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## Recommended Citation

Rodman, L., \& Spitkovsky, I. M. (2012). Compressions of linearly independent selfadjoint operators. Linear Algebra and its Applications, 436(9), 3757-3766.

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# Compressions of linearly independent selfadjoint operators 

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## ARTICLEINFO

## Article history:

Received 24 October 2011
Accepted 31 October 2011
Available online 3 December 2011
Submitted by Peter Šemrl

## AMS classification:

47A12
Keywords:
Generalized numerical range
Selfadjoint operators
Compressions
Linear independence


#### Abstract

The following question is considered: What is the smallest number $\gamma(k)$ with the property that for every family $\left\{X_{1}, \ldots, X_{k}\right\}$ of $k$ selfadjoint and linearly independent operators on a real or complex Hilbert space $\mathcal{H}$ there exists a subspace $\mathcal{H}_{0} \subset \mathcal{H}$ of dimension $\gamma(k)$ such that the compressions of $X_{1}, \ldots, X_{k}$ to $\mathcal{H}_{0}$ are still linearly independent? Upper and lower bounds for $\gamma(k)$ are established for any $k$, and the exact value is found for $k=2,3$. It is also shown that the set of all $\gamma(k)$-dimensional subspaces $\mathcal{H}_{0}$ with the desired property is open and dense in the respective Grassmannian. The $k=3$ case is used to prove that the ratio numerical range $W(A / B)$ of a pair of operators on a Hilbert space either has a non-empty interior, or lies in a line or a circle.


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## 1. General problem and statements of main results

Let $\mathcal{H}$ be a Hilbert space over $\mathbb{F}$ with the inner product $\langle\cdot, \cdot\rangle$, where $\mathbb{F}$ stands for either the filed $\mathbb{R}$ of real or $\mathbb{C}$ of complex numbers. Denote by $L(\mathcal{H})$ the Banach algebra of linear bounded operators on $\mathcal{H}$, and let $X\left[\mathcal{H}_{0}\right]$ stand for the compression of a selfadjoint operator $X \in L(\mathcal{H})$ onto a (closed) subspace $\mathcal{H}_{0}$ of $\mathcal{H}$.

For every $k=1,2, \ldots$, define the positive integer $\gamma_{\mathbb{F}}(k)$ as follows: For every linearly independent $k$-tuple of bounded selfadjoint operators $X_{1}, \ldots, X_{k}$ on $\mathcal{H}$, there is a subspace $\mathcal{H}_{0}$ of $\mathcal{H}$ of dimension $\gamma_{\mathbb{F}}(k)$ such that the compressions $X_{1}\left[\mathcal{H}_{0}\right], \ldots, X_{k}\left[\mathcal{H}_{0}\right]$ are linearly independent; and there exists a linearly independent $k$-tuple of bounded selfadjoint operators $X_{1}^{\prime}, \ldots, X_{k}^{\prime}$ on $\mathcal{H}$ such that the compressions $X_{1}^{\prime}\left[\mathcal{H}_{0}^{\prime}\right], \ldots, X_{k}^{\prime}\left[\mathcal{H}_{0}^{\prime}\right]$ are linearly dependent for every subspace $\mathcal{H}_{0}^{\prime}$ of $\mathcal{H}$ of dimension $\gamma_{\mathbb{F}}(k)-1$.

[^0]For example, it is easy to see that

$$
\begin{equation*}
\gamma_{\mathbb{C}}(1)=\gamma_{\mathbb{R}}(1)=1 \tag{1.1}
\end{equation*}
$$

In this paper we address the problem of identifying the integers $\gamma_{\mathbb{C}}(k), \gamma_{\mathbb{R}}(k)$. We came to this problem from studies of generalized numerical ranges for Hilbert space operators [1-4]. It seems to us however that the problem is independently interesting.

Bounds on $\gamma_{\mathbb{C}}(k), \gamma_{\mathbb{R}}(k)$ can be easily given:
Proposition 1. For $k=1,2, \ldots$,

$$
2 k-1 \geqslant \gamma_{\mathbb{F}}(k) \geqslant \begin{cases}\left\lceil\frac{k+1}{2}\right\rceil & \text { if } \mathbb{F}=\mathbb{C}  \tag{1.2}\\ k & \text { if } \mathbb{F}=\mathbb{R}\end{cases}
$$

where $\lceil x\rceil$ stands for the smallest integer greater than or equal to $x$.
Proof. Consider the linearly independent $k$-tuple of selfadjoint operators $X_{1}, \ldots, X_{k} \in L\left(\mathbb{F}^{k}\right)$ given in the following matrix form:

$$
\begin{aligned}
& X_{1}=\left[\begin{array}{cc}
1 & 0_{1 \times(k-1)} \\
0_{(k-1) \times 1} & 0_{(k-1) \times(k-1)}
\end{array}\right], \quad X_{2}=\left[\begin{array}{cc}
0 & e_{1}^{T} \\
e_{1} & 0_{(k-1) \times(k-1)}
\end{array}\right], \\
& \ldots, X_{k}=\left[\begin{array}{cc}
0 & e_{k-1}^{T} \\
e_{k-1} & 0_{(k-1) \times(k-1)}
\end{array}\right],
\end{aligned}
$$

where $e_{1}, \ldots, e_{k-1}$ are the standard unit coordinate vectors. Any $\ell$-dimensional subspace $\mathcal{H}_{0}$ has intersection with $\operatorname{Span}\left\{e_{2}, \ldots, e_{k}\right\} \subset \mathbb{F}^{k}$ of dimension at least $\ell-1$. Thus, in a suitable orthonormal basis for $\mathcal{H}_{0}$, the compressions $\widehat{X}_{j}=X_{j}\left[\mathcal{H}_{0}\right]$ have the matrix form

$$
\widehat{X}_{j}=\left[\begin{array}{lc}
* & * \\
* & 0_{(\ell-1) \times(\ell-1)}
\end{array}\right], \quad j=1,2, \ldots, k .
$$

Thus, for these compressions to be linearly independent, it is necessary that $k$ does not exceed the real dimension of the space of all selfadjoint operators in $L\left(\mathbb{F}^{\ell}\right)$ having the form $\left[\begin{array}{cc}* & * \\ * & 0_{(\ell-1) \times(\ell-1)}\end{array}\right]$. This real dimension is equal to $\ell$ in the real case and to $2 \ell-1$ in the complex case. So

$$
k \leqslant \begin{cases}2 \ell-1 & \text { if } \mathbb{F}=\mathbb{C}  \tag{1.3}\\ \ell & \text { if } \mathbb{F}=\mathbb{R}\end{cases}
$$

These inequalities are satisfied in particular for some subspace $\mathcal{H}_{0}$ of dimension $\gamma_{\mathbb{F}}(k)$ for which $\widehat{X}_{1}, \ldots, \widehat{X}_{k}$ are linearly independent. Letting $\ell=\gamma_{\mathbb{F}}(k)$ in (1.3), the inequalities in (1.2) on the right follow.

