i}(\operatorname{rev} \sigma)=\mathfrak{c}(\sigma)$ and $\mathfrak{c}(\operatorname{rev} \sigma)=\mathfrak{i}(\sigma)$.
Theorem 5.6. Let $P(\lambda)=\sum_{i=0}^{k} \lambda^{i} A_{i}$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, let $F_{\sigma}(\lambda)$ be the Fiedler pencil of $P(\lambda)$ associated with a bijection $\sigma$, let $\mathfrak{c}(\sigma)$ be the total number of consecutions of $\sigma$, let $\mathfrak{i}(\sigma)$ be the total number of inversions of $\sigma$, and let $\Lambda_{\mathrm{rev} ~}^{R}\left(P^{T}\right)$ be, for the $n \times m$ polynomial $P(\lambda)^{T}$ and the reversal bijection $\operatorname{rev} \sigma$, the $(m+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma)) \times m$ matrix defined in Lemma 5.2. Then the linear map

$$
\begin{aligned}
& \mathcal{L}_{\sigma}: \mathcal{N}_{\ell}(P) \longrightarrow \\
& u^{T} \longmapsto \mathcal{N}_{\ell}\left(F_{\sigma}\right) \\
& u^{T} \Lambda_{\sigma}^{L}(P),
\end{aligned}
$$

where $\Lambda_{\sigma}^{L}(P):=\left[\Lambda_{\mathrm{rev} \sigma}^{R}\left(P^{T}\right)\right]^{T}$, is an isomorphism of $\mathbb{F}(\lambda)$-vector spaces with uniform degree-shift $\mathfrak{c}(\sigma)$ on the vector polynomials in $\mathcal{N}_{\ell}(P)$. More precisely, $\mathcal{L}_{\sigma}$ induces a bijection between the subsets of vector polynomials in $\mathcal{N}_{\ell}(P)$ and $\mathcal{N}_{\ell}\left(F_{\sigma}\right)$, with the property that

$$
\begin{equation*}
\operatorname{deg} \mathcal{L}_{\sigma}\left(u^{T}\right)=\mathfrak{c}(\sigma)+\operatorname{deg}\left(u^{T}\right) \tag{36}
\end{equation*}
$$

[^2]for every nonzero vector polynomial $u^{T} \in \mathcal{N}_{\ell}(P)$. Furthermore, for any nonzero vector polynomial $u^{T}$, $\operatorname{deg} \mathcal{L}_{\sigma}\left(u^{T}\right)$ is attained only in the leftmost $1 \times m$ block of $\mathcal{L}_{\sigma}\left(u^{T}\right)$.

An immediate consequence of Theorem 5.6 is Corollary 5.7 , which establishes a very simple relationship between the left minimal indices and bases of $P(\lambda)$ and $F_{\sigma}(\lambda)$. The easy proof is also omitted. We only indicate that the fact "rev $\sigma$ has a consecution (resp. inversion) at $i$ if and only if $\sigma$ has an inversion (resp. consecution) at $i$ " implies that $\Lambda_{\mathrm{rev} \sigma}^{R}\left(P^{T}\right)$ has exactly one block equal to $I_{m}$ at block index $k$ if $c_{1}>0$ and at block index $k-i_{1}$ if $c_{1}=0$.

Corollary 5.7 (recovery of left minimal indices and bases). Let $P(\lambda)$ be an $m \times n$ matrix polynomial with degree $k \geq 2$, and let $F_{\sigma}(\lambda)$ be the Fiedler pencil of $P(\lambda)$ associated with a bijection $\sigma$ having $\operatorname{CISS}(\sigma)=\left(c_{1}, i_{1}, \ldots, c_{\ell}, i_{\ell}\right)$ and total number of consecutions and inversions $\mathfrak{c}(\sigma)$ and $\mathfrak{i}(\sigma)$, respectively. Suppose that each vector $z(\lambda)^{T} \in \mathcal{N} \ell\left(F_{\sigma}\right) \subset \mathbb{F}(\lambda)^{1 \times(m+m \mathfrak{c}(\sigma)+n \mathfrak{i}(\sigma))}$ is partitioned into $1 \times k$ blocks which are conformal for multiplication with the partition of $F_{\sigma}(\lambda)$ given by Algorithm 2.
(a) If $z(\lambda)^{T} \in \mathcal{N}_{\ell}\left(F_{\sigma}\right)$, and

