# SPECTRA, NORMS AND NUMERICAL RANGES OF GENERALIZED QUADRATIC OPERATORS 

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#### Abstract

A bounded linear operator acting on a Hilbert space is a generalized quadratic operator if it has an operator matrix of the form $$
\left[\begin{array}{cc} a I & c T \\ d T^{*} & b I \end{array}\right] .
$$

It reduces to a quadratic operator if $d=0$. In this paper, spectra, norms, and various kinds of numerical ranges of generalized quadratic operators are determined. Some operator inequalities are also obtained. In particular, it is shown that for a given generalized quadratic operator, the rank- $k$ numerical range, the essential numerical range, and the $q$-numerical range are elliptical disks; the $c$-numerical range is a sum of elliptical disks. The Davis-Wielandt shell is the convex hull of a family of ellipsoids unless the underlying Hilbert space has dimension 2.


## 1. Introduction

Let $\mathcal{B}(\mathcal{H})$ be the algebra of bounded linear operators acting on a Hilbert space $\mathcal{H}$. We identify $\mathcal{B}(\mathcal{H})$ with $M_{n}$ if $\mathcal{H}$ has dimension $n$. An operator $A \in \mathcal{B}(\mathcal{H})$ is a generalized quadratic operators if it has an operator matrix of the form

$$
\left[\begin{array}{cc}
a I & c T  \tag{1.1}\\
d T^{*} & b I
\end{array}\right],
$$

where $T$ is an operator from $\mathcal{H}_{2}$ to $\mathcal{H}_{1}$, and $a, b, c, d$ are complex numbers. [In the following discussion, we will not distinguish the operator and its operator matrix if there is no ambiguity.] When $d=0$, such an operator $A$ satisfies condition

$$
\begin{equation*}
(a I-A)(b I-A)=0 \tag{1.2}
\end{equation*}
$$

and is known as a quadratic operator. In fact, it is known that an operator $A$ satisfies (1.2) if and only if it has an operator matrix of the form (1.1) with $d=0$, by a suitable choice of orthonormal basis.

Motivated by theory and applied problems, there has been considerable interest in studying the norms and generalized numerical ranges (see the definition in later sections) of operators of the form (1.1) under the additional assumptions that (i) $a, b, c, d$ are nonnegative, or (ii) $d=0$; see $[3,13,22,27]$ and the references therein. In this paper, a complete description is given to the spectrum, the norm, and various types of generalized numerical ranges of an operator of the

[^0]form (1.1). In particular, the spectrum is a union of the spectrum of certain $2 \times 2$ matrices of the form
\[

A_{p}=\left[$$
\begin{array}{cc}
a & c p \\
d p & b
\end{array}
$$\right], \quad p \geq 0
\]

Also the norm of $A$ is the same as that of $A_{p}$ with $p=\|T\|$; the closure of the numerical range (and also many generalized numerical ranges) is always an elliptical disk. Since quadratic operators have been studied in [22], we always assume that $c d T \neq 0$ in the following discussion.

Our paper is organized as follows. In Section 2, we obtain a different operator matrix for an generalized quadratic operator $A$. We then use the result to give a description of $\sigma(A)$, which is the spectrum of $A \in \mathcal{B}(\mathcal{H})$. In Section 3, we determine the numerical range, the matricial range, and the norm of generalized quadratic operators. Furthermore, we obtain some operator inequalities concerning generalized quadratic operators that extend some results of Furuta [13] and Garcia [14]. We then give the description of various generalized numerical ranges of $A$ in Section 4. The results cover those in $[3,22,27]$ and the reference therein. Additional remarks and further research are discussed in Section 5.

We will use the following notations in our discussion. For $S \subseteq \mathbb{C}$, we will use $\operatorname{int}(S), \operatorname{cl}(S)$ and $\operatorname{conv}(S)$ to denote the relative interior, the closure and the convex hull of $S$, respectively. Note that in our discussion, it may happen that $S=\boldsymbol{\operatorname { c o n v }}\left\{\mu_{1}, \mu_{2}\right\}$ is a line segment in $\mathbb{C}$ so that $\operatorname{int}(S)=S \backslash\left\{\mu_{1}, \mu_{2}\right\}$.

For $A \in \mathcal{B}(\mathcal{H})$, let ker $A$ and range $A$ denote the null space and range space of $A$, respectively. Let $\mathcal{V}$ be a closed subspace of $\mathcal{H}$ and $Q$ the embedding of $\mathcal{V}$ into $\mathcal{H}$. Then $B=Q^{*} A Q$ is the compression of $A$ onto $\mathcal{V}$. More generally, $A$ has a compression $B$ if $A$ has an operator matrix $\left[\begin{array}{cc}B & * \\ * & *\end{array}\right]$ with respect to an orthonormal basis; alternatively, there is a closed subspace $\mathcal{V}$ of $\mathcal{H}$ and $X: \mathcal{V} \rightarrow \mathcal{H}$ such that $X^{*} X=I_{\mathcal{V}}$ and $X^{*} A X=B$. Note that, in this case, $X(\mathcal{V})$ is closed and $X^{*} A X$ is the compression of $A$ on $X(\mathcal{V})$.

## 2. A DIFFERENT OPERATOR MATRIX REPRESENTATION AND THE SPECTRUM

First, we obtain a different operator matrix for $A$ of the form (1.1). The special form reduces to that of quadratic operators in [22, Theorem 1.1] if $d=0$.

Theorem 2.1. Let $A \in \mathcal{B}(\mathcal{H})$ be an operator with an operator matrix

$$
\left[\begin{array}{cc}
a I & c T \\
d T^{*} & b I
\end{array}\right] \quad \text { with } \quad c d T \neq 0
$$

Then $\mathcal{H}$ has a decomposition $\mathcal{H}_{1} \oplus \mathcal{H}_{1} \oplus \mathcal{H}_{2}$ such that $A$ has an operator matrix of the form

$$
\left[\begin{array}{ll}
a I_{r} & c P  \tag{2.1}\\
d P & b I_{r}
\end{array}\right] \oplus \gamma I_{s} \quad \text { with } \quad c d P \neq 0
$$

where $\gamma \in\{a, b\}$, $\operatorname{dim} \mathcal{H}_{1}=r, \operatorname{dim} \mathcal{H}_{2}=s$, and $P: \mathcal{H}_{1} \rightarrow \mathcal{H}_{1}$ is a positive semidefinite operator, i.e., $\langle P x, x\rangle \geq 0$ for all $x \in \mathcal{H}_{1}$, with the additional condition that $\langle P x, x\rangle \neq 0$ for all nonzero
$x \in \mathcal{H}_{1}$ if $a=b$. Here $I_{s}$ may be vacuous. Furthermore, we have decompositions $\mathcal{H}_{1}=\oplus_{k=1}^{n} \mathcal{V}_{k}$ and $P=\oplus_{k=1}^{n} P_{k}$, where $n \leq \infty$ and for each $k, P_{k} \in \mathcal{B}\left(\mathcal{V}_{k}\right)$ such that $\sigma\left(P_{k}\right)$ is an interval.

Proof. Given $T$ on $\mathcal{H}_{1}$, there exist unitary operators $U$ and $V$ and a positive semidefinite operator $P$ such that $U T V$ is of one of the following form:
(a) $P$, if $\operatorname{dim} \operatorname{ker} T=\operatorname{dim} \operatorname{ker} T^{*}$.
(b) $\left[\begin{array}{c}P \\ 0\end{array}\right]$, if $\operatorname{dim} \operatorname{ker} T<\operatorname{dim} \operatorname{ker} T^{*}$.
(c) $\left[\begin{array}{ll}P & 0\end{array}\right]$, if $\operatorname{dim} \operatorname{ker} T>\operatorname{dim} \operatorname{ker} T^{*}$.

Since $\operatorname{ker} T^{*}=(\operatorname{range}(T))^{\perp}$, (a) follows from the polar decomposition. (b) and (c) follow from (a). Applying this to (1.1), we have (2.1).

Let $\lambda_{0}=\min \{\lambda: \lambda \in \sigma(P)\}$ and $\lambda_{1}=\max \{\lambda: \lambda \in \sigma(P)\}$. Then $\lambda_{0}, \lambda_{1} \in \sigma(P)$. If $\mu \notin \sigma(P)$ for some $\lambda_{0}<\mu<\lambda_{1}$, then $P=Q \oplus R$, where $\sigma(Q) \subseteq\left[\lambda_{0}, \mu\right)$ and $\sigma(R) \subseteq\left(\mu, \lambda_{1}\right]$. Therefore, we can decompose $P$ as $P=\oplus_{k=1}^{\infty} P_{k}$, where $\sigma\left(P_{k}\right)$ is an interval for each $k$.

By the above theorem, we can focus on an operator $A$ with an operator matrix of the form (2.1) with $c d P \neq 0$. In the following discussion, we will always identify the subspaces $\mathcal{H}_{1} \oplus 0 \oplus 0$ and $0 \oplus \mathcal{H}_{1} \oplus 0$ with $\mathcal{H}_{1}$. Also, the family of matrices

$$
A_{p}=\left[\begin{array}{cc}
a & c p  \tag{2.2}\\
d p & b
\end{array}\right], \quad p \geq 0
$$

will be very useful in our discussion.
Theorem 2.2. Suppose $A \in \mathcal{B}(\mathcal{H})$ has an operator matrix of the form (2.1) as in Theorem 2.1 and $A_{p}$ is defined as in (2.2). Then

$$
\sigma(A)=\cup_{p \in \sigma(P)} \sigma\left(A_{p}\right) \cup S=\cup_{p \in \sigma(P)}\left\{\frac{1}{2}\left[(a+b) \pm \sqrt{(a-b)^{2}+4 c d p^{2}}\right]\right\} \cup S,
$$

where $S=\{\gamma\}$ if $\gamma I_{s}$ is non-trivial and $S=\emptyset$ otherwise.
Proof. For simplicity, we may assume that $\gamma I_{s}$ is vacuous. If $\alpha=\frac{(a+b)}{2}$ and $\beta=\frac{(a-b)}{2}$, then

$$
A-\alpha I=\left[\begin{array}{cc}
\beta I & c P \\
d P & -\beta I
\end{array}\right] .
$$

By direct computation, we have $(A-\alpha I)^{2}=\left(\beta^{2} I+c d P^{2}\right) \oplus\left(\beta^{2} I+c d P^{2}\right)$ and

$$
\left[\begin{array}{cc}
0 & c I \\
-d I & 0
\end{array}\right](A-\alpha I)\left[\begin{array}{cc}
0 & c I \\
-d I & 0
\end{array}\right]^{-1}=-(A-\alpha I)
$$

Thus, $\mu \in \sigma(A-\alpha I)$ if and only if $\mu^{2} \in \sigma\left(\beta^{2} I+c d P^{2}\right) ; \mu \in \sigma(A-\alpha I)$ if and only if $-\mu \in$ $\sigma(A-\alpha I)$. So the result follows.

## 3. Numerical range, dilation, matricial range, and operator inequalities

Recall that the numerical range of $A \in \mathcal{B}(\mathcal{H})$ is defined by

$$
W(A)=\{\langle A x, x\rangle: x \in \mathcal{H},\langle x, x\rangle=1\}
$$

see $[15,16,17]$. The numerical range is useful in studying matrices and operators. One of the basic properties of the numerical range is that $W(A)$ is always convex; for example, see [16]. In particular, we have the following result, e.g., see [18] and [17, Theorem 1.3.6].

Elliptical Range Theorem If $A \in M_{2}$ has eigenvalues $\mu_{1}$ and $\mu_{2}$, then $W(A)$ is an elliptical disk with $\mu_{1}, \mu_{2}$ as foci and $\sqrt{\operatorname{tr}\left(A^{*} A\right)-\left|\mu_{1}\right|^{2}-\left|\mu_{2}\right|^{2}}$ as the length of minor axis. Furthermore, if $\hat{A}=A-(\operatorname{tr} A) I / 2$, then the lengths of minor and major axis of $W(A)$ are, respectively,

$$
\left\{\operatorname{tr}\left(\hat{A}^{*} \hat{A}\right)-2|\operatorname{det} \hat{A}|\right\}^{1 / 2} \quad \text { and } \quad\left\{\operatorname{tr}\left(\hat{A}^{*} \hat{A}\right)+2|\operatorname{det} \hat{A}|\right\}^{1 / 2}
$$

Using this theorem, one can deduce the convexity of the numerical range of a general operator; e.g., see [18]. It turns out that for an operator $A$ of the form (2.1), W(A) is also an elliptical disk with all the boundary points, two boundary points, or none of its boundary points as shown in the following.

Theorem 3.1. Suppose $A \in \mathcal{B}(\mathcal{H})$ has an operator matrix of the form (2.1) described in Theorem 2.1. Let $\tilde{p}=\|P\|, \tilde{A}=\left[\begin{array}{cc}a & c \tilde{p} \\ d \tilde{p} & b\end{array}\right]$ so that $\tilde{A}$ has eigenvalues $\mu_{ \pm}=\frac{1}{2}\left[(a+b) \pm \sqrt{(a-b)^{2}+4 c d \tilde{p}^{2}}\right]$ and $W(\tilde{A})$ is the elliptical disk with foci $\mu_{+}, \mu_{-}$and minor axis of length

$$
\sqrt{|a|^{2}+|b|^{2}+\tilde{p}^{2}\left(|c|^{2}+|d|^{2}\right)-\left|\mu_{+}\right|^{2}-\left|\mu_{-}\right|^{2}}
$$

If $\|P x\|=\|P\|$ for some unit vector $x \in \mathcal{H}_{1}$, then

$$
W(A)=W(\tilde{A})
$$

Otherwise, $W(A)=\operatorname{int}(W(\tilde{A})) \cup\{a, b\}$; more precisely, one of the following holds.
(1) $|c|=|d|$ and $\bar{d}(a-b)=c(\bar{a}-\bar{b})$, both $A$ and $\tilde{A}$ are normal, and

$$
W(A)=W(\tilde{A}) \backslash \sigma(\tilde{A})=\operatorname{conv}\left\{\mu_{+}, \mu_{-}\right\} \backslash\left\{\mu_{+}, \mu_{-}\right\}
$$

(2) $|c|=|d|$ and there is $\zeta \in(0, \pi)$ such that $\bar{d}(a-b)=e^{i 2 \zeta} c(\bar{a}-\bar{b}) \neq 0$, both numbers $a, b$ lie on the boundary $\partial W(A)$ of $W(A)$, and

$$
W(A)=\operatorname{int}(W(\tilde{A})) \cup\{a, b\}
$$

(3) $|c| \neq|d|$, and $W(A)=\operatorname{int}(W(\tilde{A}))$.

To prove Theorem 3.1, we need the following lemma, which will also be useful for later discussion.

Lemma 3.2. Let $A_{p}=\left[\begin{array}{cc}a & c p \\ d p & b\end{array}\right]$ for $p \geq 0$ so that $W\left(A_{p}\right)$ is the closed elliptical disk $\mathcal{E}(p)$ with foci $\mu_{ \pm}=\frac{1}{2}\left[(a+b) \pm \sqrt{(a-b)^{2}+4 c d p^{2}}\right]$ and minor axis of length

$$
\sqrt{|a|^{2}+|b|^{2}+p^{2}\left(|c|^{2}+|d|^{2}\right)-\left|\mu_{+}\right|^{2}-\left|\mu_{-}\right|^{2}}
$$

Suppose $p<q$. Then $W\left(A_{p}\right) \subseteq W\left(A_{q}\right)$.
(1) If $|c|=|d|$ and $\bar{d}(a-b)=c(\bar{a}-\bar{b})$, then $W\left(A_{p}\right)=\boldsymbol{\operatorname { c o n v }} \sigma\left(A_{p}\right)$ and $W\left(A_{q}\right)=\boldsymbol{\operatorname { c o n v }} \sigma\left(A_{q}\right)$ are line segments such that $W\left(A_{p}\right)$ is a subset of the relative interior of $W\left(A_{q}\right)$.
(2) If $|c|=|d|$ and there is $\zeta \in(0, \pi)$ such that $\bar{d}(a-b)=e^{i 2 \zeta} c(\bar{a}-\bar{b}) \neq 0$, then $\{a, b\}=$ $\partial W\left(A_{p}\right) \cap \partial W\left(A_{q}\right)$, and $W\left(A_{p}\right) \subseteq \operatorname{int}\left(W\left(A_{q}\right)\right) \cup\{a, b\}$.
(3) If $|c| \neq|d|$, then $W\left(A_{p}\right) \subseteq \operatorname{int} W\left(A_{q}\right)$.