To prove the inequalities on the left, let $X_{1}, \ldots, X_{k} \in L(\mathcal{H})$ be linearly independent selfadjoints, and let

$$
X_{j}=\left[\alpha_{p, q}^{(j)}\right]_{p, q \in K}, \quad j=1,2, \ldots, k
$$

be the matrix representations of the $X_{j}$ 's with respect to an orthonormal basis $\left\{e_{p}\right\}_{p \in K}$ in $\mathcal{H}$ indexed by $K$. The linear independence of the $X_{j}$ 's now implies that the (possibly infinite) linear system

$$
\sum_{j=1}^{k} x_{j} \alpha_{p, q}^{(j)}=0, \quad p, q \in K
$$

where $x_{1}, \ldots, x_{k}$ are real unknowns, has only the trivial solution. Thus, there is a subsystem of $k$ equations

$$
\begin{equation*}
\sum_{j=1}^{k} x_{j} \alpha_{p_{u}, q_{v}}^{(j)}=0, \quad u, v=1, \ldots, k \tag{1.4}
\end{equation*}
$$

that has only the trivial solution. Letting $\mathcal{H}_{0}$ be the subspace of $\mathcal{H}$ spanned by $e_{p_{1}}, \ldots, e_{p_{k}}, e_{q_{1}}, \ldots, e_{q_{k}}$ we see that $X_{1}\left[\mathcal{H}_{0}\right], \ldots, X_{k}\left[\mathcal{H}_{0}\right]$ are linearly independent, while $\operatorname{dim} \mathcal{H}_{0} \leqslant 2 k$.

If there are any repetitions in the sequence $\left\{p_{1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right\}$, then $\operatorname{dim} \mathcal{H}_{0}<2 k$ and we are done. Otherwise, let us go through another iteration of the procedure described above, with $\mathcal{H}$ replaced by $\mathcal{H}_{0}, X_{j}$ by $X_{j}\left[\mathcal{H}_{0}\right]$, and $\left\{e_{p}\right\}$ by an orthonormal basis in which the matrix of $X_{1}$ is diagonal, and moreover $\alpha_{1,1}^{(1)} \neq 0$. The respective subsystem (1.4) of $k$ equations can then be chosen in such a way that $p_{0}=q_{0}=1$, thus guaranteeing that $\operatorname{dim} \mathcal{H}_{0} \leqslant 2 k-1$.

In this paper we prove that the lower bounds for $\gamma_{\mathbb{F}}(k)$ specified in (1.2) are actually attained for $k=2,3$ (note that the case $k=1$ is trivial):

Theorem 2. Let $X_{1}, X_{2} \in L(\mathcal{H})$ be linearly independent selfadjoint operators. Then there exists a 2dimensional subspace $\mathcal{H}_{0} \subseteq \mathcal{H}$ such that $X_{1}\left[\mathcal{H}_{0}\right], X_{2}\left[\mathcal{H}_{0}\right]$ are linearly independent (over $\mathbb{R}$, or equivalently, over $\mathbb{C}$ ).

Theorem 3. Let $X_{1}, X_{2}, X_{3} \in L(\mathcal{H})$ be linearly independent selfadjoint operators. Then there exists a 2dimensional subspace $\mathcal{H}_{0} \subseteq \mathcal{H}$ in the complex case, and a 3-dimensional subspace $\mathcal{H}_{0} \subseteq \mathcal{H}$ in the real case, such that the compressions $X_{1}\left[\mathcal{H}_{0}\right], X_{2}\left[\mathcal{H}_{0}\right], X_{3}\left[\mathcal{H}_{0}\right]$ are linearly independent.

Determination of $\gamma_{\mathbb{F}}(k)$ for $k \geqslant 4$ remains an open problem.
Theorems 2 and 3 are proved in Sections 3 and 2, respectively. We return to the general problem in Section 4 where it is proved that the set of subspaces $\mathcal{H}_{0}$ of dimension $\gamma_{\mathbb{F}}(k)$ for which $X_{1}\left[\mathcal{H}_{0}\right], \ldots, X_{k}\left[\mathcal{H}_{0}\right]$ are linearly independent is open and dense. Finally, in Section 5 we apply Theorem 3 to prove a geometric property of the ratio numerical range (Theorem 6).

## 2. Proof of the main result: triples of operators

In order not to interrupt the exposition flow and for convenience of reference, we first prove an elementary auxiliary result, which is needed for the proof of the complex version only.

Lemma 4. Let $x_{1}, x_{2}, x_{3}$ be linearly independent (over $\mathbb{C}$ ) vectors in a complex inner product space. Then for every nonzero vector $y_{1}$ lying in $\mathcal{L}=\operatorname{Span}\left\{x_{1}, x_{2}, x_{3}\right\}$ there exists $y_{2} \in \mathcal{L} \backslash\{0\}, y_{2} \perp y_{1}$, such that the orthogonal projections of $x_{j}$ onto $\operatorname{Span}\left\{y_{1}, y_{2}\right\}$ are linearly independent over $\mathbb{R}$.

Proof. Choose an orthonormal basis in $\mathcal{L}$ with its first vector collinear with $y_{1}$. The first coordinate of at least one vector $x_{j}$ in this basis must be different from zero, because otherwise they would be lying in a 2 -dimensional subspace and therefore linearly dependent. Scaling and renumbering if needed, we may without loss of generality suppose that the first coordinate of $x_{1}$ is one. Further adjusting $x_{2}$ by adding real multiples of $x_{1}$ and $x_{3}$ by adding real linear combinations of $x_{1}, x_{2}$ (which has no influence on their linear dependence or independence over $\mathbb{R}$ ) we may arrange for the first coordinate of $x_{2}$ to equal $i$ or 0 while the first coordinate of $x_{3}$ is 0 .