$$
y(\lambda)^{T} \text { is the }\left\{\begin{array}{cl}
k \text { th block of } z(\lambda)^{T} & \text { if } c_{1}>0 \\
\left(k-i_{1}\right) \text { th block of } z(\lambda)^{T} & \text { if } c_{1}=0
\end{array}\right.
$$

then $y(\lambda)^{T} \in \mathcal{N}_{\ell}(P)$.
(b) If $\left\{z_{1}(\lambda)^{T}, \ldots, z_{q}(\lambda)^{T}\right\}$ is a left minimal basis of $F_{\sigma}(\lambda)$, and

$$
y_{j}(\lambda)^{T} \text { is the }\left\{\begin{array}{cl}
k \text { th block of } z_{j}(\lambda)^{T} & \text { if } c_{1}>0 \\
\left(k-i_{1}\right) \text { th block of } z_{j}(\lambda)^{T} & \text { if } c_{1}=0
\end{array}\right.
$$

for $j=1, \ldots, q$, then $\left\{y_{1}(\lambda)^{T}, \ldots, y_{q}(\lambda)^{T}\right\}$ is a left minimal basis of $P(\lambda)$.
(c) If $0 \leq \eta_{1} \leq \eta_{2} \leq \cdots \leq \eta_{q}$ are the left minimal indices of $P(\lambda)$, then

$$
\eta_{1}+\mathfrak{c}(\sigma) \leq \eta_{2}+\mathfrak{c}(\sigma) \leq \cdots \leq \eta_{q}+\mathfrak{c}(\sigma),
$$

are the left minimal indices of $F_{\sigma}(\lambda)$.
Note that these results hold for the first companion form of $P(\lambda)$ using $\left(c_{1}, i_{1}\right)=(0, k-1)$ and $\mathfrak{c}(\sigma)=0$, and for the second companion form using $\left(c_{1}, i_{1}\right)=(k-1,0)$ and $\mathfrak{c}(\sigma)=k-1$.

Next we include an example that illustrates the results presented in this section. This example extends to rectangular matrix polynomials what appears in [11, Example 5.12] only for square singular polynomials, which allows the reader to appreciate the strong similarities and the really minor differences between square and rectangular polynomials.

Example 5.8. Let us consider an $m \times n$ matrix polynomial $P(\lambda)=\sum_{i=0}^{6} \lambda^{i} A_{i}$ with degree 6 and the Fiedler pencil $F_{\tau}(\lambda)$ of $P(\lambda)$ associated with the bijection $\tau=(1,2,5,3,6,4)$. Recall that the zero degree term $M_{\tau}$ of this pencil was considered in (16) and so

$$
F_{\tau}(\lambda)=\lambda \operatorname{diag}\left(A_{6}, I_{n}, I_{m}, I_{n}, I_{m}, I_{m}\right)-M_{\tau}
$$