Proof. Let $\mathcal{E}(p)=W\left(A_{p}\right)$. Then all $\mathcal{E}(p)$ has the same center $\alpha=(a+b) / 2$. Suppose $\beta=(a-b) / 2$. Denote by $\lambda_{1}(X)$ the largest eigenvalue of a self-adjoint matrix $X$. Then

$$
W\left(A_{p}\right)=\bigcap_{\xi \in[0,2 \pi)} \Pi_{\xi}\left(A_{p}\right)
$$

where

$$
\Pi_{\xi}\left(A_{p}\right)=\left\{\mu \in \mathbb{C}: e^{i \xi} \mu+e^{-i \xi} \bar{\mu} \leq \lambda_{1}\left(e^{i \xi} A_{p}+e^{-i \xi} A_{p}^{*}\right)\right\}
$$

is a half space in $\mathbb{C}$. Since

$$
\lambda_{1}\left(e^{i \xi} A_{p}+e^{-i \xi} A_{p}^{*}\right)=e^{i \xi} \alpha+e^{-i \xi} \bar{\alpha}+\sqrt{\left|e^{i \xi} \beta+e^{-i \xi} \bar{\beta}\right|^{2}+p^{2}\left|e^{i \xi} c+e^{-i \xi} d\right|^{2}}
$$

is an increasing function of $p$, we see that $\Pi_{\xi}\left(A_{p}\right) \subseteq \Pi_{\xi}\left(A_{q}\right)$ and hence $W\left(A_{p}\right) \subseteq W\left(A_{q}\right)$ if $p \leq q$.
Case 1. Suppose $a, b, c, d$ satisfy condition (1). Then $A_{p}$ is normal and $A_{p}=\alpha I_{2}+B_{p}$, where $W\left(B_{p}\right)=\boldsymbol{\operatorname { c o n v }}\left\{ \pm \sqrt{-\operatorname{det}\left(B_{p}\right)}\right\}$ is a line segment of length $2 \sqrt{|\beta|^{2}+p^{2}|c|^{2}}=2 \sqrt{|\beta|^{2}+p^{2}|d|^{2}}$. Thus, the conclusion of (1) holds.

Case 2. Suppose $a, b, c, d$ satisfy condition (2). Then $A_{p}=\alpha I_{2}+\beta B_{p}$ with

$$
e^{i \zeta} B_{p}=\left[\begin{array}{cc}
e^{i \zeta} & \delta p \\
\bar{\delta} p & -e^{i \zeta}
\end{array}\right] \quad \delta=e^{i \zeta} \frac{2 c}{a-b}=e^{-i \zeta} \frac{2 \bar{d}}{\bar{a}-\bar{b}}
$$

Using the Elliptical Range Theorem, one readily checks that $W\left(e^{i \zeta} B_{p}\right)$ is a nondegenerate elliptical disk. Since $B_{p}=\left[\begin{array}{cc}1 & \delta p e^{-i \zeta} \\ \bar{\delta} p e^{-i \zeta} & -1\end{array}\right]$,

$$
e^{i \xi} B_{p}+e^{-i \xi} B_{p}^{*}=2\left[\begin{array}{cc}
\cos \xi & \delta p \cos (\xi-\zeta) \\
\bar{\delta} p \cos (\xi-\zeta) & -\cos \xi
\end{array}\right]
$$

we have

$$
\lambda_{1}\left(e^{i \xi} B_{p}+e^{-i \xi} B_{p}^{*}\right)=2 \sqrt{\cos ^{2} \xi+|\delta|^{2} p^{2} \cos ^{2}(\xi-\zeta)} \geq \pm 2 \cos \xi= \pm\left(e^{i \xi}+e^{-i \xi}\right),
$$

where equality holds only for $\xi=\zeta \pm \pi / 2$. Therefore, 1 and -1 are on the boundary of $W\left(B_{p}\right)$ and $\lambda_{1}\left(e^{i \xi} B_{p}+e^{-i \xi} B_{p}^{*}\right)$ is a strictly increasing function for $p \geq 0$, except for $\xi=\zeta \pm \pi / 2$. From this, we get the conclusion of (2).

Case 3. Suppose $a, b, c, d$ do not satisfy the conditions in (1) or (2). Since $|c| \neq|d|$, for every $\xi \in[0,2 \pi)$,

$$
\lambda_{1}\left(e^{i \xi} A_{p}+e^{-i \xi} A_{p}^{*}\right)=e^{i \xi} \alpha+e^{-i \xi} \bar{\alpha}+\sqrt{\left|e^{i \xi} \beta+e^{-i \xi} \bar{\beta}\right|^{2}+p^{2}\left|e^{i \xi} c+e^{-i \xi} \bar{d}\right|^{2}}
$$

is a strictly increasing function for $p \geq 0$. Thus, condition (3) holds.

## Proof of Theorem 3.1

Since $W(X \oplus Y)=\boldsymbol{\operatorname { c o n v }}\{W(X) \cup W(Y)\}=W(X)$ if $W(Y) \subseteq W(X)$, we may assume that $\gamma I_{s}$ is vacuous.

Suppose $x \in \mathcal{H}=\mathcal{H}_{1} \oplus \mathcal{H}_{1}$ is a unit vector and $\mu=\langle A x, x\rangle \in W(A)$. Let $x=\left[\begin{array}{c}\cos \theta x_{1} \\ \sin \theta x_{2}\end{array}\right]$ for some unit vectors $x_{1}, x_{2} \in \mathcal{H}_{1}$. Let $\left\langle P x_{1}, x_{2}\right\rangle=p e^{i \phi}$ with $p \in[0, \tilde{p}]$ and $\phi \in[0,2 \pi)$. Then

$$
\mu=\left[\cos \theta, e^{-i \phi} \sin \theta\right] A_{p}\left[\begin{array}{c}
\cos \theta \\
e^{i \phi} \sin \theta
\end{array}\right] \in W\left(A_{p}\right) \subseteq W(\tilde{A})
$$

by Lemma 3.2.
If there is a unit vector $x \in \mathcal{H}_{1}$ such that $\|P\|=\|P x\|$, then

$$
\|P\|^{2}=\left\langle P^{2} x, x\right\rangle \leq\left\|P^{2} x\right\|\|x\| \leq\left\|P^{2}\right\|=\|P\|^{2} .
$$

Thus, $P^{2} x=\|P\|^{2} x$ and hence $P x=\|P\| x$ as $P$ is positive semi-definite. Then the operator matrix of $A$ with respect to $\mathcal{H}=\mathcal{H}_{0} \oplus \mathcal{H}_{0}^{\perp}$, where

$$
\mathcal{H}_{0}=\operatorname{span}\left\{\left[\begin{array}{l}
x \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
x
\end{array}\right]\right\}
$$

has the form $\tilde{A} \oplus \hat{A} \in \mathcal{B}(\mathcal{H})$. Thus, $W(\tilde{A}) \subseteq W(A)$, and the equality holds.
Suppose there is no unit vector $z \in \mathcal{H}_{1}$ such that $\|P\|=\|P z\|$. Then for any unit vector $x \in \mathcal{H}$, let $x=\left[\begin{array}{c}\cos \theta x_{1} \\ \sin \theta x_{2}\end{array}\right]$ for some unit vectors $x_{1}, x_{2} \in \mathcal{H}_{1}$. If $\left\langle P x_{1}, x_{2}\right\rangle=p e^{i \phi}$ with $p \in[0, \tilde{p}]$ and $\phi \in[0,2 \pi)$, then $p<\tilde{p}$. By Lemma 3.2, we see that $\mu \in \operatorname{int}(W(\tilde{A}))$ if (a) or (c) holds, and $\mu \in \operatorname{int}(W(\tilde{A})) \cup\{a, b\}$ if (b) holds.

To prove the reverse set equalities, note that there is a sequence of unit vectors $\left\{x_{m}\right\}$ in $\mathcal{H}_{1}$ such that $\left\langle P x_{m}, x_{m}\right\rangle=p_{m}$ converges to $\tilde{p}$. Then the compression of $A$ on the subspace

$$
\mathcal{V}_{m}=\operatorname{span}\left\{\left[\begin{array}{c}
x_{m} \\
0
\end{array}\right],\left[\begin{array}{c}
0 \\
x_{m}
\end{array}\right]\right\} \subseteq \mathcal{H}=\mathcal{H}_{1} \oplus \mathcal{H}_{1}
$$

has the form $A_{p_{m}}$. Since $W\left(A_{p_{m}}\right) \rightarrow W(\tilde{A})$, we see that $\operatorname{int}(W(\tilde{A})) \subseteq W(A)$. It is also clear that $\{a, b\} \subseteq W(A)$. Thus, the set equalities in (1) - (3) hold.

By Theorem 3.1 and [7, Theorem 2.1], we have the following corollary.
Corollary 3.3. Suppose $A$ and $\tilde{A}$ satisfy the condition in Theorem 3.1. Then $A$ has a dilation of the form $I \otimes \tilde{A}$.

The $n$-th matricial range $W^{n}(L)$ of $L$ is the set of $n \times n$ matrices of the form $\phi(L)$, where $\phi: \mathcal{B}(\mathcal{H}) \rightarrow M_{n}$ is a unital completely positive map.

Theorem 3.4. Suppose $A \in \mathcal{B}(\mathcal{H})$ and $\tilde{A} \in M_{2}$ satisfy the hypothesis of Theorem 3.1. Then for every $n \geq 1, W^{n}(A)$ consists of all $n \times n$ matrices $B$ with $W(B) \subseteq \overline{W(A)}=W(\tilde{A})$.

Proof. By [7, Theorem 2.1], an $n \times n$ matrix $B \in W^{n}(\tilde{A})$ if and only if $W(B) \subseteq W(\tilde{A})$. The proof of the result is similar to the proof of Theorem 3.1 in [29].

We will consider the other kinds of numerical ranges for generalized quadratic operators in Section 4. We consider some operator inequalities in the following.

Denote by $\rho(A)=\max \{|\mu|: \mu \in \sigma(A)\}$ and $w(A)=\sup \{|\mu|: \mu \in W(A)\}$ the spectral radius and numerical radius of $A \in \mathcal{B}(\mathcal{H})$. It follows readily from Theorems 2.2 and 3.1, that

$$
\rho(A)=\rho(\tilde{A}) \quad \text { and } \quad w(A)=w(\tilde{A})
$$

if $A$ and $\tilde{A}$ are defined as in Theorem 3.1. Since $A$ has a dilation of the form $I \otimes \tilde{A}$ by Corollary 3.3, we have $\|A\| \leq\|\tilde{A}\|$. As shown in the proof of Theorem 3.1, there is a sequence of two dimensional subspaces $\left\{\mathcal{V}_{m}\right\}$ such that the compression of $A$ on $\mathcal{V}_{m}$ is $A_{p_{m}}$, which converges to $\tilde{A}$. Thus, we have $\|A\|=\|\tilde{A}\|$. Suppose $\tilde{A}$ has singular values $s_{1} \geq s_{2}$. Then $\|\tilde{A}\|=s_{1}$, $\operatorname{tr}\left(\tilde{A}^{*} \tilde{A}\right)=s_{1}^{2}+s_{2}^{2}$ and $|\operatorname{det}(\tilde{A})|=s_{1} s_{2}$. Hence, for $\tilde{p}=\|P\|$,

$$
\begin{aligned}
\|\tilde{A}\|= & \frac{1}{2}\left\{\sqrt{\operatorname{tr}\left(\tilde{A}^{*} \tilde{A}\right)+2|\operatorname{det}(\tilde{A})|}+\sqrt{\operatorname{tr}\left(\tilde{A}^{*} \tilde{A}\right)-2|\operatorname{det}(\tilde{A})|}\right\} \\
= & \frac{1}{2}\left\{\sqrt{|a|^{2}+|b|^{2}+\left(|c|^{2}+|d|^{2}\right) \tilde{p}^{2}+2\left|a b-c d \tilde{p}^{2}\right|}\right. \\
& \left.\quad+\sqrt{|a|^{2}+|b|^{2}+\left(|c|^{2}+|d|^{2}\right) \tilde{p}^{2}-2\left|a b-c d \tilde{p}^{2}\right|}\right\} .
\end{aligned}
$$

By the fact that $s_{1}^{2}$ is the larger zero of $\operatorname{det}\left(\lambda I-\tilde{A}^{*} \tilde{A}\right)$ and that $\operatorname{det}\left(\tilde{A}^{*} \tilde{A}\right)=|\operatorname{det}(\tilde{A})|^{2}$, we have

$$
\begin{aligned}
\|\tilde{A}\| & =\frac{1}{\sqrt{2}}\left\{\sqrt{\left.\operatorname{tr}\left(\tilde{A}^{*} \tilde{A}\right)+\sqrt{\left[\operatorname{tr}\left(\tilde{A}^{*} \tilde{A}\right)\right]^{2}-4|\operatorname{det}(\tilde{A})|^{2}}\right\}}\right. \\
& =\frac{1}{\sqrt{2}} \sqrt{|a|^{2}+|b|^{2}+\left(|c|^{2}+|d|^{2}\right) \tilde{p}^{2}+\sqrt{\left(|a|^{2}+|b|^{2}+\left(|d|^{2}+|c|^{2}\right) \tilde{p}^{2}\right)^{2}-4\left|a b-c d \tilde{p}^{2}\right|^{2}}} \\
& =\frac{1}{\sqrt{2}} \sqrt{|a|^{2}+|b|^{2}+\left(|c|^{2}+|d|^{2}\right) \tilde{p}^{2}+\sqrt{\left(|a|^{2}-|b|^{2}+\left(|d|^{2}-|c|^{2}\right) \tilde{p}^{2}\right)^{2}+4|a \bar{c}+\bar{b} d|^{2} \tilde{p}^{2}}}
\end{aligned}
$$

We summarize the above discussion in the following corollary, which also covers the result of Furuta [13] on $w(A)$ for $A$ of the form (1.1) for $a, b, c, d \geq 0$.

Corollary 3.5. Suppose $A$ and $\tilde{A}$ satisfy the hypothesis of Theorem 3.1. Then $\rho(A)=\rho(\tilde{A})$, $w(A)=w(\tilde{A})$, and $\|A\|=\|\tilde{A}\|$. In particular, if $a, b \in \mathbb{R}$ and $c, d \in \mathbb{C}$ satisfy $c d \geq 0$, then $\mathbf{c l}(W(A))=W(\tilde{A})$ is symmetric about the real axis, and

$$
\begin{aligned}
w(A) & =w\left(\left(A+A^{*}\right) / 2\right)=w(\tilde{A})=w\left(\left(\tilde{A}+\tilde{A}^{*}\right) / 2\right)=\rho\left(\left(\tilde{A}+\tilde{A}^{*}\right) / 2\right) \\
& =\frac{1}{2}\left\{|a+b|+\sqrt{(a-b)^{2}+(|c|+|d|)^{2}\|P\|^{2}}\right\}
\end{aligned}
$$

and

$$
\|A\|=\|\tilde{A}\|=\frac{1}{2}\left\{\sqrt{(a+b)^{2}+(|b|-|c|)^{2}\|P\|^{2}}+\sqrt{(a-b)^{2}+(|b|+|c|)^{2}\|P\|^{2}}\right\} .
$$

Proof. The first assertion follows readily from Theorem 3.1. Suppose $a, b \in \mathbb{R}$ and $c, d \in \mathbb{C}$ with $c d \geq 0$. Then there is a diagonal unitary matrix $D=\operatorname{diag}(1, \mu)$ such that $D^{*} \tilde{A} D=$ $\left[\begin{array}{cc}a & |c|\|P\| \\ |d|\|P\| & b\end{array}\right]$. It is then easy to get the equalities.