If the first coordinate of $x_{2}$ is $i$, then we may choose $y_{2}$ as any vector in $\mathcal{L}_{1}^{\perp}$ which is not orthogonal to $x_{3}$. Otherwise (that is, if the first coordinates of $x_{2}$ and $x_{3}$ are both equal to zero) the linear independence over $\mathbb{R}$ of the projections of $x_{1}, x_{2}, x_{3}$ onto $\operatorname{Span}\left\{y_{1}, y_{2}\right\}$ will be guaranteed if $y_{2} \perp y_{1}$ and the scalar products $\left\langle x_{2}, y_{2}\right\rangle,\left\langle x_{3}, y_{2}\right\rangle$ are $\mathbb{R}$-independent, that is, they are both nonzero and their ratio is not real. In the notation

$$
x_{2}=\left[\begin{array}{c}
0 \\
\xi_{1} \\
\xi_{2}
\end{array}\right], \quad x_{3}=\left[\begin{array}{c}
0 \\
\eta_{1} \\
\eta_{2}
\end{array}\right], \quad y_{2}=\left[\begin{array}{c}
0 \\
\zeta_{1} \\
\zeta_{2}
\end{array}\right],
$$

$\mathbb{R}$-independence of $\left\langle x_{2}, y_{2}\right\rangle$, and $\left\langle x_{3}, y_{2}\right\rangle$ holds whenever $\xi_{1}+\xi_{2} z$ and $\eta_{1}+\eta_{2} z$ are $\mathbb{R}$-independent, where $z=\overline{\zeta_{2}} / \overline{\zeta_{1}}, \zeta_{1} \neq 0$, and also for $\zeta_{1}=0, \zeta_{2} \neq 0$, provided that $\xi_{2}$ and $\eta_{2}$ are $\mathbb{R}$-independent. Since $x_{2}$ and $x_{3}$ are linearly independent, there is a plenitude of vectors $y_{2}$ satisfying these requirements.

We now start the proof of Theorem 3. Let us consider separately two scenarios, depending on whether or not the operators $X_{j}$ are locally linearly independent, that is, whether or not there exist $u \in \mathcal{H}$ such that the vectors $X_{j} u, j=1,2,3$, are linearly independent (over $\mathbb{C}$ ). The interested reader may consult, e.g., [5] and references therein for some general structural results on systems of locally linearly dependent operators and the history of the subject. For our purposes, however, we merely need the definition.

Case 1. Locally linearly dependent $X_{j}, j=1,2,3$. Then

$$
\begin{equation*}
\text { for all } u \in \mathcal{H} \text {, the vectors } X_{1} u, X_{2} u, X_{3} u \text { are linearly dependent. } \tag{2.1}
\end{equation*}
$$

According to Proposition 1, it suffices to consider the case of finite (at most 5-) dimensional $\mathcal{H}$. Of course, we may also suppose that $\operatorname{dim} \mathcal{H} \geqslant 3$, because otherwise there is nothing to prove.

Denote by $m$ be the maximal rank of all operators in the real span of $X_{1}, X_{2}, X_{3}$. Passing to a different basis of this span if needed, we may without loss of generality suppose that $X_{1}$ has rank $m$. With an appropriate choice of an orthonormal basis $E=\left\{e_{j}\right\}$ in $\mathcal{H}, X_{1}$ can be represented in the block matrix form as

$$
X_{1}=\left[\begin{array}{ll}
\Lambda & 0 \\
0 & 0
\end{array}\right]
$$

where $\Lambda=\operatorname{diag}\left[\lambda_{1}, \ldots, \lambda_{m}\right]$, the diagonal matrix with real nonzero $\lambda_{1}, \ldots, \lambda_{m}$ on the main diagonal. Let

$$
X=\left[\begin{array}{ll}
A & B  \tag{2.2}\\
B^{*} & D
\end{array}\right]
$$

be the matching block matrix representation of any $X$ from the real span of $X_{j}$. Then for sufficiently small $t$ the left upper block $\Lambda+t A$ of $X_{1}+t X$ is invertible along with $\Lambda$, and from the Schur complement formula we conclude that $X_{1}+t X$ is congruent to the direct sum of $\Lambda+t A$ with

$$
Z(t)=t D-t^{2} B^{*}(\Lambda+t A)^{-1} B .
$$

Since the rank of $X_{1}+t X$ should not exceed $m$, this is only possible if $Z(t)=0$ for all $t$ in some neighborhood of zero. In particular, $D=0$.

Consequently, the $(i, j)$ entries of $X_{2}$ and $X_{3}$ are equal to zero if $i, j>m$.
Subcase 1a. There exists a pair $(i, j)$ such that the $(i, j)$ entries of the matrices of $X_{2}$ and $X_{3}$ (in the same basis $E$ that diagonalizes $X_{1}$ ) are $\mathbb{R}$-independent. Then $i \leqslant m$ or $j \leqslant m$ (since otherwise both entries are zeros) and $i \neq j$ (since otherwise they are both real). Let $\mathcal{H}_{0}=\operatorname{Span}\left\{e_{i}, e_{j}\right\}$. Then the matrix
of $X_{1}\left[\mathcal{H}_{0}\right]$ is nonzero and diagonal, while the matrices of $X_{2}\left[\mathcal{H}_{0}\right]$ and $X_{3}\left[\mathcal{H}_{0}\right]$ have $\mathbb{R}$-independent offdiagonal entries. This makes the triple $\left\{X_{j}\left[\mathcal{H}_{0}\right]: j=1,2,3\right\}$ linearly independent, which settles the subcase under consideration.

In the rest of the proof for Case 1 we may therefore suppose that the $(i, j)$ entries of $X_{2}$ and $X_{3}$ are $\mathbb{R}$-dependent for any given pair $(i, j)$.

Subcase $1 b$. The matrix of at least one of the operators $X_{2}, X_{3}$ in the basis $\left\{e_{j}\right\}$ is not diagonal.
There must be a nonzero entry in the off-diagonal part of the block $\left[\begin{array}{c}A \\ B^{*}\end{array}\right]$ in representation (2.2) of
either $X_{2}$ or $X_{3}$. Relabeling vectors $e_{1}, \ldots, e_{m}$, switching between $X_{2}$ and $X_{3}$ if needed, and denoting the matrix of $X_{3}$ by $\left(\xi_{i j}\right)$, we may without loss of generality suppose that $\xi_{i 1} \neq 0$ for some $i>1$. Then the $(i, 1)$ entry of $X_{2}$ equals $r \xi_{i 1}$ for some real $r$. Substituting $X_{2}$ by $X_{2}-r X_{3}$ we may suppose that this entry is actually equal to zero. But then, from linear dependence of $X_{j} e_{1}, j=1,2,3$ we conclude that all entries in the first column of $X_{2}$ (except maybe for its ( 1,1 )-entry) are equal to zero. Due to selfadjointness, this also holds for the elements in the first row of $X_{2}$. Finally, subtracting from $X_{2}$ an appropriate real scalar multiple of $X_{1}$, we can set its $(1,1)$ entry to equal zero as well.