Observe that $\operatorname{CISS}(\tau)=(2,1,1,1)$. So, for $\tau$, the parameters in (32)-(33) are $\ell=2, s_{\ell-1}=s_{1}=3$, and $m_{\ell-1}=m_{1}=1$. In addition, $\operatorname{rev} \tau=(6,5,2,4,1,3)$, hence $\operatorname{CISS}(\operatorname{rev} \tau)=(0,2,1,1,1,0)$, and, for $\operatorname{rev} \tau$, $\ell=3, s_{1}=2, s_{\ell-1}=s_{2}=4$, and $m_{1}=2, m_{\ell-1}=m_{2}=3$. Therefore

$$
\left.\begin{array}{rl}
\Lambda_{\tau}^{L}(P)=\left[\begin{array}{ll}
\Lambda_{\mathrm{rev} \tau}^{R}\left(P^{T}\right)
\end{array}\right]^{T} & =\left[\begin{array}{lllll}
\lambda^{3} I_{m} & \lambda^{3} P_{1}(\lambda) & \lambda^{2} I_{m} & \lambda^{2} P_{3}(\lambda) & \lambda I_{m}
\end{array} I_{m}\right.
\end{array}\right] .
$$

The relationships between the minimal indices and bases of $F_{\tau}(\lambda)$ and those of $P(\lambda)$ may now be summarized as follows:

- Right minimal indices of $F_{\tau}(\lambda)$ are shifted from those of $P(\lambda)$ by $\mathfrak{i}(\tau)=2$.
- Left minimal indices of $F_{\tau}(\lambda)$ are shifted from those of $P(\lambda)$ by $\mathfrak{c}(\tau)=3$.
- A right minimal basis of $P(\lambda)$ is recovered from the 4 th $=\left(k-c_{1}\right)$ th blocks (of size $n \times 1$ ) of any right minimal basis of $F_{\tau}(\lambda)$.
- A left minimal basis of $P(\lambda)$ is recovered from the 6 th $=k$ th blocks (of size $1 \times m$ ) of any left minimal basis of $F_{\tau}(\lambda)$.


## 6. Conclusions and future work

In the last decade several new classes of linearizations for square matrix polynomials have been introduced by various authors $[1,2,11,12,23,27,28,34]$. Among them, the class of Fiedler companion linearizations, which includes the classical first and second Frobenius companion forms, is a privileged class as a consequence of possessing the many valuable properties described in the Introduction. In this paper, we have extended Fiedler linearizations from square to rectangular matrix polynomials. To achieve this we have followed a completely different approach than the one followed in $[2,11]$ for regular and singular square polynomials, which cannot be easily generalized to the rectangular case. This new approach is based on a constructive definition via Algorithm 2, and has allowed us to prove that Fiedler pencils of rectangular matrix polynomials satisfy the same properties as Fiedler pencils of square matrix polynomials. More precisely, we have proved that every Fiedler pencil of a given rectangular polynomial $P(\lambda)$ is always a strong linearization for $P(\lambda)$, and that Fiedler pencils of rectangular matrix polynomials allow us to recover minimal indices and bases of matrix polynomials with essentially the same extremely simple rules as for Fiedler pencils of square polynomials. As far as we know, the class of Fiedler linearizations is the first of the new classes of linearizations introduced in the last decade that has been extended from square to rectangular polynomials. The most natural open problem in this context is to try to extend other classes of linearizations from square to rectangular matrix polynomials, e.g., the classes related to Fiedler pencils considered in $[2,5,12,34]$, or the vector spaces of linearizations introduced in [27]. Investigating the possibility of such extensions will be the subject of future work.

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[^1]:    ${ }^{1}$ Polynomials with large degrees may appear, for instance, in the computation of the roots of scalar polynomials as the eigenvalues of a Fiedler pencil [14].

[^2]:    ${ }^{2}$ In $\left[11\right.$, Theorem 5.9] the matrix $\left[\Lambda_{\mathrm{rev} \sigma}^{R}(P)\right]^{\mathcal{B}}$ was used, while in Theorem 5.6 we use $\left[\Lambda_{\mathrm{rev} \sigma}^{R}\left(P^{T}\right)\right]^{T}$. Note that both expressions coincide for square matrix polynomials, but that $\left[\Lambda_{\mathrm{rev} \sigma}^{R}(P)\right]^{\mathcal{B}}$ is not defined for rectangular polynomials.