Corollary 3.6. Let $A_{i}$ be self-adjoint operators on $\mathcal{H}_{i}$ with $\sigma\left(A_{i}\right) \subseteq[m, M]$ for $i=1,2$, and let $T$ be an operator from $\mathcal{H}_{2}$ to $\mathcal{H}_{1}$. Then

$$
w\left(\left[\begin{array}{cc}
A_{1} & T  \tag{3.1}\\
T^{*} & -A_{2}
\end{array}\right]\right) \leq \frac{1}{2}(M-m)+\frac{1}{2} \sqrt{(M+m)^{2}+4\|T\|^{2}} .
$$

Proof. For two self-adjoint operators $X, Y \in \mathcal{B}(\mathcal{H})$, we write $X \leq Y$ if $Y-X$ is positive semidefinite. Since $m I \leq A_{i} \leq M I$ for $i=1,2$, we have

$$
\left[\begin{array}{cc}
m I & T \\
T^{*} & -M I
\end{array}\right] \leq\left[\begin{array}{cc}
A_{1} & T \\
T^{*} & -A_{2}
\end{array}\right] \leq\left[\begin{array}{cc}
M I & T \\
T^{*} & -m I
\end{array}\right] .
$$

By Theorem 3.1,

$$
\left\|\left[\begin{array}{cc}
m I & T \\
T^{*} & -M I
\end{array}\right]\right\|=\left\|\left[\begin{array}{cc}
M I & T \\
T^{*} & -m I
\end{array}\right]\right\|=\frac{1}{2}(M-m)+\frac{1}{2} \sqrt{(M+m)^{2}+4\|T\|^{2}} .
$$

The desired inequality (3.1) holds.
Note that if $X, Y \in \mathcal{B}(\mathcal{H})$, then we have the unitary similarity relations

$$
\begin{aligned}
{\left[\begin{array}{cc}
S+i Y & 0 \\
0 & X-i Y
\end{array}\right] } & =\frac{1}{\sqrt{2}}\left[\begin{array}{cc}
I & i I \\
i I & I
\end{array}\right]\left[\begin{array}{cc}
X & -Y \\
Y & X
\end{array}\right]\left[\begin{array}{cc}
I & -i I \\
-i I & I
\end{array}\right] \frac{1}{\sqrt{2}} \\
& =\frac{1}{\sqrt{2}}\left[\begin{array}{cc}
I & I \\
-I & I
\end{array}\right]\left[\begin{array}{cc}
X & i Y \\
i Y & X
\end{array}\right]\left[\begin{array}{cc}
I & -I \\
I & I
\end{array}\right] \frac{1}{\sqrt{2}} .
\end{aligned}
$$

Thus,

$$
\max \{\|X+i Y\|,\|X-i Y\|\}=\left\|\left[\begin{array}{cc}
X & -Y \\
Y & X
\end{array}\right]\right\|=\left\|\left[\begin{array}{cc}
X & i Y \\
i Y & X
\end{array}\right]\right\| .
$$

Consequently, if $X, Y \in \mathcal{B}(\mathcal{H})$ are self-adjoint with $\sigma(X) \subseteq[m, M]$, then using Corollary 3.6, we have

$$
\begin{aligned}
\|X+i Y\| & =\|X-i Y\|=\left\|\left[\begin{array}{cc}
X & i Y \\
i Y & X
\end{array}\right]\right\|=\left\|\left[\begin{array}{cc}
X & -Y \\
Y & X
\end{array}\right]\right\|=\left\|\left[\begin{array}{cc}
X & Y \\
-Y & X
\end{array}\right]\right\| \\
& \leq \frac{1}{2}(M-m)+\frac{1}{2} \sqrt{(M+m)^{2}+4\|Y\|^{2}} .
\end{aligned}
$$

This covers a result in [14]. Of course, one can deduce the optimal norm bounds of the sum of normal operators using the results in [7], and the bounds of the sum of two general operators using the results in [8].

## 4. Generalized numerical ranges

Motivated by theoretical study and applications, there has been many generalizations of the numerical range such as the the $q$-numerical range, $k$-numerical range, the $c$-numerical range, and the essential numerical range; for example, see $[2,10,11,16,17,19,24,30]$ and their references. Recently, researchers have studied the higher rank numerical range in connection to quantum error correction; see $[4,5,6,21,23]$ and Section 2.1. Each of these generalizations encodes certain specific information of the operator that leads to interesting applications. To advance the study of these generalized numerical ranges, it is useful to have concrete descriptions of the
numerical ranges of certain operators. In most cases, it is relatively easy to solve the problem for self-adjoint or normal operators. The task is more challenging for general operators. In the following, we consider $A$ of the form (2.1). Note that some of the proofs are modifications of those in [22], but some other require different treatments. Moreover, some of the results actually are different. We will include remarks indicating these different situations.
4.1. $q$-numerical range. For $q \in[0,1]$, the $q$-numerical range of $A$ is the set

$$
W_{q}(A)=\{\langle A x, y\rangle: x, y \in \mathcal{H},\langle x, x\rangle=\langle y, y\rangle=1,\langle x, y\rangle=q\} .
$$

It is known [19, 28] that

$$
\begin{equation*}
W_{q}(A)=\left\{q\langle A x, x\rangle+\sqrt{1-q^{2}}\langle A x, y\rangle: \exists \text { orthonormal }\{x, y\} \subseteq \mathcal{H}\right\}, \tag{4.1}
\end{equation*}
$$

and also

$$
\begin{equation*}
W_{q}(A)=\left\{q \mu+\sqrt{1-q^{2}} \nu: \exists x \in \mathcal{H} \text { with }\|x\|=1, \mu=\langle A x, x\rangle,|\mu|^{2}+|\nu|^{2} \leq\|A x\|^{2}\right\} . \tag{4.2}
\end{equation*}
$$

If $q=1, W_{q}(A)=W(A)$. For $0 \leq q<1$, we have the following description of $W_{q}(A)$ for a generalized quadratic operator $A \in \mathcal{B}(\mathcal{H})$. In particular, $W_{q}(A)$ will always be an open or closed elliptical disk, which may degenerate to a line segment or a point.

Theorem 4.1. Suppose $A$ and $\tilde{A}$ satisfy the condition in Theorem 3.1. For any $q \in[0,1)$, if there is a unit vector $z \in \mathcal{H}_{1}$ such that $\|P z\|=\|P\|$, then $W_{q}(A)=W_{q}(\tilde{A})$; otherwise $W_{q}(A)=\operatorname{int}\left(W_{q}(\tilde{A})\right)$.

We need the following lemma.
Lemma 4.2. Let $A_{p}$ be defined as in (2.2). If $p<q$, then for any unit vector $x \in \mathbb{C}^{2}$ there is a unit vector $y \in \mathbb{C}^{2}$ such that $\left\langle A_{p} x, x\right\rangle=\left\langle A_{q} y, y\right\rangle$ and $\left\|A_{p} x\right\|<\left\|A_{q} y\right\|$.

Proof. Choose a unit vector $y$ orthogonal to $x$ such that $A_{p} x=\mu_{1} x+\nu_{1} y$. Let $U=[x \mid y]$. Then $U^{*} A_{p} U$ has the form

$$
\hat{A}_{p}=\left[\begin{array}{ll}
\mu_{1} & \mu_{2} \\
\nu_{1} & \nu_{2}
\end{array}\right]
$$

with $\left\langle A_{p} x, x\right\rangle=\mu_{1}$ and $\left\|A_{p} x\right\|^{2}=\left|\mu_{1}\right|^{2}+\left|\nu_{1}\right|^{2}$. Since $W\left(A_{p}\right) \subseteq W\left(A_{q}\right)$ by Lemma 3.2 and $\operatorname{tr} A_{p}=\operatorname{tr} A_{q}$, we see that $A_{q}$ is unitarily similar to a matrix of the form

$$
\hat{A}_{q}=\left[\begin{array}{ll}
\mu_{1} & \hat{\mu}_{2} \\
\hat{\nu}_{1} & \nu_{2}
\end{array}\right] .
$$

Since $X \in M_{2}$ is unitarily similar to $X^{t}$, we may assume that $\left|\hat{\nu}_{1}\right| \geq\left|\hat{\mu}_{2}\right|$. Note that

$$
\begin{gather*}
\left|\hat{\nu}_{1}\right|^{2}+\left|\hat{\mu}_{2}\right|^{2}-\left|\nu_{1}\right|^{2}-\left|\mu_{2}\right|^{2}=\operatorname{tr}\left(\hat{A}_{q}^{*} \hat{A}_{q}-\hat{A}_{p}^{*} \hat{A}_{p}\right)  \tag{4.3}\\
=\operatorname{tr}\left(A_{q}^{*} A_{q}-A_{p}^{*} A_{p}\right)=\left(|c|^{2}+|d|^{2}\right)\left(|q|^{2}-|p|^{2}\right) \geq 0
\end{gather*}
$$

and

$$
\begin{aligned}
& \left|\left|\hat{\nu}_{1} \hat{\mu}_{2}\right|-\left|\nu_{1} \mu_{2}\right|\right| \leq\left|\hat{\nu}_{1} \hat{\mu}_{2}-\nu_{1} \mu_{2}\right|=\left|\operatorname{det}\left(\hat{A}_{p}\right)-\operatorname{det}\left(\hat{A}_{q}\right)\right| \\
= & \left|\operatorname{det}\left(A_{p}\right)-\operatorname{det}\left(A_{q}\right)\right|=|c d|\left(|q|^{2}-|p|^{2}\right) .
\end{aligned}
$$

Consequently,

$$
\left(\left|\hat{\nu}_{1}\right|+\left|\hat{\mu}_{2}\right|\right)^{2}-\left(\left|\nu_{1}\right|+\left|\mu_{2}\right|\right)^{2} \geq(|c|-|d|)^{2}\left(|q|^{2}-|p|^{2}\right) \geq 0
$$

and

$$
\left(\left|\hat{\nu}_{1}\right|-\left|\hat{\mu}_{2}\right|\right)^{2}-\left(\left|\nu_{1}\right|-\left|\mu_{2}\right|\right)^{2} \geq(|c|-|d|)^{2}\left(|q|^{2}-|p|^{2}\right) \geq 0 .
$$

So we have

$$
\left|\hat{\nu}_{1}\right|+\left|\hat{\mu}_{2}\right| \geq\left|\nu_{1}\right|+\left|\mu_{2}\right| \quad \text { and } \quad\left|\hat{\nu}_{1}\right|-\left|\hat{\mu}_{2}\right| \geq\left|\left|\nu_{1}\right|-\left|\mu_{2}\right|\right| \geq\left|\nu_{1}\right|-\left|\mu_{2}\right|
$$

which implies that $\left|\hat{\nu}_{1}\right| \geq\left|\nu_{1}\right|$. From the proof, we can see that if $\left|\hat{\nu}_{1}\right|=\left|\nu_{1}\right|$, then we must have $|c|=|d|$ and $\left|\hat{\mu}_{2}\right|=\left|\mu_{2}\right|$. Then the left hand side of (4.3) is 0 , a contradiction.

Therefore, we must have $\left|\hat{\nu}_{1}\right|>\left|\nu_{1}\right|$ and the result follows.

## Proof of Theorem 4.1

By Corollary 3.3, $A$ has a dilation of the form $I \otimes \tilde{A}$. It is then easy to check that

$$
W_{q}(A) \subseteq W_{q}(I \otimes \tilde{A})=W_{q}(\tilde{A})
$$

Let $\left\{z_{m}\right\}$ be a sequence of unit vectors in $\mathcal{H}_{1}$ such that $\left\langle P z_{m}, z_{m}\right\rangle=p_{m} \rightarrow\|P\|$. Then the compression of $A$ on the subspace $\mathcal{V}_{m}=$ span $\left\{\left[\begin{array}{c}z_{m} \\ 0\end{array}\right],\left[\begin{array}{c}0 \\ z_{m}\end{array}\right]\right\}$ equals $A_{p_{m}}$ as defined in (2.2). Thus, $W_{q}\left(A_{p_{m}}\right) \subseteq W(A)$ for all $m$. Since $A_{p_{m}} \rightarrow \tilde{A}$, we see that $\operatorname{int}\left(W_{q}(\tilde{A})\right) \subseteq W_{q}(A)$. Suppose there is a unit vector $z \in \mathcal{H}_{1}$ such that $\|P z\|=\|P\|$. We may assume that $z_{m}=z$ for each $m$ so that $W_{q}(A)=W_{q}(\tilde{A})$.

Suppose there is no unit vector $z \in \mathcal{H}_{1}$ such that $\|P z\|=\|P\|$. For any unit vectors $x, y \in \mathcal{H}$ with $\langle x, y\rangle=q$, we show that $\zeta=\langle A x, y\rangle \in \operatorname{int}\left(W_{q}(\tilde{A})\right)$ in the following. To prove our claim, let $x=\left[\begin{array}{l}\alpha_{1} u_{1} \\ \alpha_{2} u_{2} \\ \alpha_{3} u_{3}\end{array}\right], y=\left[\begin{array}{l}\beta_{1} u_{1}+\gamma_{1} v_{1} \\ \beta_{2} u_{1}+\gamma_{2} v_{2} \\ \beta_{3} u_{3}+\gamma_{3} v_{3}\end{array}\right] \in \mathcal{H}_{1} \oplus \mathcal{H}_{1} \oplus \mathcal{H}_{2}$ such that $u_{1}, u_{2}, v_{1}, v_{2} \in \mathcal{H}_{1}, u_{3}, v_{3} \in \mathcal{H}_{2}$ are unit vectors, $u_{i} \perp v_{i}$ for $i=1,2,3$, and $\left(\alpha_{1}, \alpha_{2}, \alpha_{3}\right)^{t},\left(\beta_{1}, \beta_{2}, \beta_{3}\right)^{t},\left(\gamma_{1}, \gamma_{2}, \gamma_{3}\right)^{t} \in \mathbb{C}^{3}$. Then the compression of $A$ on

$$
\mathcal{V}=\operatorname{span}\left\{\left[\begin{array}{c}
u_{1} \\
0 \\
0
\end{array}\right],\left[\begin{array}{c}
0 \\
u_{2} \\
0
\end{array}\right],\left[\begin{array}{c}
0 \\
0 \\
u_{3}
\end{array}\right],\left[\begin{array}{c}
v_{1} \\
0 \\
0
\end{array}\right],\left[\begin{array}{c}
0 \\
v_{2} \\
0
\end{array}\right],\left[\begin{array}{l}
0 \\
0 \\
v_{3}
\end{array}\right]\right\}
$$

has the form

$$
B=\left[\begin{array}{ccc}
a I_{2} & c P_{0} & 0 \\
d P_{0}^{*} & b I_{2} & 0 \\
0 & 0 & \gamma I_{2}
\end{array}\right]
$$

where $P_{0} \in M_{2}$ satisfies $\left\|P_{0}\right\|<\|P\|$. Here note that the $(3,3)$ entry may be vacuous. Then $B$ is unitarily similar to

$$
\tilde{B}=\left[\begin{array}{ccc}
a I_{2} & c D_{0} & 0 \\
d D_{0} & b I_{2} & 0 \\
0 & 0 & \gamma I_{2}
\end{array}\right] \text {, where } D_{0}=\left[\begin{array}{cc}
d_{1} & 0 \\
0 & d_{2}
\end{array}\right] \text { with } d_{2} \leq d_{1}<\|P\|
$$

Since $W(\tilde{B}) \subseteq W\left(A_{d_{1}}\right), \tilde{B}$ has a dilation $I \otimes A_{d_{1}}$. Therefore, $W_{q}(\tilde{B}) \subseteq W_{q}\left(I \otimes A_{d_{1}}\right)=W_{q}\left(A_{d_{1}}\right)$ Since $\zeta \in W_{q}(\tilde{B}) \subseteq W_{q}\left(A_{d_{1}}\right), A_{d_{1}}$ is unitarily similar to a matrix of the form

$$
\hat{A}=\left[\begin{array}{ll}
\mu_{1} & \mu_{2} \\
\nu_{1} & \nu_{2}
\end{array}\right]
$$

such that $\zeta=q \mu_{1}+\sqrt{1-q^{2}} \nu_{1}$, where $\mu_{1}=\left\langle A_{d_{1}} x^{\prime}, x^{\prime}\right\rangle$ and $\left|\nu_{1}\right|^{2} \leq\left\|A_{d_{1}} x_{1}\right\|^{2}-\left|\mu_{1}\right|^{2}$ for some unit vector $x^{\prime} \in \mathbb{C}^{2}$. By Lemma 4.2, there is a unit vector $y^{\prime} \in \mathbb{C}^{2}$ such that $\left\langle\tilde{A} y^{\prime}, y^{\prime}\right\rangle=\mu_{1}$ and $\left\|\tilde{A} y^{\prime}\right\|>\left\|A_{d_{1}} x^{\prime}\right\|$. Hence, by (4.1) and (4.2),

$$
\zeta \in\left\{q \mu_{1}+\sqrt{1-q^{2}} \nu:\left|\mu_{1}\right|^{2}+|\nu|^{2}<\left\|\tilde{A} y^{\prime}\right\|^{2}\right\} \subseteq \operatorname{int}\left(W_{q}(\tilde{A})\right) .
$$

4.2. Rank- $k$ numerical ranges and essential numerical ranges. For a positive integer $k$, define the rank-k numerical range of $A \in \mathcal{B}(\mathcal{H})$ by

$$
\Lambda_{k}(A)=\{\lambda \in \mathbb{C}: P A P=\lambda P \text { for some rank- } k \text { orthogonal projection } P\}
$$

This generalized numerical range is motivated by the study of quantum error correction; see $[4,5,6]$.