Now, if the $i$ th diagonal entry of $X_{2}$ (which we will denote by $\mu_{i}$ ) is nonzero, then for $\mathcal{H}_{0}=$ $\operatorname{Span}\left\{e_{1}, e_{i}\right\}$ the matrices of $X_{j}\left[\mathcal{H}_{0}\right]$ for $j=1,2,3$ take the form

$$
\left[\begin{array}{cc}
\lambda_{1} & 0 \\
0 & \lambda_{i}
\end{array}\right],\left[\begin{array}{cc}
0 & 0 \\
0 & \mu_{i}
\end{array}\right], \quad\left[\begin{array}{cc}
* & \xi_{1 i} \\
\xi_{i 1} & *
\end{array}\right]
$$

respectively, and are therefore linearly independent.
We will show that the alternative (that is, the situation when $\mu_{i}=0$ whenever $\xi_{i 1} \neq 0$ ) leads to a contradiction, and therefore does not materialize. To this end, consider first $i \leqslant m$ and apply (2.1) to $u=e_{i}$. Since $X_{1} u=\lambda_{i} e_{i}$ and $X_{3} u$ having the first coordinate $\overline{\xi_{1 i}}$ are then linearly independent, it follows that $X_{2} u$ must be their linear combination. Making use of the fact that the first and $i$ th coordinates of $X_{2} e_{i}$ are equal to zero, we conclude that the whole $i$ th column of $X_{2}$ is equal to zero.

If $i>m$, we reach the same conclusion invoking (2.1) with $u=e_{1}+t e_{i}$ and varying $t$. Indeed, then $X_{1} u=\lambda_{1} e_{1}$ while $X_{3} u$ has a nonzero $i$ th coordinate for all except for maybe one value of $t$. Since $X_{2} u=t X_{2} e_{i}$, the $i$ th column of $X_{2}$ lies in the span of $e_{1}$ and a vector with nonzero $i$ th coordinate. This is again possible only if the combination is trivial.

We have thus shown that the $i$ th column (and therefore the $i$ th row) of $X_{2}$ is zero for any $i$ such that $\xi_{i 1} \neq 0$.

Suppose that the matrix of $X_{2}$ is nevertheless not diagonal. Then it must contain a nonzero entry say in $(s, j)$ position with $1<j \leqslant m, j \neq i$ and $s \neq 1, i, j$. (Here $i \neq 1$ is any index such that $\xi_{i 1} \neq 0$.) Let us invoke (2.1) once again, this time for $u=e_{1}+t e_{j}$. For $t \neq 0$ the vector $X_{2} u=t X_{2} e_{j}$ has a nonzero sth entry, while $X_{1} u=\lambda_{1} e_{1}+t \lambda_{j} e_{j}$. Consequently, $X_{3} u=X_{3} e_{1}+t X_{3} e_{j}$ should be a linear combination of $X_{2} e_{j}$ and $\lambda_{1} e_{1}+t \lambda_{j} e_{j}$, and therefore its $i$ th coordinate is zero. However, for $t$ small enough it is different from zero.

The contradiction obtained shows that the matrix of $X_{2}$ in the basis $\left\{e_{j}\right\}$ is $\operatorname{diag}\left[0, \mu_{2}, \ldots, \mu_{m}, 0, \ldots, 0\right]$. Note that $\mu_{j} \neq 0$ for at least one value of $j$. From (2.1) with $u=e_{1}+t e_{j}$ we conclude that for $t \neq 0, X_{3} u=X_{3} e_{1}+t X_{3} e_{j}$ must be a linear combination of the linearly independent vectors $X_{1} u=\lambda_{1} e_{1}+t \lambda_{j} e_{j}$ and $X_{2} u=t \mu_{j} e_{j}$. Let $i \neq 1$ be any index such that $\xi_{i 1} \neq 0$. Then $\mu_{i}=0$, hence $j \neq i$. Thus, the $i$ th coordinate of $X_{1} u$ and $X_{2} u=t \mu_{j} e_{j}$ is zero, while for $X_{3} u$ it differs from zero. The contradiction obtained concludes the consideration of this subcase.

The only remaining situation in Case 1 is as follows.

Subcase 1c. The matrices of all $X_{j}$ in the basis $E$ are diagonal.
Consistently with the notation and conventions made earlier, let $X_{1}$ and $X_{2}$ have the matrices $\operatorname{diag}\left[\lambda_{1}, \ldots, \lambda_{m}, 0, \ldots, 0\right]$ and $\operatorname{diag}\left[0, \mu_{2}, \ldots, \mu_{m}, 0, \ldots, 0\right]$, respectively. By changing the order of
vectors in $E$ if necessary, we may suppose that $\mu_{2} \neq 0$. Subtracting then a suitable linear combination of $X_{1}$ and $X_{2}$ from $X_{3}$, we can reduce the matrix of the latter to

$$
\operatorname{diag}\left[0,0, \nu_{3}, \ldots, v_{m}, 0, \ldots, 0\right]
$$

Now take $u=e_{1}+e_{2}+e_{k}$ with $k \neq 1$, 2 . Since $X_{1} u=\lambda_{1} e_{1}+\lambda_{2} e_{2}+\lambda_{k} e_{k}, X_{2} u=\mu_{2} e_{2}+\mu_{k} e_{k}$ and $X_{3} u=v_{k} e_{k}$, property (2.1) for this choice of $u$ implies that $v_{k}=0$. In other words, $X_{3}=0$, which is in contradiction with linear independence of the set $\left\{X_{1}, X_{2}, X_{3}\right\}$.

This concludes the consideration of Case 1 . Note that in this case we were able to come up with 2-dimensional compressions, both in real and in complex setting.

Case 2. Locally linearly independent $X_{j}, j=1,2,3$ i.e. $X_{1} u, X_{2} u, X_{3} u$ are linearly independent for some $u \in \mathcal{H}$.

We claim that then there exists a 3 -dimensional subspace $\mathcal{H}_{1}$ of $\mathcal{H}$ and $v \in \mathcal{H}_{1}$ such that the vectors $X_{j}\left[\mathcal{H}_{1}\right] v, j=1,2,3$, are linearly independent. Indeed, if $u \in \operatorname{Span}\left\{X_{1} u, X_{2} u, X_{3} u\right\}$, we may simply take

$$
\mathcal{H}_{1}=\operatorname{Span}\left\{X_{1} u, X_{2} u, X_{3} u\right\} \text { and } v=u
$$

Otherwise, let $v$ be a perturbation of $u$ so small that the set $\left\{v, X_{1} v, X_{2} v, X_{3} v\right\}=: \Gamma$ is still linearly independent, while in addition $\left\langle X_{1} v, v\right\rangle \neq 0$ (such a perturbation exists, since for any nonzero $A \in$ $L(\mathcal{H})$ the set of vectors $x$ for which $\langle A x, x\rangle \neq 0$ is dense in $\mathcal{H})$. Let $Z$ be the $4 \times 3$ matrix whose $j$ th column equals the coordinate vector of $X_{j} v$ in the orthonormal basis $\left\{v_{1}, \ldots, v_{4}\right\}$ obtained from $\Gamma$ by the Gram-Schmidt procedure. Since the rank of $Z$ equals three and its $(1,1)$ entry is nonzero, for some $j_{0}=2,3$ or 4 the deletion of the $j$ th row of $Z$ yields a nonsingular $3 \times 3$ submatrix. It remains to choose $\mathcal{H}_{1}$ as the span of $v_{j}, j \in\{1,2,3,4\} \backslash\left\{j_{0}\right\}$.

In the real setting, $\mathcal{H}_{0}=\mathcal{H}_{1}$ is the desired 3-dimensional subspace. In the complex setting, we need to go one step further. Namely, let us apply Lemma 4 to $x_{j}=X_{j}\left[\mathcal{H}_{1}\right] v, j=1,2,3$, and $y_{1}=v$. We end up with a 2 -dimensional subspace $\mathcal{H}_{0}=\operatorname{Span}\left\{y_{1}, y_{2}\right\}$ and an orthonormal basis in it such that the first columns of the matrices of $X_{j}\left[\mathcal{H}_{0}\right]$ with respect to this basis are linearly independent over $\mathbb{R}$. But then $X_{j}\left[\mathcal{H}_{0}\right]$ themselves also are linearly independent (over $\mathbb{R}$ or equivalently, over $\mathbb{C}$ ).

## 3. Proof of the main result: pairs of operators

In the complex setting, if $X_{1}, X_{2} \in L(\mathcal{H})$ are linearly independent selfadjoints, we simply adjoin a selfadjoint $X_{3} \in L(\mathcal{H})$ so that $X_{1}, X_{2}, X_{3}$ are linearly independent, and use Theorem 3 for $X_{1}, X_{2}, X_{3}$.

Consider now the real case. Let $X_{1}, X_{2} \in L(\mathcal{H})$ be linearly independent selfadjoints, where $\mathcal{H}$ is a real Hilbert space. By Proposition 1, we may (and do) assume that $\mathcal{H}$ is at most 3-dimensional. In fact, we may suppose that $\operatorname{dim} \mathcal{H}=3$, since otherwise $\operatorname{dim} \mathcal{H}=2$ and there is nothing to prove.

Step 1. Let us first show that $X_{1}$ and $X_{2}$ are locally linearly independent.
As in Section 2, with respect to a suitable orthonormal basis $\left\{e_{j}\right\}$,

$$
X_{1}=\left[\begin{array}{ll}
\Lambda & 0 \\
0 & 0
\end{array}\right], \quad X_{2}=\left[\begin{array}{cc}
A & B \\
B^{*} & 0
\end{array}\right]
$$

where $\Lambda=\operatorname{diag}\left[\lambda_{1}, \ldots, \lambda_{m}\right]$ with real nonzero $\lambda_{1}, \ldots, \lambda_{m}$. Arguing by contradiction, assume that for every $u \in \operatorname{Span}\left\{e_{1}, \ldots, e_{m}\right\} \backslash\{0\}$ there exists a real number $\mu_{u}$ such that $X_{2} u=\mu_{u} X_{1} u$. Take linearly independent $u, v \in \operatorname{Span}\left\{e_{1}, \ldots, e_{m}\right\}$ (we leave aside for the time being the case $m=1$ ), and write:

$$
X_{2} u=\mu_{u} X_{1} u, \quad X_{2} u=\mu_{v} X_{1} v, \quad X_{2}(u+v)=\mu_{u+v} X_{1}(u+v) .
$$

Subtracting the first two equalities from the third, we obtain

$$
0=\left(\mu_{u+v}-\mu_{u}\right) X_{1} u+\left(\mu_{u+v}-\mu_{v}\right) X_{1} v
$$

Linear independence of $u$ and $v$ and injectivity of $X_{1}$ on Span $\left\{e_{1}, \ldots, e_{m}\right\}$ yield

$$
\mu_{u+v}-\mu_{u}=0, \quad \mu_{u+v}-\mu_{v}=0
$$

Thus, $\mu_{u}=\mu_{v}$, and since $u, v$ is an arbitrary linearly independent pair of vectors in Span $\left\{e_{1}, \ldots, e_{m}\right\}$, it follows that in fact $\mu:=\mu_{u}$ is independent of $u \in \operatorname{Span}\left\{e_{1}, \ldots, e_{m}\right\} \backslash\{0\}$. The same conclusion obviously holds also in the case $m=1$. We now have $A=\mu \Lambda, B=0$, and $X_{2}=\mu X_{1}$, a contradiction with the linear independence of $X_{1}$ and $X_{2}$.

Step 2. Analogously to the proof in Section 2, the real setting of Case 2, it follows from the linear independence of $\left\{X_{1} u, X_{2} u\right\}$ that there exist a 2-dimensional subspace $\mathcal{H}_{1}$ of $\mathcal{H}$ and $v \in \mathcal{H}_{1}$ such that $X_{1}\left[\mathcal{H}_{1}\right] v$ and $X_{2}\left[\mathcal{H}_{1}\right] v$ are linearly independent. Thus, $X_{1}\left[\mathcal{H}_{1}\right]$ and $X_{2}\left[\mathcal{H}_{1}\right]$ are linearly independent, and the proof of Theorem 2 is complete.