To describe some basic results of $\Lambda_{k}(A)$, we need the following notation. Let $H \in \mathcal{B}(\mathcal{H})$ be a self-adjoint operator. If $\operatorname{dim} \mathcal{H}=n$, denote by $\lambda_{1}(H) \geq \cdots \geq \lambda_{n}(H)$ the eigenvalues of $H$. If $\mathcal{H}$ is infinite dimensional, define

$$
\lambda_{m}(H)=\sup \left\{\lambda_{m}\left(X^{*} H X\right): X \text { is an isometry from } \mathbb{C}^{m} \text { to } \mathcal{H}\right\} .
$$

It is known (see [26]) and not hard to verify that $\lambda_{m}(H)$ of an infinite dimensional operator $H$ can be determined as follows. Let

$$
\sigma_{e}(A)=\cap\{\sigma(A+F): F \in \mathcal{B}(\mathcal{H}) \text { has finite rank }\}
$$

be the essential spectrum of $A \in \mathcal{B}(\mathcal{H})$, and let

$$
\lambda_{\infty}(H)=\sup \sigma_{e}(H),
$$

which also equals the supremum of the set

$$
\sigma(H) \backslash\{\mu \in \mathbb{C}: H-\mu I \text { has a non-trivial finite dimensional null space }\} .
$$

Then $\mathcal{S}=\sigma(H) \cap\left(\lambda_{\infty}(H), \infty\right)$ has only isolated points, and we can arrange the elements in descending order, say, $\lambda_{1} \geq \lambda_{2} \geq \cdots$ counting multiplicities, i.e., each element repeats according to the dimension of its eigenspace. If $\mathcal{S}$ is infinite, then $\lambda_{j}(H)=\lambda_{j}$ for each positive integer $j$. If $\mathcal{S}$ has $m$ elements, then $\lambda_{j}(H)=\lambda_{j}$ for $j=1, \ldots, m$, and $\lambda_{j}(H)=\lambda_{\infty}(H)$ for $j>m$.

Let

$$
\Omega_{k}(A)=\bigcap_{\xi \in[0,2 \pi)}\left\{\mu \in \mathbb{C}: \operatorname{Re}\left(e^{i \xi} \mu\right) \leq \lambda_{k}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)\right\}
$$

where $\operatorname{Re}(B)=\left(B+B^{*}\right) / 2$ is the real part of $B$. It was shown in [20] that

$$
\operatorname{int}\left(\Omega_{k}(A)\right) \subseteq \Lambda_{k}(A) \subseteq \Omega_{k}(A)=\mathbf{c l}\left(\Lambda_{k}(A)\right)
$$

In particular, $\Lambda_{k}(A)=\Omega_{k}(A)$ if $A \in M_{n}$; see also [23].
The rank- $k$ numerical range of a generalized quadratic operator $A$ of the form (2.1) can be an empty set, a singleton, a line segment or an elliptical disk with all or part of its boundary. The following theorem gives the precise description of the set using Theorem 3.1 and Lemma 3.2.

Theorem 4.3. Suppose $A \in \mathcal{B}(\mathcal{H})$ satisfies the hypothesis of Theorem 3.1. Let $k$ be a positive integer not larger than $\operatorname{dim} \mathcal{H}$. Suppose $\mathcal{E}$ is the closed elliptical disk

$$
\mathcal{E}=W\left(\left[\begin{array}{cc}
a & c \lambda_{k}(P) \\
d \lambda_{k}(P) & b
\end{array}\right]\right)
$$

with foci $\mu_{ \pm}=\frac{1}{2}\left[(a+b) \pm \sqrt{(a-b)^{2}+4 c d \lambda_{k}(P)^{2}}\right]$ and minor axis of length

$$
\sqrt{|a|^{2}+|b|^{2}+\lambda_{k}(P)^{2}\left(|c|^{2}+|d|^{2}\right)-\left|\mu_{+}\right|^{2}-\left|\mu_{-}\right|^{2}} .
$$

(a) If $r+s<k$, then $\Lambda_{k}(A)=\emptyset$.
(b) If $r<k \leq r+s$, then $\Lambda_{k}(A)=\{\gamma\}$.
(c) Suppose $k \leq r$. If $P: \mathcal{H}_{1} \rightarrow \mathcal{H}_{1}$ has a compression $\operatorname{diag}\left(p_{1}, \ldots, p_{k}\right)$ with $p_{1} \geq \cdots \geq p_{k}=$ $\lambda_{k}(P)$, then $\Lambda_{k}(A)=\mathcal{E}$. Otherwise, $\Lambda_{k}(A)=\operatorname{int}(\mathcal{E}) \cup\{a, b\}$; more precisely, one of the following holds.
(1) If $|c|=|d|$ and $\bar{d}(a-b)=c(\bar{a}-\bar{b})$, then $\mathcal{E}=\operatorname{conv}\left\{\mu_{+}, \mu_{-}\right\}$is a line segment and $\Lambda_{k}(A)=\mathcal{E} \backslash\left\{\mu_{+}, \mu_{-}\right\}$.
(2) If $|c|=|d|$ and there is $\zeta \in(0, \pi)$ such that $\bar{d}(a-b)=e^{i 2 \zeta} c(\bar{a}-\bar{b}) \neq 0$, then $a, b \in \partial \Lambda_{k}(A)$ and $\Lambda_{k}(A)=\operatorname{int}(\mathcal{E}) \cup\{a, b\}$.
(3) If $|c| \neq|d|$, then $\Lambda_{k}(A)=\operatorname{int}(\mathcal{E})$.

Remark 4.4. In (c) of Theorem 4.3, it is not hard to show that another equivalent condition for $\Lambda_{k}(A)=\mathcal{E}$ is that

$$
P: \mathcal{H}_{1} \rightarrow \mathcal{H}_{1} \text { has a compression } \operatorname{diag}\left(\lambda_{1}(P), \ldots, \lambda_{k}(P)\right),
$$

and therefore $\Lambda_{\ell}(A)$ is an elliptical disk with foci $\mu_{ \pm}=\frac{1}{2}\left[(a+b) \pm \sqrt{(a-b)^{2}+4 c d \lambda_{\ell}(P)^{2}}\right]$ and minor axis of length

$$
\sqrt{|a|^{2}+|b|^{2}+\lambda_{\ell}(P)^{2}\left(|c|^{2}+|d|^{2}\right)-\left|\mu_{+}\right|^{2}-\left|\mu_{-}\right|^{2}}
$$

for any $\ell \in\{1, \ldots, k\}$. Also if $\lambda_{k}(P)=0$, then $\mathcal{E}$ becomes the line segment joining $a$ and $b$, i.e., $\mathcal{E}=\boldsymbol{\operatorname { c o n v }}\{a, b\}$. In this case, $\Lambda_{k}(A)$ equals $\operatorname{conv}\{a, b\}$.

The proof of Theorem 4.3 is very similar to Theorem 2.1 in [22]. First, we state Lemma 2.4 in [22], and then prove a lemma analogous to Lemma 2.5 in [22].

Lemma 4.5. Suppose $P$ is a positive semidefinite operator in $\mathcal{B}\left(\mathcal{H}_{1}\right)$ with $\operatorname{dim}\left(\mathcal{H}_{1}\right) \geq k$. For any $\varepsilon>0$, there exist $p_{1}, \ldots, p_{k} \in[0, \infty)$ with $\lambda_{j}(P)-\varepsilon<p_{j} \leq \lambda_{j}(P)$ for $j=1, \ldots, k$, such that $P$ has a compression of the form $\operatorname{diag}\left(p_{1}, \ldots, p_{k}\right)$.

Lemma 4.6. Let $A \in \mathcal{B}(\mathcal{H})$ be an operator of the form described in Theorem 3.1 with the additional assumption that $r=\infty$. Suppose $\mathcal{V}_{1}$ is a $k$-dimensional subspace of $\mathcal{H}$. Then there is a $(4 k+\ell)$-dimensional subspace $\mathcal{V}_{2}$ of $\mathcal{H}$ containing $\mathcal{V}_{1}$ with $\ell=\min \{s, k\}$ such that the compression of $A$ on $\mathcal{V}_{2}$ has the form

$$
\left[\begin{array}{cc}
a I_{2 k} & c P^{\prime} \\
d P^{\prime} & b I_{2 k}
\end{array}\right] \oplus \gamma I_{\ell},
$$

where $P^{\prime}=\operatorname{diag}\left(p_{1}, \ldots, p_{2 k}\right)$ is a compression of $P$, with $p_{1} \geq \cdots \geq p_{2 k}$ and $p_{i} \leq \lambda_{i}(P)$ for $1 \leq i \leq 2 k$.

Proof. Suppose $A$ has the form described in Theorem 3.1, with respect to the decomposition $\mathcal{H}=\mathcal{H}_{1} \oplus \mathcal{H}_{1} \oplus \mathcal{H}_{2}$ and $\operatorname{dim} \mathcal{H}_{1}=r=\infty$. Let $\mathcal{K}_{1}$ and $\mathcal{K}_{2}$ be $k$-dimensional subspaces of $\mathcal{H}_{1}$ such that $\mathcal{K}_{1} \oplus 0 \oplus 0$ and $0 \oplus \mathcal{K}_{2} \oplus 0$ contain the orthogonal projections of $\mathcal{V}_{1}$ on $\mathcal{H}_{1} \oplus 0 \oplus 0$ and $0 \oplus \mathcal{H}_{1} \oplus 0$, respectively. Also let $\mathcal{K}_{3}$ be a $\ell$-dimensional subspace of $\mathcal{H}_{2}$, with $\ell=\min \{s, k\}$, such that $0 \oplus 0 \oplus \mathcal{K}_{3}$ contains the orthogonal projection of $\mathcal{V}_{1}$ on $0 \oplus 0 \oplus \mathcal{H}_{2}$. Clearly, $\mathcal{K}_{1} \oplus \mathcal{K}_{2} \oplus \mathcal{K}_{3}$ contains $\mathcal{V}_{1}$. Take a $2 k$-dimensional subspace $\mathcal{K}$ of $\mathcal{H}_{1}$ containing $\mathcal{K}_{1}+\mathcal{K}_{2}$ and $\mathcal{V}_{2}=\mathcal{K} \oplus \mathcal{K} \oplus \mathcal{K}_{3}$. Then $\mathcal{V}_{2}$ also contains $\mathcal{V}_{1}$. Let $S: \mathcal{V}_{2} \hookrightarrow \mathcal{H}$ be the imbedding of $\mathcal{V}_{2}$ into $\mathcal{H}$. Then $S^{*} A S$ has an operator matrix of the form

$$
\left[\begin{array}{cc}
a I_{2 k} & c X^{*} P X \\
d X^{*} P X & b I_{2 k}
\end{array}\right] \oplus \gamma I_{\ell},
$$

where $X$ is the imbedding of $\mathcal{K}$ into $\mathcal{H}_{1}$. Furthermore, we can find a unitary operator $U$ such that

$$
U^{*} X^{*} P X U=\operatorname{diag}\left(p_{1}, \ldots, p_{2 k}\right) \quad \text { with } \quad p_{1} \geq \cdots \geq p_{2 k}
$$

Let $R=S\left(U \oplus U \oplus I_{\ell}\right)$. Then $R^{*} R=I_{4 k+\ell}$ and $R^{*} A R$ has the asserted form.

## Proof of Theorem 4.3.

We first consider the finite dimensional case. Let $n=\operatorname{dim} \mathcal{H}=2 r+s$. Assume that $P$ has eigenvalues $p_{1} \geq \cdots \geq p_{r} \geq 0$. By Theorem 2.1, $A$ is unitarily similar to

$$
B_{1} \oplus B_{2} \oplus \cdots \oplus B_{r} \oplus \gamma I_{s},
$$

where $B_{j}=A_{p_{j}}$, where $A_{p}$ is defined as in (2.2).
We note that if $a=b$, then our assumption ensures that $p_{j} \neq 0$. Therefore, $B_{j}$ is never a scalar matrix and $\Omega_{2}\left(B_{j}\right)=\emptyset$.

By the argument in the preceding paragraph,

$$
\begin{aligned}
\lambda_{1}\left(\operatorname{Re}\left(e^{i \xi} B_{1}\right)\right) & \geq \cdots \geq \lambda_{1}\left(\operatorname{Re}\left(e^{i \xi} B_{r}\right)\right) \\
& \geq \operatorname{Re}\left(e^{i \xi} \gamma\right) \geq \lambda_{2}\left(\operatorname{Re}\left(e^{i \xi} B_{r}\right)\right) \geq \cdots \geq \lambda_{2}\left(\operatorname{Re}\left(e^{i \xi} B_{1}\right)\right)
\end{aligned}
$$

Then $\lambda_{k}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)$ equals

$$
\begin{cases}\lambda_{1}\left(\operatorname{Re}\left(e^{i \xi} B_{k}\right)\right) & \text { if } k \leq r, \\ \operatorname{Re}\left(e^{i \xi} \gamma\right) & \text { if } r<k \leq r+s \\ \lambda_{2}\left(\operatorname{Re}\left(e^{i \xi} B_{n-k+1}\right)\right) & \text { if } r+s<k \leq n\end{cases}
$$

Recall that $\mu \in \Omega_{k}(A)$ if and only if $\operatorname{Re}\left(e^{i \xi} \mu\right) \leq \lambda_{k}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)$ for all $\xi \in[0,2 \pi)$. We have

$$
\Lambda_{k}(A)=\Omega_{k}(A)= \begin{cases}\Omega_{1}\left(B_{k}\right) & \text { if } k \leq r \\ \{\gamma\} & \text { if } r<k \leq r+s \\ \Omega_{2}\left(B_{n-k+1}\right)=\emptyset & \text { if } r+s<k \leq n\end{cases}
$$

Then the assertion holds when $k>r$. If $k \leq r$, then

$$
\Lambda_{k}(A)=\Omega_{k}(A)=\Omega_{1}\left(B_{k}\right)=\Lambda_{1}\left(B_{k}\right)=\mathcal{E}\left(p_{k}\right)
$$

Thus, the result holds in the finite dimensional case.
Next, suppose $\mathcal{H}$ is an infinite dimensional Hilbert space. If $r<k$, then $\Omega_{k}(A)=\{\gamma\}$ and hence $\Lambda_{k}(A)=\{\gamma\}$.