## 4. Openness and denseness

Let $\mathcal{H}$ be real or complex Hilbert space. Fix a positive integer $\ell \leqslant \operatorname{dim} \mathcal{H}$. It is well known (the basic references in the context of Banach space are [6-8], see also the English translation [9] of [7]) that the set $\mathrm{Gr}_{\ell} \mathcal{H}$ of all $\ell$-dimensional subspaces of $\mathcal{H}$ is a connected complete metric space, called the Grassmannian. In the Hilbert space setting, the standard metric on the Grassmannian is given by the $\operatorname{gap} \theta(\mathcal{M}, \mathcal{N})$ between two subspaces $\mathcal{M}, \mathcal{N} \subseteq \mathcal{H}$ :

$$
\theta(\mathcal{M}, \mathcal{N})=\left\|P_{\mathcal{M}}-P_{\mathcal{N}}\right\|,
$$

where $P_{\mathcal{M}}$ stands for the orthogonal projection onto $\mathcal{M}$. Equivalently, the metric topology in $\mathrm{Gr}_{\ell} \mathcal{H}$ can be given in terms of orthonormal bases: If $f_{1}, \ldots, f_{\ell}$ is an orthonormal basis in $\mathcal{M} \in \operatorname{Gr}_{\ell} \mathcal{H}$, then a basis of open neighborhoods of $\mathcal{M}$ is given by

$$
\left\{\operatorname{Span}\left\{f_{1}^{\prime}, \ldots, f_{\ell}^{\prime}\right\}:\left\{f_{1}^{\prime}, \ldots, f_{\ell}^{\prime}\right\} \text { is orthonormal, }\left\|f_{j}^{\prime}-f_{j}\right\|<1 / m, j=1,2, \ldots, \ell\right\}
$$

for $m=1,2, \ldots$. Indeed, if $\left\{f_{j}\right\}_{j=1}^{\ell},\left\{f_{j}^{\prime}\right\}_{j=1}^{\ell}$ are orthonormal bases for $\mathcal{M}, \mathcal{M} \in \operatorname{Gr}_{\ell} \mathcal{H}$, respectively, and if $\left\|f_{j}^{\prime}-f_{j}\right\|<\alpha, j=1,2, \ldots, \ell$, for some positive $\alpha$, then the formulas

$$
P_{\mathcal{M}} x=\sum_{j=1}^{\ell}\left\langle x, f_{j}\right\rangle f_{j}, \quad P_{\mathcal{M}^{\prime} x}=\sum_{j=1}^{\ell}\left\langle x, f_{j}^{\prime}\right\rangle f_{j}^{\prime}, \quad x \in \mathcal{H}
$$

easily lead to the inequality $\left\|P_{\mathcal{M}}-P_{\mathcal{M}^{\prime}}\right\| \leqslant 2 \ell \alpha$. Conversely, if $\left\|P_{\mathcal{M}}-P_{\mathcal{M}^{\prime}}\right\| \leqslant \beta$ for $\mathcal{M}^{\prime} \in \operatorname{Gr}_{\ell} \mathcal{H}$ and for some positive $\beta$, then let $u_{j}=P_{\mathcal{M}^{\prime}} f_{j}, j=1,2, \ldots, \ell$, where $\left\{f_{1}, \ldots, f_{\ell}\right\}$ is a given orthonormal basis for $\mathcal{M}$. Using the inequalities $\left|\left\langle u_{i}, u_{j}\right\rangle-\left\langle f_{i}, f_{j}\right\rangle\right| \leqslant 2 \beta$, we see that if $\beta$ is sufficiently small, then the Gram matrix of $\left\{u_{1}, \ldots, u_{\ell}\right\}$ is invertible, hence $\left\{u_{1}, \ldots, u_{\ell}\right\}$ are linearly independent. Now the Gram-Schmidt procedure applied to $\left\{u_{1}, \ldots, u_{\ell}\right\}$ yields an orthonormal basis $\left\{f_{j}^{\prime}\right\}_{j=1}^{\ell}$ for $\mathcal{M}^{\prime}$ satisfying $\left\|f_{j}^{\prime}-f_{j}\right\| \leqslant C \beta$, where the positive constant $C$ depends on $\ell$ only, as it follows from the well-known formulas for this procedure, see e.g. [10, Theorem 9.1].

Theorem 5. Let $X_{1}, \ldots, X_{k} \in L(\mathcal{H})$ be linearly independent selfadjoint operators, where $\mathcal{H}$ is an $\mathbb{F}$-Hilbert space, $\mathbb{F}=\mathbb{R}$ or $\mathbb{F}=\mathbb{C}$. Then the set $\operatorname{LinInd}_{\mathbb{F}}\left\{X_{1}, \ldots, X_{k}\right\}$ of all $\gamma_{\mathbb{F}}(k)$-dimensional subspaces $\mathcal{H}_{0} \subseteq \mathcal{H}$ such that the compressions $X_{1}\left[\mathcal{H}_{0}\right], X_{2}\left[\mathcal{H}_{0}\right], \ldots, X_{k}\left[\mathcal{H}_{0}\right]$ are linearly independent is open and dense in $\operatorname{Gr}_{\gamma_{\mathbb{F}}(k)} \mathcal{H}$.

Proof. We prove the result for the complex case only; the proof in the real case is completely analogous. We begin with denseness. Let

$$
\mathcal{H}_{1} \in \operatorname{LinInd}_{\mathbb{C}}\left\{X_{1}, \ldots, X_{k}\right\}
$$

and let $\mathcal{H}_{0}$ be any subspace of $\mathcal{H}$ of dimension $\ell:=\gamma_{\mathbb{C}}(k)$. Assuming the intersection $\mathcal{H}_{1} \cap \mathcal{H}_{0}$ is p-dimensional, let $f_{1}, \ldots f_{\ell}$ and $g_{1}, \ldots, g_{\ell}$ be orthonormal bases in $\mathcal{H}_{1}$ and $\mathcal{H}_{0}$, respectively, such that
$f_{1}=g_{1}, \ldots, f_{p}=g_{p}$ is a basis for $\mathcal{H}_{1} \cap \mathcal{H}_{0}$. (The situation when $\mathcal{H}_{1} \cap \mathcal{H}_{0}=\{0\}$ is not excluded; in this case the proof goes through with obvious changes.) Consider

$$
f_{j}(t):= \begin{cases}f_{j} & j=1,2, \ldots, p ; \quad-\infty<t<\infty \\ t f_{j}+(1-t) g_{j} j=p+1, \ldots, \ell ; \quad-\infty<t<\infty .\end{cases}
$$

It is easy to see that $\left\{f_{1}(t), \ldots, f_{\ell}(t)\right\}$ are linearly independent for every $t \in \mathbb{R}$. Let $\left\{h_{1}(t), \ldots, h_{\ell}(t)\right\}$ be an orthonormal system obtained from $\left\{f_{1}(t), \ldots, f_{\ell}(t)\right\}$ by using the Gram-Schmidt process. Clearly,

$$
\left\{h_{1}(1), \ldots, h_{\ell}(1)\right\}=\left\{f_{1}, \ldots, f_{\ell}\right\} \text { and }\left\{h_{1}(0), \ldots, h_{\ell}(0)\right\}=\left\{g_{1}, \ldots, g_{\ell}\right\} .
$$

Moreover, $\mathcal{H}(t):=\operatorname{Span}\left\{h_{1}(t), \ldots, h_{\ell}(t)\right\}$ belongs to $\operatorname{LinInd}_{\mathbb{C}}\left\{X_{1}, \ldots, X_{k}\right\}$ if and only if the $k \times \ell^{2}$ matrix $Z(t)$ whose $j$ th row is

$$
\begin{aligned}
& {\left[\left\langle X_{j} h_{1}(t), h_{1}(t)\right\rangle, \ldots,\left\langle X_{j} h_{\ell}(t), h_{\ell}(t)\right\rangle,\right.} \\
& \quad \Re\left\langle X_{j} h_{1}(t), h_{2}(t)\right\rangle, \mathfrak{\Im}\left\langle X_{j} h_{1}(t), h_{2}(t)\right\rangle, \ldots, \mathfrak{R}\left\langle X_{j} h_{1}(t), h_{\ell}(t)\right\rangle, \mathfrak{\Im}\left\langle X_{j} h_{1}(t), h_{\ell}(t)\right\rangle, \\
& \left.\quad \ldots, \mathfrak{\Re}\left\langle X_{j} h_{\ell-1}(t), h_{\ell}(t)\right\rangle, \mathfrak{\Im}\left\langle X_{j} h_{\ell-1}(t), h_{\ell}(t)\right\rangle\right]
\end{aligned}
$$

has rank $k$. (We denote by $\Re z$ and $\Im z$ the real and imaginary parts of a complex number $z$.) Observe that the determinants of the $k \times k$ submatrices of $Z(t)$ are real analytic functions of the real variable $t$ and that at least one of those determinants is nonzero for $t=1$. Since the zeros of a real analytic function cannot accumulate within its domain of definition (unless the function is identically zero), it follows that the same determinant is nonzero for all $t \in(0, \epsilon)$, where $\epsilon>0$ is sufficiently small. Thus,

$$
\operatorname{Span}\left\{h_{1}(t), \ldots, h_{\ell}(t)\right\} \in \operatorname{LinInd}_{\mathbb{C}}\left\{X_{1}, \ldots, X_{k}\right\}, \quad t \in(0, \epsilon),
$$

and

$$
\lim _{t \rightarrow 0+} \theta\left(\mathcal{H}_{0}, \operatorname{Span}\left\{h_{1}(t), \ldots, h_{\ell}(t)\right\}\right)=0,
$$

which proves the denseness of LinInd $\mathbb{C}_{\mathbb{C}}\left\{X_{1}, \ldots, X_{k}\right\}$.
Next, we prove that LinInd $\mathbb{C}_{\mathbb{C}}\left\{X_{1}, \ldots, X_{k}\right\}$ is open.
Let $\mathcal{H}_{1} \in \operatorname{LinInd}_{\mathbb{C}}\left\{X_{1}, \ldots, X_{k}\right\}$, and let $f_{1}, \ldots f_{\ell}$ be an orthonormal basis for $\mathcal{H}_{1}$. For any orthonormal system $\left\{h_{1}, \ldots, h_{\ell}\right\}$ we denote by $Z\left(\left\{h_{1}, \ldots, h_{\ell}\right\}\right)$ the $k \times \ell^{2}$ matrix whose $j$ th row is

$$
\begin{aligned}
& {\left[\left\langle X_{j} h_{1}, h_{1}\right\rangle, \ldots,\left\langle X_{j} h_{\ell}, h_{\ell}\right\rangle, \mathfrak{\Re}\left\langle X_{j} h_{1}, h_{2}\right\rangle, \mathfrak{\Im}\left\langle X_{j} h_{1}, h_{2}\right\rangle, \ldots,\right.} \\
& \\
& \left.\quad \Re\left\langle X_{j} h_{1}, h_{\ell}\right\rangle, \Im\left\langle X_{j} h_{1}, h_{\ell}\right\rangle, \ldots, \Re\left\langle X_{j} h_{\ell-1}, h_{\ell}\right\rangle, \Im\left\{X_{j} h_{\ell-1}, h_{\ell}\right\rangle\right],
\end{aligned}
$$

for $j=1,2, \ldots, k$. Then at least one $k \times k$ minor of $Z\left(\left\{f_{1}, \ldots f_{\ell}\right\}\right)$ is nonzero. Clearly, there exists $\delta>0$ such that the determinant of the same $k \times k$ submatrix of $Z\left(\left\{h_{1}, \ldots, h_{\ell}\right\}\right)$ is nonzero for every orthonormal system $\left\{h_{1}, \ldots, h_{\ell}\right\}$ satisfying $\left\|h_{j}-f_{j}\right\|<\delta, j=1,2, \ldots, k$. From our description of a basis of open neighborhoods given at the beginning of this section, the openness of LinInd $\mathbb{C}\left\{X_{1}, \ldots, X_{k}\right\}$ follows.

## 5. Application to ratio numerical ranges

For $A, B \in L(\mathcal{H}), B \neq 0$, the ratio numerical range $W(A / B)$ was introduced in [3] as

$$
W(A / B)=\{\langle A x, x\rangle /\langle B x, x\rangle:\langle B x, x\rangle \neq 0\} .
$$

Note that $W(B / I)$ is simply the classical numerical range (or the field of values) $W(A)$ of $A$ and that in the particular case of finite dimensional $\mathcal{H}$ and $0 \notin W(B)$ the set $W(A / B)$ was studied earlier in [1]. On the other hand, if $\operatorname{dim} \mathcal{H}<\infty$ and $A, B$ are such that

$$
\begin{equation*}
\langle B x, x\rangle=\langle A x, x\rangle=0 \quad \Longrightarrow \quad x=0 \tag{5.1}
\end{equation*}
$$

then $W(B / A)$ coincides with the numerical range of the matrix pencil $\lambda B-A$ treated in [11,2].
From [11] it follows in particular that, under condition (5.1) and in the finite dimensional case, $W(A / B)$ is either contained in a line or a circle, or has a non-empty interior. We will show that condition (5.1) is redundant and the statement actually holds in any dimension.