If $r \geq k$ is finite or $\lambda_{k}(P)=0$, then $P: \mathcal{H}_{1} \rightarrow \mathcal{H}_{1}$ has a compression $\operatorname{diag}\left(\lambda_{1}(P), \ldots, \lambda_{k}(P)\right)$.
Let

$$
\tilde{A}=B_{1} \oplus \cdots \oplus B_{k} \in M_{2 k}
$$

with $B_{j}=\left[\begin{array}{cc}a & c \lambda_{j}(P) \\ d \lambda_{j}(P) & b\end{array}\right]$ for $j=1, \ldots, k$. Notice that $\tilde{A}$ is a compression of $A$ and

$$
\lambda_{k}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)=\lambda_{k}\left(\operatorname{Re}\left(e^{i \xi} \tilde{A}\right)\right) \quad \text { for all } \xi \in[0,2 \pi)
$$

Hence,

$$
\Lambda_{k}(\tilde{A}) \subseteq \Lambda_{k}(A) \subseteq \Omega_{k}(A)=\Omega_{k}(\tilde{A})=\Lambda_{k}(\tilde{A})
$$

Thus, $\Lambda_{k}(A)=\Lambda_{k}(\tilde{A})$ so that the result follows from the finite dimensional case.
Suppose $r$ is infinite and $\lambda_{k}(P)>0$. We prove that (c) holds with $\mathcal{E}=\mathcal{E}\left(\lambda_{k}(P)\right)$. Let $\mu$ be an interior point of $\mathcal{E}$. Then there exists $\varepsilon>0$ such that $\mu \in \mathcal{E}\left(\lambda_{k}(P)-\varepsilon\right)$. By Lemma 4.5, there exist a $k$-dimensional subspace $\mathcal{V}$ of $\mathcal{H}$ and $X: \mathcal{V} \rightarrow \mathcal{H}_{1}$ satisfying $X^{*} X=I_{k}$ and

$$
\lambda_{k}\left(X^{*} P X\right)>\lambda_{k}(P)-\varepsilon
$$

Let $Z=\left[\begin{array}{cc}X & 0 \\ 0 & X\end{array}\right] \oplus I_{s}$. Then we have $Z^{*} A Z=\left[\begin{array}{cc}a I_{k} & c X^{*} P X \\ d X^{*} P X & b I_{k}\end{array}\right] \oplus \gamma I_{s}$ and

$$
\mu \in \mathcal{E}\left(\lambda_{k}(P)-\varepsilon\right) \subseteq \Lambda_{k}\left(Z^{*} A Z\right) \subseteq \Lambda_{k}(A)
$$

Conversely, suppose $\mu \in \Lambda_{k}(A)$. Then there exist a $k$-dimensional subspace $\mathcal{V}_{1}$ of $\mathcal{H}$ and $X: \mathcal{V}_{1} \rightarrow \mathcal{H}$ such that $X^{*} X=I_{\mathcal{V}_{1}}$ and $X^{*} A X=\mu I_{\mathcal{V}_{1}}$. By Lemma 4.6, there is a $(4 k+\ell)$ dimensional subspace $\mathcal{V}_{2}$ containing $\mathcal{V}_{1}$ such that the compression of $A$ on $\mathcal{V}_{2}$ has an operator matrix

$$
A^{\prime}=\left[\begin{array}{cc}
a I_{2 k} & c P^{\prime} \\
d P^{\prime} & b I_{2 k}
\end{array}\right] \oplus \gamma I_{\ell} \in M_{4 k+\ell}
$$

where $P^{\prime}=\operatorname{diag}\left(p_{1}, \ldots, p_{2 k}\right)$ is a $2 k$-dimensional compression of $P$, with $p_{1} \geq \cdots \geq p_{2 k}$ and $p_{i} \leq \lambda_{i}(P)$ for $1 \leq i \leq 2 k$. By the result in the finite dimensional case, we have

$$
\mu \in \Lambda_{k}\left(A^{\prime}\right)=\mathcal{E}\left(\lambda_{k}\left(P^{\prime}\right)\right) \subseteq \mathcal{E}\left(\lambda_{k}(P)\right) .
$$

So, we have shown that

$$
\operatorname{int}\left(\mathcal{E}\left(\lambda_{k}(P)\right)\right) \subseteq \Lambda_{k}(A) \subseteq \mathcal{E}\left(\lambda_{k}(P)\right)
$$

Also, if $P$ has a $k$-dimensional compression $\operatorname{diag}\left(p_{1}, \ldots, p_{k}\right)$ with $p_{k}=\lambda_{k}(P)$ and $\Lambda_{k}(A)=$ $\mathcal{E}\left(\lambda_{k}(P)\right)$. Otherwise, $\Lambda_{k}(A)$ can only contain points in the relative interior of $\mathcal{E}\left(\lambda_{k}(P)\right)$ unless (2) holds so that $a, b \in \Lambda_{k}(A) \cap \partial \mathcal{E}\left(\lambda_{k}(P)\right)$. The proof is complete.

For an infinite dimensional operator $A$, one can extend the definition of rank- $k$ numerical range to $\Lambda_{\infty}(A)$ defined as the set of scalars $\lambda \in \mathbb{C}$ such that $P A P=\lambda P$ for an infinite rank orthogonal projection $P$ on $\mathcal{H}$, see [20, 25]. Evidently, $\Lambda_{\infty}(A)$ consists of those $\lambda \in \mathbb{C}$ for which there exists an infinite orthonormal set $\left\{x_{i} \in \mathcal{H}: i \geq 1\right\}$ such that $\left\langle A x_{i}, x_{j}\right\rangle=\delta_{i j} \lambda$ for all $i, j \geq 1$. It is shown in [20] that

$$
\Lambda_{\infty}(A)=\bigcap_{k \geq 1} \Lambda_{k}(A)=\bigcap\{W(A+F): F \in \mathcal{B}(\mathcal{H}) \text { has a finite rank }\}
$$

Recall that $\lambda_{\infty}(H)$ is the supremum of the set

$$
\sigma(H) \backslash\{\mu \in \mathbb{C}: H-\mu I \text { has a non-trivial finite dimensional null space }\} .
$$

One can extend the definition of $\Omega_{k}(A)$ to

$$
\Omega_{\infty}(A)=\bigcap_{k \geq 1} \Omega_{k}(A)
$$

By Theorem 5.1 in [20] (see also [1, Theorem 4]),

$$
\Omega_{\infty}(A)=\bigcap\{\mathbf{c l}(W(A+F)): F \in \mathcal{B}(\mathcal{H}) \text { has a finite rank }\}
$$

is the essential numerical range $W_{e}(A)$ of $A ; \Omega_{\infty}(A)=\mathbf{c l}\left(\Lambda_{\infty}(A)\right)$ if and only if $\Lambda_{\infty}(A)$ is non-empty.

By Theorem 4.3, we have the following corollary, which gives a complete description of $\Lambda_{\infty}(A)$ and the essential numerical range of a quadratic operator $A$. It turns out that each of them can be a singleton, a line segment or an elliptical disk.

Corollary 4.7. Suppose $A \in \mathcal{B}(\mathcal{H})$ is an infinite dimensional generalized quadratic operator with operator matrix in the form in Theorem 3.1. Suppose

$$
\mathcal{E}=W\left(\left[\begin{array}{cc}
a & c \lambda_{\infty}(P) \\
d \lambda_{\infty}(P) & b
\end{array}\right]\right)
$$

with foci $\mu_{ \pm}=\frac{1}{2}\left[(a+b) \pm \sqrt{(a-b)^{2}+4 c d \lambda_{\infty}(P)^{2}}\right]$ and minor axis of length

$$
\sqrt{|a|^{2}+|b|^{2}+\lambda_{\infty}(P)^{2}\left(|c|^{2}+|d|^{2}\right)-\left|\mu_{+}\right|^{2}-\left|\mu_{-}\right|^{2}} .
$$

(a) If $r<\infty$, then $\Lambda_{\infty}(A)=\{\gamma\}$.
(b) Suppose $r=\infty$. If $\sigma(P) \cap\left(\lambda_{\infty}(P), \infty\right)$ is infinite, equivalently, $P-\lambda_{\infty}(P) I$ has an infinite dimensional null space, then $\Lambda_{\infty}(A)=\mathcal{E}$. Otherwise, $\Lambda_{\infty}(A)=\operatorname{int}(\mathcal{E}) \cup\{a, b\}$; more precisely, one of the following holds.
(1) If $|c|=|d|$ and $\bar{d}(a-b)=c(\bar{a}-\bar{b})$, then $\mathcal{E}=\operatorname{conv}\left\{\mu_{+}, \mu_{-}\right\}$is a line segment and $\Lambda_{\infty}(A)=\mathcal{E} \backslash\left\{\mu_{+}, \mu_{-}\right\}$.
(2) If $|c|=|d|$ and there is $\zeta \in(0, \pi)$ such that $\bar{d}(a-b)=e^{i 2 \zeta} c(\bar{a}-\bar{b}) \neq 0$, then $a, b \in \partial \Lambda_{k}(A)$ and $\Lambda_{\infty}(A)=\operatorname{int}(\mathcal{E}) \cup\{a, b\}$.
(3) If $|c| \neq|d|$, then $\Lambda_{\infty}(A)=\operatorname{int}(\mathcal{E})$.

Consequently, $W_{e}(A)=\Omega_{\infty}(A)=\mathbf{c l}\left(\Lambda_{\infty}(A)\right)$ is a singleton, a line segment or a closed elliptical disk.
4.3. $c$-numerical ranges. For $c=\left(c_{1}, \ldots, c_{k}\right)$ with $c_{1} \geq \cdots \geq c_{k}$ and $k \leq \operatorname{dim} \mathcal{H}$, the $c$ numerical range of $A$ is

$$
W_{c}(A)=\left\{\sum_{j=1}^{k} c_{j}\left\langle A x_{j}, x_{j}\right\rangle:\left\{x_{1}, \ldots, x_{k}\right\} \subseteq \mathcal{H} \text { is an orthonormal set }\right\} .
$$

If $\left(c_{1}, \ldots, c_{k}\right)=(1, \ldots, 1)$, then $W_{c}(A)$ reduces to the $k$-numerical range; see [16].
Suppose $A \in M_{2}$ has eigenvalues $\lambda_{1}$ and $\lambda_{2}$, and $c=\left(c_{1}, c_{2}\right)$. Then

$$
W_{c}(A)=\left(c_{1}-c_{2}\right) W(A)+c_{2} \operatorname{tr} A=W\left(\left(c_{1}-c_{2}\right) A+\left(c_{2} \operatorname{tr} A\right) I_{2}\right)
$$

is the elliptical disk with foci $c_{1} \lambda_{1}+c_{2} \lambda_{2}$ and $c_{2} \lambda_{1}+c_{1} \lambda_{2}$, and the lengths of minor and major axis of $W_{c}(A)$ are, respectively,

$$
\left\{\operatorname{tr}\left(\hat{A}^{*} \hat{A}\right)-2|\operatorname{det} \hat{A}|\right\}^{1 / 2} \quad \text { and } \quad\left\{\operatorname{tr}\left(\hat{A}^{*} \hat{A}\right)+2|\operatorname{det} \hat{A}|\right\}^{1 / 2},
$$

where $\hat{A}=\frac{\left(c_{1}-c_{2}\right)}{2}\left(A-2(\operatorname{tr} A) I_{2}\right)$.
For a self-adjoint operator $H \in B(\mathcal{H})$, we have

$$
\mathbf{c l}\left(W_{c}(H)\right)=\left[m_{c}(H), M_{c}(H)\right],
$$

where

$$
m_{c}(H)=\inf \left\{-\sum_{j=1}^{\ell} c_{j} \lambda_{j}(-H)+\sum_{j=1}^{k-\ell} c_{k-j+1} \lambda_{j}(H): 0 \leq \ell \leq k\right\}
$$

and

$$
M_{c}(H)=\sup \left\{\sum_{j=1}^{\ell} c_{j} \lambda_{j}(H)-\sum_{j=1}^{k-\ell} c_{k-j+1} \lambda_{j}(-H): 0 \leq \ell \leq k\right\} .
$$

For a general operator $A \in \mathcal{B}(\mathcal{H})$, we have

$$
\begin{equation*}
\mathbf{c l}\left(W_{c}(A)\right)=\bigcap_{t \in[0,2 \pi)}\left\{\mu \in \mathbb{C}: \operatorname{Re}\left(e^{i t} \mu\right) \leq M_{c}\left(\operatorname{Re}\left(e^{i t} A\right)\right)\right\} \tag{4.4}
\end{equation*}
$$

For a generalized quadratic operator $A \in \mathcal{B}(\mathcal{H})$, it is easy to determine $\lambda_{m}\left(\operatorname{Re}\left(e^{i t} A\right)\right)$. Thus, it is not hard to determine $W_{c}(A)$ using (4.4). It turns out that $\mathbf{c l}\left(W_{c}(A)\right)$ can always be expressed as the sum of a finite number of elliptical disks, namely,

$$
\mathbf{c l}\left(W_{c}(A)\right)=W\left(C_{1}\right)+\cdots+W\left(C_{t}\right)+d
$$

for some constant $d \in \mathbb{C}$ and $C_{1}, \ldots, C_{t} \in M_{2}$ with $t \leq k$.
To simplify the statement of our results, we will impose the following assumption on the vector $c=\left(c_{1}, \ldots, c_{k}\right)$ :

$$
\begin{align*}
& c_{1} \geq \cdots \geq c_{k} \quad \text { with } \quad c_{m+1}=0, \quad \text { where } \\
& \operatorname{dim} \mathcal{H}=\infty>k=2 m \quad \text { or } \quad \operatorname{dim} \mathcal{H}=k \in\{2 m, 2 m+1\} . \tag{4.5}
\end{align*}
$$

Note that it is easy to reduce the general case to the study of the special vector $c$ with assumption (4.5). In the infinite dimensional case, this can be achieved by adding zeros to the vector $c=\left(c_{1}, \ldots, c_{k}\right)$. In the finite dimensional case, we can first assume that $k=\operatorname{dim} \mathcal{H}$ by adding zeros to the vector $c$, and then replace $c$ with $\hat{c}=c-c_{m+1}(1, \ldots, 1)$. One can then use the fact that $W_{c}(A)=W_{\hat{c}}(A)+c_{m+1} \operatorname{tr} A$ to determine the shape of $W_{c}(A)$. Note also that the advantage of this assumption on $c$ is that the supremum in the definition of $M_{c}(H)$ is always attained at $\ell=m$.

For notational convenience, we assume that the generalized quadratic operator $A$ in the form $(2.2)$ has $(1,2)$ block equal to $P$ instead of $c P$ in the following theorem.