Theorem 6. The ratio numerical range $W(A / B)$ either lies in a line, or in a circle, or has a non-empty interior.

Proof. Using the standard notation

$$
\mathfrak{R}(X)=\frac{1}{2}\left(X+X^{*}\right), \quad \Im(X)=\frac{1}{2 i}\left(X-X^{*}\right)
$$

for $X \in L(\mathcal{H})$, let us denote by $r(A, B)$ the number of linearly independent (over $\mathbb{R}$ ) operators among $\mathfrak{R}(A), \mathfrak{S}(A), \mathfrak{R}(B), \mathfrak{J}(B)$. According to [3, Theorem 3.1] $W(A / B)$ lies in a line or a circle if and only if $r(A, B) \leqslant 2$. So, we need only to verify that if $r(A, B) \geqslant 3$, then $W(A / B)$ has non-empty interior.

By Theorem 3, it is possible to find a 2-dimensional subspace $\mathcal{H}_{0}$ such that for the compressions $A_{0}, B_{0}$ of $A, B$ onto $\mathcal{H}_{0}$ we still have $r\left(A_{0}, B_{0}\right) \geqslant 3$. Since $W\left(A_{0} / B_{0}\right) \subset W(A / B)$, it suffices to consider the case $\operatorname{dim} \mathcal{H}=2$.

Recall the obvious properties $W((A-\lambda B) / B)=W(A / B)-\lambda$ and $W(A / B) \backslash\{0\}=1 / W(B / A)$ due to which the interiors of $W(A / B)$ and $W(B /(A-\lambda B)$ ) are (or are not) empty simultaneously. Considering $W(B /(A-\lambda B))$ with $\lambda=\operatorname{tr}(A) / \operatorname{tr}(B)$ instead of $W(A / B)$ if $\operatorname{tr}(B) \neq 0$, we therefore may without loss of generality suppose that $B$ is traceless.

We may suppose in addition that for some unit vector $x \in \mathcal{H},\langle A x, x\rangle=\langle B x, x\rangle=0$, since otherwise (5.1) holds. Choosing this $x$ as the first vector of an orthonormal basis of $\mathcal{H}$, we obtain the following matrix representations:

$$
A=\left[\begin{array}{cc}
0 & a_{12}  \tag{5.2}\\
a_{21} & a_{22}
\end{array}\right], \quad B=\left[\begin{array}{cc}
0 & b_{12} \\
b_{21} & 0
\end{array}\right]
$$

Note that $a_{22} \neq 0$ due to the condition imposed on $r(A, B)$.
In the definition of $W(A / B)$ only unit vectors in $\mathcal{H}$ both coordinates of which are nonzero will have an impact. Denoting by $s$ the ratio of their absolute values and by $\theta$ the difference of arguments, we think of $W(A / B)$ as the range of the function

$$
\begin{equation*}
f(s, \theta)=\frac{a_{22} s+a_{12} e^{i \theta}+a_{21} e^{-i \theta}}{b_{12} e^{i \theta}+b_{21} e^{-i \theta}} \tag{5.3}
\end{equation*}
$$

on the set $\left\{s>0, \theta \neq \arg \left(-b_{21} / b_{12}\right) \bmod \pi\right.$ if $\left.\left|b_{12}\right|=\left|b_{21}\right|\right\}$.
Suppose the interior of $W(A / B)$ is empty. Then the Jacobian of the mapping (5.3) is identically zero, that is, the ratio $\frac{\partial f}{\partial \theta} / \frac{\partial f}{\partial s}$ must be real for all admissible values of $s$ and $\theta$. A direct computation shows that this ratio is a linear function of $s$, namely

$$
\frac{\partial f}{\partial \theta}(s, \theta) / \frac{\partial f}{\partial s}(s, \theta)=m(\theta) s+n(\theta)
$$

where

$$
\begin{equation*}
m(\theta)=-i \frac{b_{12} e^{i \theta}-b_{21} e^{-i \theta}}{b_{12} e^{i \theta}+b_{21} e^{-i \theta}} \tag{5.4}
\end{equation*}
$$

and

$$
\begin{equation*}
n(\theta)=\frac{i}{a_{22}}\left(a_{12} e^{i \theta}-a_{21} e^{-i \theta}-\left(a_{12} e^{i \theta}+a_{21} e^{-i \theta}\right) \frac{b_{12} e^{i \theta}-b_{21} e^{-i \theta}}{b_{12} e^{i \theta}+b_{21} e^{-i \theta}}\right) \tag{5.5}
\end{equation*}
$$

So, both $m(\theta)$ and $n(\theta)$ must be real for all $\theta$ such that $b_{12} e^{i \theta}+b_{21} e^{-i \theta} \neq 0$.

From (5.4) we conclude then that $\left|b_{12}\right|=\left|b_{21}\right|$. Applying a diagonal unitary similarity to $A$ and $B$ given by (5.2), which of course does not change $W(A / B)$ and $r(A, B)$, we may then arrange that $b_{12}=b_{21}$. Scaling both $A$ and $B$ by their inverse (which once again will not change the ratio numerical range and the number $r(A, B)$ ) we may even suppose that $b_{12}=b_{21}=1$. Under this condition (5.5) simplifies to

$$
\begin{aligned}
& \frac{i}{a_{22}}\left(a_{12} e^{i \theta}-a_{21} e^{-i \theta}-i\left(a_{12} e^{i \theta}+a_{21} e^{-i \theta}\right) \tan \theta\right) \\
& \quad=\frac{i}{a_{22}}\left(a_{12} e^{i \theta}(1-i \tan \theta)-a_{21} e^{-i \theta}(1+i \tan \theta)\right) \\
& =\frac{i}{a_{22}}(1-i \tan \theta) e^{i \theta}\left(a_{12}-a_{21}\right),
\end{aligned}
$$

which can be real for all the admissible (and thus, by continuity, for all) values of $\theta$ only if $a_{12}=a_{21}$. But then $\mathfrak{R}(B)(=B), \mathfrak{J}(B)(=0), \mathfrak{R}(A)$ and $\Im(A)$ lie in the span of two matrices, $\left[\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right]$ and $\left[\begin{array}{ll}0 & 0 \\ 0 & 1\end{array}\right]$. This contradiction with $r(A, B) \geqslant 3$ completes the proof.

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