Theorem 4.8. Let $A \in \mathcal{B}(\mathcal{H})$ be a generalized quadratic operator with an operator matrix in the form

$$
\left[\begin{array}{cc}
a I_{r} & P  \tag{4.6}\\
d P & b I_{r}
\end{array}\right] \oplus \gamma I_{s} \quad \text { with } \quad d P \neq 0
$$

Suppose $c=\left(c_{1}, \ldots, c_{k}\right)$ satisfies (4.5) and $t=\min \{m, r\}$. Let

$$
\mathcal{E}=W\left(C_{1}\right)+\cdots+W\left(C_{t}\right)+\gamma \sum_{j=t+1}^{k-t} c_{j},
$$

where

$$
C_{j}=\left(c_{j}-c_{k-j+1}\right)\left[\begin{array}{cc}
a & \lambda_{j}(P) \\
d \lambda_{j}(P) & b
\end{array}\right]+c_{k-j+1}(a+b) I_{2}, \quad j=1, \ldots, t .
$$

If $P: \mathcal{H}_{1} \rightarrow \mathcal{H}_{1}$ has a compression $\operatorname{diag}\left(\lambda_{1}(P), \ldots, \lambda_{t}(P)\right)$, then $W_{c}(A)=\mathcal{E}$. Otherwise,

$$
W_{c}(A)=\operatorname{int}(\mathcal{E}) \cup\left\{\mathbf{q}_{1}, \mathbf{q}_{2}\right\}
$$

where

$$
\mathbf{q}_{1}=\sum_{j=1}^{t}\left(a c_{j}-b c_{k-j+1}\right)+\gamma \sum_{j=t+1}^{k-\ell} c_{j}, \quad \mathbf{q}_{2}=\sum_{j=1}^{t}\left(b c_{j}-a c_{k-j+1}\right)+\gamma \sum_{j=t+1}^{k-\ell} c_{j} ;
$$

more precisely, one of the following holds.
(1) If $|d|=1$ and $\bar{d}(a-b)=(\bar{a}-\bar{b})$, then $\mathcal{E}$ is a line segment and $W_{c}(A)$ is the relative interior of $\mathcal{E}$.
(2) If $|d|=1$ and there is $\zeta \in(0, \pi)$ such that $\bar{d}(a-b)=e^{i 2 \zeta}(\bar{a}-\bar{b}) \neq 0$, then $W_{c}(A)$ consists of the interior of the non-degenerate elliptical disk $\mathcal{E}$ and its two boundary points $\mathbf{q}_{1}$ and $\mathbf{q}_{2}$.
(3) If $|d| \neq 1$, then $\Lambda_{k}(A)=\operatorname{int}(\mathcal{E})$.

Proof. Suppose $\operatorname{dim} \mathcal{H}=n$ is finite. So we have $k=n$ and $r \leq m$. Notice that $A$ is unitarily similar to

$$
B_{1} \oplus \cdots \oplus B_{r} \oplus \gamma I_{s}
$$

where $B_{j}=\left[\begin{array}{cc}a & \lambda_{j}(P) \\ d \lambda_{j}(P) & b\end{array}\right]$ for $j=1, \ldots, r$. By the argument in the proof of Theorem 4.3, we have

$$
\lambda_{j}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)= \begin{cases}\lambda_{1}\left(\operatorname{Re}\left(e^{i \xi} B_{j}\right)\right) & \text { if } j \leq r \\ \operatorname{Re}\left(e^{i \xi} \gamma\right) & \text { if } r<j \leq r+s \\ \lambda_{2}\left(\operatorname{Re}\left(e^{i \xi} B_{n-j+1}\right)\right) & \text { if } r+s<j \leq n\end{cases}
$$

Under assumption (4.5) and $k=n$, we have

$$
M_{c}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)=\sum_{j=1}^{n} c_{j} \lambda_{j}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)
$$

On the other hand,

$$
\operatorname{Re}\left(e^{i \xi} \gamma\right) \sum_{j=r+1}^{n-r} c_{j}=\sum_{j=r+1}^{n-r} c_{j} \lambda_{j}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)
$$

and

$$
\begin{aligned}
& \sum_{j=1}^{r} M_{\left(c_{j}, c_{n-j+1}\right)}\left(\operatorname{Re}\left(e^{i \xi} B_{j}\right)\right) \\
= & \sum_{j=1}^{r}\left[c_{j} \lambda_{1}\left(\operatorname{Re}\left(e^{i \xi} B_{j}\right)\right)+c_{n-j+1} \lambda_{2}\left(\operatorname{Re}\left(e^{i \xi} B_{j}\right)\right)\right] \\
= & \sum_{j=1}^{r} c_{j} \lambda_{j}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)+\sum_{j=n-r+1}^{n} c_{j} \lambda_{j}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right) .
\end{aligned}
$$

Thus, $M_{c}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)$ equals

$$
\sum_{j=1}^{r} M_{\left(c_{j}, c_{n-j+1}\right)}\left(\operatorname{Re}\left(e^{i \xi} B_{j}\right)\right)+\operatorname{Re}\left(e^{i \xi} \gamma\right) \sum_{j=r+1}^{n-r} c_{j}
$$

By (4.4) and the above equation, the two compact convex sets

$$
W_{c}(A) \quad \text { and } \quad W_{\left(c_{1}, c_{n}\right)}\left(B_{1}\right)+\cdots+W_{\left(c_{r}, c_{n-r+1}\right)}\left(B_{r}\right)+\gamma \sum_{j=r+1}^{n-r} c_{j}
$$

always share the same support line in each direction. Thus, the two sets are the same. Since $W_{\left(c_{j}, c_{n-r+j}\right)}\left(B_{j}\right)=W\left(C_{j}\right)$ for $j=1, \ldots, r$, it follows that

$$
W_{c}(A)=W\left(C_{1}\right)+\cdots+W\left(C_{r}\right)+\gamma \sum_{j=r+1}^{n-r} c_{j} .
$$

Next, suppose $\operatorname{dim} \mathcal{H}$ is infinite. Suppose $r$ is finite or $\lambda_{m}(P)=0$. Let $t=\min \{m, r\}$. Then $P: \mathcal{H}_{1} \rightarrow \mathcal{H}_{1}$ has a compression $\operatorname{diag}\left(\lambda_{1}(P), \ldots, \lambda_{t}(P)\right)$. Take

$$
\tilde{A}=B_{1} \oplus \cdots \oplus B_{t} \oplus \gamma I_{k-2 t} \in M_{k}
$$

with $B_{j}=\left[\begin{array}{cc}a & \lambda_{j}(P) \\ d \lambda_{j}(P) & b\end{array}\right]$ for $j=1, \ldots, t$. Then we have $\lambda_{j}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)=\lambda_{j}\left(\operatorname{Re}\left(e^{i \xi} \tilde{A}\right)\right)$ for each $\xi \in[0,2 \pi)$ and $j=1, \ldots, m$. Thus, $M_{c}\left(\operatorname{Re}\left(e^{i \xi} A\right)\right)=M_{c}\left(\operatorname{Re}\left(e^{i \xi} \tilde{A}\right)\right)$ for all $\xi \in[0,2 \pi)$ and so $W_{c}(A)=W_{c}(\tilde{A})$. The result follows from the finite dimensional case.

Suppose $r$ is infinite and $\lambda_{m}(P)>0$. For $\mu_{1} \geq \cdots \geq \mu_{m}>0$, let

$$
\mathcal{E}\left(\mu_{1}, \ldots, \mu_{m}\right)=W_{\left(c_{1}, c_{k}\right)}\left(\left[\begin{array}{cc}
a & \mu_{1} \\
d \mu_{1} & b
\end{array}\right]\right)+\cdots+W_{\left(c_{m}, c_{k-m+1}\right)}\left(\left[\begin{array}{cc}
a & \mu_{m} \\
d \mu_{m} & b
\end{array}\right]\right)
$$

Notice that $\mathcal{E}\left(\lambda_{1}(P), \cdots \lambda_{m}(P)\right)=W\left(C_{1}\right)+\cdots+W\left(C_{m}\right)$.
By Lemma 4.5, there exist an $m$-dimensional subspace $\mathcal{V}$ of $\mathcal{H}$ and $X: \mathcal{V} \rightarrow \mathcal{H}_{1}$ satisfying $X^{*} X=I_{m}$ and $X^{*} P X=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{m}\right)$ with $\lambda_{j}(P)-\varepsilon<\lambda_{j} \leq \lambda_{j}(P)$ for $j=1, \ldots, m$. Let $Z=\left[\begin{array}{cc}X & 0 \\ 0 & X\end{array}\right] \oplus I_{s}$. Then $Z^{*} A Z$ is unitary similar to

$$
\left[\begin{array}{cc}
a & \lambda_{1} \\
d \lambda_{1} & b
\end{array}\right] \oplus \cdots \oplus\left[\begin{array}{cc}
a & \lambda_{m} \\
d \lambda_{m} & b
\end{array}\right] \oplus \gamma I_{s} .
$$

Note that $W_{c}(B) \subseteq W_{c}(A)$ if $B$ is a compression of $A$. Applying the result for finite $r=m$, we have

$$
\mathcal{E}\left(\lambda_{1}, \ldots, \lambda_{m}\right)=W_{c}\left(Z^{*} A Z\right) \subseteq W_{c}(A)
$$

As $\lambda_{j} \rightarrow \lambda_{j}(P)$ and hence $\left[\begin{array}{cc}a & \lambda_{j} \\ d \lambda_{j} & b\end{array}\right] \rightarrow\left[\begin{array}{cc}a & \lambda_{j}(P) \\ d \lambda_{j}(P) & b\end{array}\right]$ when $\varepsilon \rightarrow 0$, we see that all the interior points of $\mathcal{E}\left(\lambda_{1}(P), \ldots, \lambda_{m}(P)\right)$ lie in $W_{c}(A)$.

Conversely, suppose $\mu \in W_{c}(A)$. Then there exist a $k$-dimensional subspace $\mathcal{V}_{1}$ of $\mathcal{H}$ and $X: \mathcal{V}_{1} \rightarrow \mathcal{H}$ such that $X^{*} X=I_{k}$ and $\mu \in W_{c}\left(X^{*} A X\right)$. By Lemma 4.6, there are a $(4 k+\ell)$ dimensional subspace $\mathcal{V}_{2}$, containing $\mathcal{V}_{1}$ and $Y: \mathcal{V}_{2} \rightarrow \mathcal{H}$ such that $Y^{*} Y=I_{\mathcal{V}_{2}}$ and $Y^{*} A Y$ has an operator matrix

$$
\left[\begin{array}{cc}
a I_{2 k} & P^{\prime} \\
d P^{\prime} & b I_{2 k}
\end{array}\right] \oplus \gamma I_{\ell} \in M_{4 k+\ell},
$$

where $P^{\prime}=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{2 k}\right)$ is a $2 k$-dimensional compression of $P$, with $\lambda_{1} \geq \cdots \geq \lambda_{2 k}$ and $\lambda_{i} \leq \lambda_{i}(P)$ for $1 \leq i \leq 2 k$. Since $X^{*} A X$ is a compression of $Y^{*} A Y$, we have $\mu \in W_{c}\left(X^{*} A X\right) \subseteq$ $W_{c}\left(Y^{*} A Y\right)$. By the finite dimensional result, we have

$$
\mu \in W_{c}\left(Y^{*} A Y\right)=\mathcal{E}\left(\lambda_{1}, \ldots, \lambda_{p}\right) \subseteq \mathcal{E}\left(\lambda_{1}(P), \ldots, \lambda_{p}(P)\right) .
$$

So, we have shown that

$$
\operatorname{int}\left(\mathcal{E}\left(\lambda_{1}(P), \ldots, \lambda_{p}(P)\right)\right) \subseteq W_{c}(A) \subseteq \mathcal{E}\left(\lambda_{1}(P), \ldots, \lambda_{p}(P)\right)
$$

If $P$ has a $p$-dimensional compression $\tilde{P}=\operatorname{diag}\left(\lambda_{1}(P), \ldots, \lambda_{p}(P)\right)$, then $A$ is a compression of the form

$$
\left[\begin{array}{cc}
I_{p} & \tilde{P} \\
d \tilde{P} & I_{p}
\end{array}\right]
$$

Thus, $\mathcal{E}\left(\lambda_{1}(P), \ldots, \lambda_{p}(P)\right) \subseteq W_{c}(A)$. Hence, $W_{c}(A)=\mathcal{E}$.
Suppose $P$ does not have an compression of the above form. Then the $W_{c}(A)$ can only include boundary points of $\mathcal{E}$ if $A$ satisfies condition (2) in Theorem 3.1, i.e., condition (2) in Lemma 3.2. Moreover, the two boundary points of $\mathcal{E}$ included in $W_{c}(A)$ will be the points described in (2).

In Theorem 4.8, if $\lambda_{j}(P)=0$ for some $j \leq t$, then $W\left(C_{j}\right)+\cdots+W\left(C_{t}\right)$ becomes a line segment joining

$$
a \sum_{i=j}^{t} c_{i}+b \sum_{i=k-t+1}^{k-m+1} c_{i} \text { and } b \sum_{i=j}^{t} c_{i}+a \sum_{i=k-t+1}^{k-m+1} c_{i} .
$$

Thus, $W_{c}(A)$ is a sum of $j-1$ nondegenerate elliptical disks with one line segment. Therefore, we have the following corollary.

Corollary 4.9. Let $c=\left(c_{1}, \ldots, c_{k}\right)$ and $A \in \mathcal{B}(\mathcal{H})$ satisfy the hypotheses of Theorem 4.8. Then the boundary of $\mathbf{c l}\left(W_{c}(A)\right)$ is differentiable. If $\lambda_{j}(P)=0$ for some $j \leq \min \{m, r\}$, then there are exactly two flat portions on the boundary. Otherwise, there is no flat portion on the boundary.

Corollary 4.10. Suppose $A$ is a generalized quadratic operator with an operator matrix in the form described in Theorem 3.1. Let $t=\min \{k, r\}$ and

$$
B_{j}=\left[\begin{array}{cc}
a & c \lambda_{j}(P) \\
d \lambda_{j}(P) & b
\end{array}\right] \quad j=1, \ldots, t .
$$

Let

$$
\mathcal{E}=W\left(B_{1}\right)+\cdots+W\left(B_{t}\right)+(k-t) \gamma .
$$

(a) Suppose $k \leq r+s$. If $P: \mathcal{H}_{1} \rightarrow \mathcal{H}_{1}$ has a compression $\operatorname{diag}\left(\lambda_{1}(P), \ldots, \lambda_{t}(P)\right)$, then $W_{k}(A)=\mathcal{E}$. Otherwise, $W_{k}(A)=\operatorname{int}(\mathcal{E}) \cup\{t a+(k-t) \gamma, t b+(k-t) \gamma\}$.
(b) If $k>r+s$, then $W_{k}(A)$ equals

$$
W\left(B_{1}\right)+\cdots+W\left(B_{2 r+s-k}\right)+(k-r-s)(a+b)+s \gamma .
$$

4.4. Davis-Wielandt shell. Define the Davis-Wielandt shell of $A$ by

$$
D W(A)=\{(\langle A x, x\rangle,\langle A x, A x\rangle): x \in \mathcal{H},\langle x, x\rangle=1\}
$$

see $[10,11,30]$. Evidently, the projection of the set $D W(A)$ on the first co-ordinate is $W(A)$. So, $D W(A)$ captures more information about the operator $A$ than $W(A)$. For example, in the finite dimensional case, normality of operators can be completely determined by the geometrical shape of their Davis-Wielandt shells, namely, $A \in M_{n}$ is normal if and only if $D W(A)$ is a polyhedron in $\mathbb{C} \times \mathbb{R}$ identified with $\mathbb{R}^{3}$. Suppose $A \in \mathcal{B}(\mathcal{H})$. It is known that if $\operatorname{dim} \mathcal{H} \geq 3$ then $D W(A)$ is always convex. Suppose $A=\left[\begin{array}{ll}a & c \\ d & b\end{array}\right] \in M_{2}$ is non-scalar and has eigenvalues $\lambda_{1}, \lambda_{2}$. If $A$ is normal, then $D W(A)$ degenerates to the line segment joining the points $\left(\lambda_{1},\left|\lambda_{1}\right|^{2}\right)$ and $\left(\lambda_{2},\left|\lambda_{2}\right|^{2}\right)$; otherwise, $D W(A)$ is an $D W(A)$ is an ellipsoid (without its interior) centered at $\left(\left(\lambda_{1}+\lambda_{2}\right) / 2, \operatorname{tr}\left(A^{*} A\right) / 2\right)$.

Suppose $\operatorname{dim} \mathcal{H} \geq 3$. In [22], a complete description of $D W(A)$ for a quadratic operator $A$ was given. For generalized quadratic operators, we have the following.

Theorem 4.11. Suppose $\operatorname{dim} \mathcal{H} \geq 3$ and $A \in \mathcal{B}(\mathcal{H})$ is a generalized quadratic operator with operator matrix in the form in Theorem 3.1 and let $A_{p}$ be defined as in (2.2). Let $S=\left\{\left(\gamma,|\gamma|^{2}\right)\right\}$ if $\gamma I_{s}$ is non-trivial and $S=\emptyset$ otherwise. Suppose $\sigma(P)=\sigma_{1}(P) \cup \sigma_{2}(P)$, where $\sigma_{1}(P)$ is the set of eigenvalues of $P$. Then

$$
\begin{equation*}
D W(A)=\operatorname{conv}\left[\bigcup_{p \in \sigma_{1}(P)} D W\left(A_{p}\right) \cup \bigcup_{p \in \sigma_{2}(P)} \operatorname{int}\left(D W\left(A_{p}\right)\right) \cup S\right] \cup \mathcal{L} \tag{4.7}
\end{equation*}
$$

where
(i) $\mathcal{L}=\left\{\left(\mu,|\mu|^{2}+\eta^{2}\right): \mu \in\{a, b\}, \eta \in W(|c| P \oplus|d| P)\right\} \subseteq \partial W(A)$ if $a, b \in \partial D W(A)$, i.e., condition (2) in Theorem 3.1 holds, and
(ii) $\mathcal{L}=\emptyset$ otherwise, i.e., condition (1) or (3) in Theorem 3.1 holds.

Consequently,

$$
\begin{equation*}
\mathbf{c l}(D W(A))=\mathbf{c l}\left(\operatorname{conv}\left[\bigcup_{p \in \sigma(P)} D W\left(A_{p}\right) \cup S\right]\right) \tag{4.8}
\end{equation*}
$$

As pointed out by the referee, for a generalized quadratic operator $A$, the $q$-numerical range is just an elliptical disk (with all or without any boundary points), see Theorem 4.1, whereas the description of the Davis Wielandt shell is more involved. This shows the subtlety of the study.

To prove Theorem 4.11, we need the following lemmas.
Lemma 4.12. Let $A_{p}$ be defined as in (2.2). If $T \in \mathcal{B}(\mathcal{H})$ such that $W(T) \subseteq W\left(A_{p}\right)$, then for any unit vector $x \in \mathcal{H}$ there is a unit vector $y \in \mathbb{C}^{2}$ such that $\langle T x, x\rangle=\left\langle A_{p} y, y\right\rangle$ and $\|T x\| \leq\left\|A_{p} y\right\|$.

Proof. Suppose $W(T) \subseteq W\left(A_{p}\right)$. By Corollary 3.3, there exists an isometry $U$ such that $T=U^{*}\left(I \otimes A_{p}\right) U$. For any unit vector $x \in \mathcal{H}$, if $z=U x$, then $\mu=\langle T x, x\rangle=\left\langle\left(I \otimes A_{p}\right) z, z\right\rangle$ and $\|T x\|^{2} \leq\left\|\left(I \otimes A_{p}\right) z\right\|^{2}=\nu$. Since $(\mu, \nu) \in D W\left(I \otimes A_{p}\right) \subseteq \operatorname{conv} D W\left(A_{p}\right)$ and $D W(A)$ is either a line segment or an ellipsoid (without its interior), there exists ( $\mu, \hat{\nu}) \in D W\left(A_{p}\right)$ such that $\hat{\nu} \geq \nu$. The result follows.

Lemma 4.13. Suppose

$$
B_{k}=\left[\begin{array}{cc}
a I & c P_{k}  \tag{4.9}\\
d P_{k} & b I
\end{array}\right] \in \mathcal{B}\left(\mathcal{V}_{k} \oplus \mathcal{V}_{k}\right)
$$

where $P_{k}$ is one of the summand in $P=\oplus_{k=1}^{n} P_{k}$ described in Theorem 2.1. Then

$$
\begin{equation*}
D W\left(B_{k}\right) \subseteq \operatorname{conv}\left[\cup_{p \in \sigma\left(P_{k}\right)} D W\left(A_{p}\right)\right] . \tag{4.10}
\end{equation*}
$$

Furthermore, if $(\mu, \nu) \in D W\left(B_{k}\right)$ is an extreme point of $D W\left(B_{k}\right)$, then $(\mu, \nu) \in D W\left(A_{p}\right)$ for some eigenvalue $p$ of $P_{k}$.

Proof. First, we prove (4.10). The result is clear if $P_{k}$ is a singleton. So, we assume that $P_{k}$ is a non-degenerate interval. Hence, $\mathcal{V}_{k}$ is infinite dimensional.

Suppose $(\mu, \nu) \in D W\left(B_{k}\right)$. Then there exists a unit vector $x \in \mathcal{V}_{k} \oplus \mathcal{V}_{k}$, such that $(\mu, \nu)=$ $\left(\left\langle B_{k} x, x\right\rangle,\left\|B_{k} x\right\|^{2}\right)$. Let $x=\left[\begin{array}{c}\cos \theta x_{1} \\ \sin \theta x_{2}\end{array}\right]$ for some unit vectors, $x_{1}, x_{2} \in \mathcal{V}_{k}$. Let $\left\{u_{1}, u_{2}\right\} \in \mathcal{V}_{k}$ be an orthogonal normal family such that $x_{1}, x_{2} \in \operatorname{span}\left\{u_{1}, u_{2}\right\}$. Let

$$
U=\left[\begin{array}{ll}
u_{1} & u_{2}
\end{array}\right]: \mathbb{C}^{2} \rightarrow \mathcal{V}_{k} \text { and } X=\left[\begin{array}{cc}
U & 0 \\
0 & U
\end{array}\right]: \mathbb{C}^{4} \rightarrow \mathcal{V}_{k} \oplus \mathcal{V}_{k}
$$

Then $X^{*} X=I_{4}$ and

$$
C_{k}=X^{*} B_{k} X=\left[\begin{array}{cc}
a I_{2} & c U^{*} P_{k} U \\
d U^{*} P_{k} U & b I_{2}
\end{array}\right] \in M_{4}
$$

is a compression of $A$. Clearly, $\sigma\left(U^{*} P_{k} U\right) \subseteq \sigma\left(P_{k}\right)$. Now, suppose $v_{1}, v_{2} \in \mathbb{C}^{2}$ are unit vectors such that $x_{1}=U v_{1}$ and $x_{2}=U v_{2}$. Then for $v=\left[\begin{array}{c}\cos \theta v_{1} \\ \sin \theta v_{2}\end{array}\right]$, we have

$$
\left\langle C_{k} v, v\right\rangle=\left\langle B_{k} x, x\right\rangle=\mu .
$$

Furthermore, if $\tilde{\nu}=\left\|C_{k} v\right\|^{2}$, then $(\mu, \nu) \in D W\left(C_{k}\right)$ with

$$
\begin{aligned}
\nu-\tilde{\nu} & =\left\langle B_{k}^{*} B_{k} x, x\right\rangle-\left\langle C_{k}^{*} C_{k} x, x\right\rangle \\
& =|d|^{2}\left(v_{1}^{*} U^{*} P_{k}^{2} U v_{1}-v_{1}^{*}\left(U^{*} P_{k} U\right)^{2} v_{1}\right)+|c|^{2}\left(v_{2}^{*} U^{*} P_{k}^{2} U v_{2}-v_{2}^{*}\left(U^{*} P_{k} U\right)^{2} v_{2}\right) \\
& =|d|^{2}\left(\left\|P_{k} x_{1}\right\|^{2}-\left\|U^{*} P_{k} x_{1}\right\|^{2}\right)+|c|^{2}\left(\left\|P_{k} x_{2}\right\|^{2}-\left\|U^{*} P_{k} x_{2}\right\|^{2}\right) \\
& \geq 0 .
\end{aligned}
$$

Let $p=\max \sigma\left(P_{k}\right)$. Then by Lemma 3.2, $W\left(B_{k}\right) \subseteq W\left(A_{p}\right)$. By Lemma 4.12, there is $\eta \geq \nu$ such that $(\mu, \eta) \in D W\left(A_{p}\right)$. Since $C_{k}$ is unitarily similar to $A_{q_{1}} \oplus A_{q_{2}}$ for some $q_{1}, q_{2} \in \sigma\left(P_{k}\right)$, it follows that $(\mu, \eta) \in D W\left(A_{p}\right),(\mu, \tilde{\nu}) \in \operatorname{conv}\left[D W\left(A_{q_{1}}\right) \cup D W\left(A_{q_{2}}\right)\right]$, and

$$
(\mu, \nu) \in \operatorname{conv}\{(\mu, \eta),(\mu, \tilde{\nu})\} \subseteq \operatorname{conv}\left[\cup_{q \in \sigma\left(P_{k}\right)} D W\left(A_{q}\right)\right] .
$$

Suppose $(\mu, \nu) \in D W\left(B_{k}\right)$ is an extreme point of $D W\left(B_{k}\right)$. Following the argument in the previous paragraph, we see that

$$
\text { (1) } \nu=\tilde{\nu}, \quad \text { or } \quad \text { (2) } \eta=\nu \text { and } p \in\left\{q_{1}, q_{2}\right\} .
$$

Suppose (1) holds. Let $\mathcal{S}$ be the subspace spanned by $\left\{u_{1}, u_{2}\right\}$. Analyzing the equality condition for $\nu-\tilde{\nu} \geq 0$, we see that both $P_{k} x_{1}$ and $P_{k} x_{2}$ lie in $\mathcal{S}$. So $\mathcal{S}$ is a reducing subspace of $P_{k}$, and consequently, a reducing subspace of $P$. We can choose $u_{1}$ and $u_{2}$ to be eigenvectors of $P$ corresponding to eigenvalues $p_{1}$ and $p_{2}$ respectively. Then $C_{k}$ is a direct summand of $A$ and unitarily similar to $A_{p_{1}} \oplus A_{p_{2}}$. Since $(\mu, \nu)$ is an extreme point of $\mathbf{c l}(D W(A))$, it must lie in $D W\left(A_{p_{i}}\right)$ for some $i=1,2$.

Suppose (2) holds, and assume that $p=q_{1}$. In particular, we may assume that $A_{p}=\left[\begin{array}{cc}a & c p \\ d p & b\end{array}\right]$ is a principal submatrix of $C_{k}$, and $P$ is unitarily similar to $[p] \oplus \hat{P}$. Thus, $p$ is an eigenvalue of $P$ and $(\mu, \nu) \in D W\left(A_{p}\right)$.

## Proof of Theorem 4.11.

We will first prove the equality (4.8) and then use the result to prove (4.7).
Step 1. We prove the inclusion "?" of (4.8). Suppose $(\mu, \nu) \in D W\left(A_{p}\right)$, with $p \in \sigma(P)$. Then there is a sequence of unit vectors $\left\{x_{m}\right\}$ in $\mathcal{H}_{1}$ such that $\left\|P x_{m}-p x_{m}\right\| \rightarrow 0$ so that $\left\langle P x_{m}, x_{m}\right\rangle=p_{m} \rightarrow p$. Thus, for each $m, P$ is unitarily similar to

$$
P_{m}=\left[\begin{array}{ccc}
p_{m} & \delta_{m} & 0 \\
\delta_{m} & * & * \\
0 & * & *
\end{array}\right],
$$

where $\delta_{m} \geq 0$ and $\delta_{m} \rightarrow 0$. Hence, $A$ is untarily similar to

$$
T_{m}=\left[\begin{array}{cc}
a I_{r} & c P_{m} \\
d P_{m} & b I_{r}
\end{array}\right] \oplus \gamma I_{s} .
$$

Then, for any unit vector $u=\left[\begin{array}{l}u_{1} \\ u_{2}\end{array}\right] \in \mathbb{C}^{2}$, we can extend it to $\tilde{u}_{1}=\left[\begin{array}{c}u_{1} \\ 0\end{array}\right], \tilde{u}_{2}=\left[\begin{array}{c}u_{2} \\ 0\end{array}\right] \in \mathcal{H}_{1}$ so that for $\tilde{u}=\left[\begin{array}{c}\tilde{u}_{1} \\ \tilde{u}_{2} \\ 0\end{array}\right] \in \mathcal{H}_{1} \oplus \mathcal{H}_{1} \oplus \mathcal{H}_{2}$. It follows that

$$
\left(\left\langle T_{m} \tilde{u}, \tilde{u}\right\rangle,\left\|T_{m} \tilde{u}\right\|^{2}\right) \rightarrow\left(\left\langle A_{p} u, u\right\rangle,\|A u\|^{2}\right) .
$$

Thus, $\mathbf{c l} D W(A) \supseteq D W\left(A_{p}\right)$. By convexity of $D W(A)$, we get the inclusion "?" of (4.8).
Step 2. We prove the inclusion " $\subseteq$ " of (4.8). Since $D W(X \oplus Y)=\operatorname{conv}\{D W(X) \cup D W(Y)\}$, we may assume that $\gamma I_{s}$ is vacuous. So, it suffices to show that

$$
D W\left(B_{k}\right) \subseteq \operatorname{conv}\left[\cup_{p \in \sigma\left(P_{k}\right)} D W\left(A_{p}\right)\right],
$$

which follows from Lemma 4.13.
By Step 1 and Step 2, we get the equality (4.8). Next, we turn to the equality (4.7).

Step 3. We prove the inclusion " $\supseteq$ " of (4.7). By the description of $\mathbf{c l}(D W(A))$ and the fact that $D W(A)$ has non-empty relative interior as $c d P \neq 0$, we see that

$$
D W(A) \supseteq \operatorname{conv}\left[\bigcup_{p \in \sigma_{1}(P)} D W\left(A_{p}\right) \cup \bigcup_{p \in \sigma_{2}(P)} \operatorname{int}\left(D W\left(A_{p}\right)\right) \cup S\right] .
$$

Suppose condition (2) of Theorem 3.1 holds. We will show that $\mathcal{L} \subseteq D W(A)$. Note that one can apply the construction in Step 1 to get the sequence of operators $\left\{T_{m}\right\}$ so that

$$
\left\{p_{m}\right\} \rightarrow \hat{p} \in\{\max \sigma(P), \min \sigma(P)\} .
$$

Let $e_{1}=\left[\begin{array}{l}1 \\ 0\end{array}\right] \in \mathcal{H}_{1}$. Then for $v_{1}=\left[\begin{array}{c}e_{1} \\ 0 \\ 0\end{array}\right], v_{2}=\left[\begin{array}{c}0 \\ e_{1} \\ 0\end{array}\right] \in \mathcal{H}_{1} \oplus \mathcal{H}_{1} \oplus \mathcal{H}_{2}$, we have

$$
\left(\left\langle T_{m} v_{1}, v_{1}\right\rangle,\left\|T_{m} v_{1}\right\|^{2}\right) \rightarrow\left(a,|a|^{2}+|d \hat{p}|^{2}\right) \quad \text { and } \quad\left(\left\langle T_{m} v_{2}, v_{2}\right\rangle,\left\|T_{m} v_{2}\right\|^{2}\right) \rightarrow\left(b,|b|^{2}+|c \hat{p}|^{2}\right) .
$$

Note that there is a unitary $U=\left(u_{i j}\right) \in M_{2}$ such that

$$
U\left[\begin{array}{cc}
a & c p_{m} \\
d p_{m} & b
\end{array}\right] U^{*}=\left[\begin{array}{cc}
b & c p_{m} \\
d p_{m} & a
\end{array}\right] .
$$

Extend $U$ to

$$
V=\left[\begin{array}{cccc}
u_{11} & 0 & u_{12} & 0 \\
0 & I_{r-1} & 0 & 0 \\
u_{21} & 0 & u_{22} & 0 \\
0 & 0 & 0 & I_{r-1}
\end{array}\right] \oplus I_{s} \in \mathcal{B}\left(\mathcal{H}_{1} \oplus \mathcal{H}_{1} \oplus \mathcal{H}_{2}\right)
$$

Then

$$
\hat{T}_{m}=V T_{m} V^{*}=\left[\begin{array}{cccccc}
b & d u_{12} \delta_{m} & 0 & c p_{m} & c u_{11} \delta_{m} & 0 \\
c \bar{u}_{12} \delta_{m} & a & 0 & c u_{22} \delta_{m} & * & * \\
0 & 0 & a I_{r-2} & 0 & * & * \\
d p_{m} & d u_{22} \delta_{m} & 0 & a & c u_{21} \delta_{m} & 0 \\
d \bar{u}_{11} \delta_{m} & * & * & d \bar{u}_{21} \delta_{m} & b & 0 \\
0 & * & * & 0 & 0 & b I_{r-2}
\end{array}\right] \oplus \gamma I_{s} .
$$

For $e_{1}=\left[\begin{array}{l}1 \\ 0\end{array}\right] \in \mathcal{H}_{1}$, and $v_{1}=\left[\begin{array}{c}e_{1} \\ 0 \\ 0\end{array}\right], v_{2}=\left[\begin{array}{c}0 \\ e_{1} \\ 0\end{array}\right] \in \mathcal{H}_{1} \oplus \mathcal{H}_{1} \oplus \mathcal{H}_{2}$, we have

$$
\left(\left\langle\hat{T}_{m} v_{1}, v_{1}\right\rangle,\left\|\hat{T}_{m} v_{1}\right\|^{2}\right) \rightarrow\left(b,|b|^{2}+|d \hat{p}|^{2}\right) \quad \text { and } \quad\left(\left\langle\hat{T}_{m} v_{2}, v_{2}\right\rangle,\left\|\hat{T}_{m} v_{2}\right\|^{2}\right) \rightarrow\left(a,|a|^{2}+|c \hat{p}|^{2}\right) .
$$

Hence, $\left(\mu,|\mu|^{2}+\eta^{2}\right) \in \mathbf{c l}(D W(A))$, for $\mu \in\{a, b\}$ and $\eta \in\{\max \sigma(P), \min \sigma(P)\}$. By the convexity of $D W(A)$, the relative interior of $\mathcal{L} \subseteq D W(A)$. Moreover, if $\hat{p} \in\{\max \sigma(A), \min \sigma(A)\}$ and $\hat{p} \in W(P)$, then there is a unit vector $u \in \mathcal{H}_{1}$ such that $\|P\|=\langle P u, u\rangle=\|P u\|\|u\| \leq\|P\|$. As a result, $A$ is unitarily similar to $A_{\hat{p}} \oplus \hat{A}$. In particular, $D W\left(A_{\hat{p}}\right) \subseteq D W(A)$ so that $\left(\mu,|\mu|^{2}+\hat{p}^{2}\right) \in D W(A)$. Thus, we see that $\mathcal{L} \subseteq D W(A)$.

Step 4. We prove the inclusion " $\subseteq$ " of (4.7). Suppose $\left(\mu_{0}, \nu_{0}\right) \in \partial D W(A) \cap D W(A)$. Then there is a support plane $\mathbf{P}$ of the convex set $\mathbf{c l}(D W(A))$ passing through $\left(\mu_{0}, \nu_{0}\right) \subseteq \mathbb{C} \times \mathbb{R} \sim \mathbb{R}^{3}$.

So, there exist real numbers $f, g, h, \lambda$ such that

$$
\lambda=\operatorname{Re}\left((f+i g) \mu_{0}\right)+h \nu_{0}=\max \{\operatorname{Re}((f+i g) \mu)+h \nu:(\mu, \nu) \in D W(A)\} .
$$

Now, if $\mathbf{P}$ is a support plane for $D W\left(A_{p}\right)$, then there is a unit vector $x \in \mathbb{C}^{2}$ such that

$$
\lambda=\operatorname{Re}\left((f+i g)\left(x^{*} A_{p} x\right)\right)+h x^{*} A_{p}^{*} A_{p} x \geq \operatorname{Re}\left((f+i g)\left(y^{*} A_{p} y\right)\right)+h y^{*} A_{p}^{*} A_{p} y, \quad y \in \mathbb{C}^{2}, y^{*} y=1
$$

Thus, $\lambda$ is the largest eigenvalue of

$$
L_{p}=\operatorname{Re}\left((f+i g) A_{p}\right)+h A_{p}^{*} A_{p}=\left[\begin{array}{cc}
a_{11}+h|d|^{2} p^{2} & a_{12} p \\
a_{21} p & a_{22}+h|c|^{2} p^{2}
\end{array}\right]
$$

and hence

$$
\begin{equation*}
0=\operatorname{det}\left(\lambda I-L_{p}\right)=|h c d|^{2} p^{4}+t_{1} p^{2}+t_{2} \tag{4.11}
\end{equation*}
$$

where $t_{1}=-h\left[\left(\lambda-a_{11}\right)|c|^{2}+\left(\lambda-a_{22}\right)|d|^{2}\right]-a_{12} a_{21}, t_{2}=\left(\lambda-a_{11}\right)\left(\lambda-a_{22}\right) \in \mathbb{R}$. If $h \neq 0$, then $p^{2}=\left[-t_{1} \pm \sqrt{t_{1}^{2}-4 t_{2}|h c d|^{2}}\right] / 2|h c d|^{2}$, and there are at most two distinct $p \geq 0$ satisfying (4.11). If $h=0$ and $t_{1} \neq 0$, then $p^{2}=-t_{2} / t_{1}$ and there is at most one $p \geq 0$ satisfying (4.11). If $h=t_{1}=0$ and $t_{2} \neq 0$, then there is no $p \geq 0$ satisfying (4.11). Finally, if $h=t_{1}=t_{2}=0$, we see that (4.11) holds for all $p \geq 0$.

Note that for any unit vector $x \in \mathcal{H}$, we can decomposed $x$ according to the range space of $B_{k}$. It follows that every element $(\mu, \nu)=\left(\langle A x, x\rangle,\|A x\|^{2}\right)$ in $D W(A)$ can be written as $(\mu, \nu)=\sum_{k} t_{k}\left(\mu_{k}, \nu_{k}\right)$ with $\left(\mu_{k}, \nu_{k}\right) \in D W\left(B_{k}\right)$, where the sum could be infinite. In any event, if $\left(\mu_{0}, \nu_{0}\right) \in \partial D W(A) \cap \mathbf{P}$, where $\mathbf{P}$ is a support plane of $D W(A)$, then $\left(\mu_{0}, \nu_{0}\right)$ is a convex combination of elements in $D W\left(A_{p}\right) \cap \mathbf{P}$ for some $p \in \sigma(A)$ such that $D W\left(A_{p}\right) \cap P \neq \emptyset$. By the previous discussion, we have the following three cases.

Case 1 If the support plane $\mathbf{P}$ of $D W(A)$ has intersection with $D W\left(A_{p}\right)$ for only one $p \in \sigma(P)$ and if $p \in \sigma\left(P_{k}\right)$, then we have $\mathbf{P} \cap D W\left(B_{k}\right)=\left\{\left(\mu_{0}, \nu_{0}\right)\right\}$. Hence, $\left(\mu_{0}, \nu_{0}\right)$ is an extreme point of $D W\left(B_{k}\right)$. By Lemma 4.13, $p$ is an eigenvalue of $P_{k}$.

Case 2 If the support plane $\mathbf{P}$ of $D W(A)$ has intersection with $D W\left(A_{p}\right)$ and $D W\left(A_{q}\right)$ for two different $p, q \in \sigma(P)$, then $\mathbf{P} \cap D W(A)=\mathbf{c o n v}\left[\left(\mathbf{P} \cap D W\left(A_{p}\right)\right) \cup\left(\mathbf{P} \cap D W\left(A_{q}\right)\right)\right]$. If $\left(\mu_{0}, \nu_{0}\right)$ is one of the two points in the set $\left(\mathbf{P} \cap D W\left(A_{p}\right)\right) \cup\left(\mathbf{P} \cap D W\left(A_{q}\right)\right)$, then $\left(\mu_{0}, \nu_{0}\right)$ is the extreme point of $D W\left(B_{k}\right)$ for some $B_{k}$ defined as in Lemma 4.13. If $\left(\mu_{0}, \nu_{0}\right)$ is a non-trivial combination of the two points in the set $\left(\mathbf{P} \cap D W\left(A_{p}\right)\right) \cup\left(\mathbf{P} \cap D W\left(A_{q}\right)\right)$, then the two points must lie in $D W(A)$. Moreover, they must be extreme points of $D W\left(B_{k}\right)$ and $D W\left(B_{\ell}\right)$, where $B_{k}$ and $B_{\ell}$ defined as in Lemma 4.13 such that $p \in \sigma\left(P_{k}\right)$ and $q \in \sigma\left(P_{\ell}\right)$. By Lemma 4.13, $p, q$ are eigenvalues of $P_{k}$ and $P_{\ell}$, respectively.

In both Cases 1 and 2, we conclude that $\left(\mu_{0}, \nu_{0}\right) \in \boldsymbol{\operatorname { c o n v }}\left[\bigcup_{p \in \sigma_{1}(P)} D W\left(A_{p}\right)\right]$.
Case 3 If $\left(\mu_{0}, \nu_{0}\right) \in D W\left(A_{p}\right)$ for all $p \geq 0$, then $h=0$ so that $\mu_{0} \in \cap_{p \geq 0} W\left(A_{p}\right)$. It follows that condition (2) of Lemma 3.2 holds, and $\mu_{0} \in\{a, b\}$. By the argument in Step 3, we have $\left(\mu_{0}, \nu_{0}\right) \in \mathcal{L} \subseteq \operatorname{conv}\left[\bigcup_{p \in \sigma_{1}(P)} D W\left(A_{p}\right)\right]$. The result follows.

## 5. Additional remarks and further research

We may extend most of our results to $A \in \mathcal{B}(\mathcal{H})$ of the form in (2.1) without requiring that $\gamma \in$ $\{a, b\}$, using the simple fact that $\sigma(X \oplus Y)=\sigma(X) \cup \sigma(Y), W(X \oplus Y)=\operatorname{conv}\{W(X) \cup W(Y)\}$, $\|X \oplus Y\|=\max \{\|X\|,\|Y\|\}, D W(X \oplus Y)=\operatorname{conv}\{D W(X) \cup D W(Y)\}$, etc.

We may consider $A \in \mathcal{B}(\mathcal{H})$ with an operator matrix of the form

$$
\left[\begin{array}{cc}
a I & R  \tag{5.1}\\
S & b I
\end{array}\right]
$$

so that $R S$ and $S R$ are normal. If $A$ is a Hilbert-Schmidt operator, i.e., $A^{*} A$ has a bounded trace, then $\mathcal{H}=\mathcal{H}_{1} \oplus \mathcal{H}_{1} \oplus \mathcal{H}_{2}$ and $A$ is unitarily similar to

$$
\left[\begin{array}{ll}
a I_{r} & D_{1}  \tag{5.2}\\
D_{2} & b I_{r}
\end{array}\right] \oplus \gamma I_{s},
$$

where $D_{1}, D_{2} \in \mathcal{B}\left(\mathcal{H}_{1}\right)$ are commuting normal operators, and $\gamma \in\{a, b\}$.
To see this, we may first find $U, V$ such that $U R V$ has the form (a), (b), or (c) as in the proof of Theorem 2.1, we see that $A$ is unitarily similar to

$$
\left[\begin{array}{cc}
a I_{r} & P \\
Q & b I_{r}
\end{array}\right] \oplus \gamma I_{s}
$$

where $P, Q \in \mathcal{B}\left(\mathcal{H}_{1}\right)$ with $\gamma \in\{a, b\}$. We may further assume that $P=P_{1} \oplus 0$ so that $\operatorname{ker}\left(P_{1}\right)=$ $\{0\}$. Note that $P_{1}$ is compact and has discrete eigenvalues. Thus, we may assume that $P_{1}$ is a direct sum of $p_{j} I_{r_{j}}$ such that $p_{1}>p_{2}>\cdots$. Note that $P Q$ and $Q P$ are normal. Thus, $Q=Q_{1} \oplus 0$ so that $P_{1} Q_{1}$ and $Q_{1} P_{1}$ are normal. Now, decompose $Q_{1}$ according to the decomposition of $P_{1}$. Using the equalities $P_{1} Q_{1} Q_{1}^{*} P_{1}=Q_{1}^{*} P_{1} P_{1} Q_{1}$ and $Q_{1} P_{1} P_{1} Q_{1}^{*}=P_{1} Q_{1}^{*} Q_{1} P$, we see that $Q_{1}$ is also a direct sum of $\hat{Q}_{j}$ and each $\hat{Q}_{j}$ is normal. We get the decomposition (5.2).

With the decomposition (5.2), we easily show that $\sigma(A), W(A), \rho(A), w(A),\|A\|$, etc. are completely determined by matrices of the form

$$
A(u, v)=\left[\begin{array}{ll}
a & u \\
v & b
\end{array}\right], \quad(u, v) \in \sigma\left(D_{1}, D_{2}\right)
$$

where $\sigma\left(D_{1}, D_{2}\right)$ is the joint spectrum of $\left(D_{1}, D_{2}\right)$ defined as the collection of $(u, v) \in \mathbb{C} \times \mathbb{C}$ such that $\left\|R x_{m}-u x_{m}\right\|+\left\|S x_{m}-v x_{m}\right\| \rightarrow 0$ for a sequence of unit vectors $\left\{x_{m}\right\} \subseteq \mathcal{H}_{1}$. For instance, one has

$$
\mathbf{c l}\left(W_{q}(A)\right)=\mathbf{\operatorname { c o n v }}\left[\mathbf{c l}\left(\cup_{(u, v) \in \sigma\left(D_{1}, D_{2}\right)} W_{q}(A(u, v))\right)\right], \quad q \in[0,1],
$$

and

$$
\mathbf{c l}(D W(A))=\mathbf{\operatorname { c o n v }}\left[\mathbf{c l}\left(\cup_{(u, v) \in \sigma\left(D_{1}, D_{2}\right)} D W(A(u, v))\right)\right] .
$$

If there is $(\tilde{u}, \tilde{v}) \in \sigma\left(D_{1}, D_{2}\right)$ such that

$$
W(A(u, v)) \subseteq W(A(\tilde{u}, \tilde{v})) \quad \text { for all }(u, v) \in \sigma\left(D_{1}, D_{2}\right)
$$

then we have $\mathbf{c l}(W(A))=W(A(\tilde{u}, \tilde{v})), \mathbf{c l}\left(W_{q}(A)\right)=W\left(A_{q}(\tilde{u}, \tilde{v})\right), A$ has a dilation of the form $I \otimes A(\tilde{u}, \tilde{v}),\|A\|=\|A(\tilde{u}, \tilde{v})\|$, and the matricial range result holds, namely, $B \in W^{n}(A)$ if and
only if $W(B) \subseteq W(A(\tilde{u}, \tilde{v}))$. Furthermore, if $A$ is unitarily similar to a direct sum of $A\left(u_{j}, v_{j}\right)$ for $j=1, \ldots, m$, and $\hat{A}$ such that

$$
W(\hat{A}) \subseteq W\left(A\left(u_{m}, v_{m}\right)\right) \subseteq W\left(A\left(u_{m-1}, v_{m-1}\right)\right) \subseteq \cdots \subseteq W\left(A\left(u_{1}, v_{1}\right)\right),
$$

then we can obtain results for the rank- $k$ numerical range, essential numerical range (if $m=\infty$ ), $c$-numerical range.

It would be nice to show that the decomposition (5.2) holds also for general operators with an operator matrix of the form (5.1).

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