

OPERATORS OF THE FORM $PAQ - QAP$

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1. In this note the Hilbert spaces under consideration are complex, and the operators referred to are bounded, linear operators. If \mathfrak{H} is a Hilbert space, then the algebra of all operators on \mathfrak{H} is denoted by $\mathfrak{L}(\mathfrak{H})$.

It is known (**1**) that if \mathfrak{H} is any Hilbert space, then the class of *commutators* on \mathfrak{H} , i.e., the class of all operators that can be written in the form $PQ - QP$ for some $P, Q \in \mathfrak{L}(\mathfrak{H})$, can be exactly described. A similar problem is that of characterizing all operators on \mathfrak{H} that can be written in the form $PAQ - QAP$ for some $P, A, Q \in \mathfrak{L}(\mathfrak{H})$. If no restrictions are placed on the operators P, A , and Q , it is relatively easy to see that for $\dim \mathfrak{H} > 1$, every operator in $\mathfrak{L}(\mathfrak{H})$ can be written in this form. (A very brief and pretty proof of this fact for infinite-dimensional \mathfrak{H} was shown to us by Paul Federbush; it is reproduced in Remark 3.3.)

Since every commutator $PQ - QP$ is automatically a commutator of invertible operators by virtue of the identity

$$PQ - QP = (P + \lambda)(Q + \mu) - (Q + \mu)(P + \lambda),$$

valid for every pair of scalars λ and μ , it is natural to ask which operators can be written in the form $PAQ - QAP$ with *invertible* P, A , and Q . This problem is somewhat more difficult, and it is the purpose of this note to furnish the solution by proving the following theorem.

THEOREM. *If \mathfrak{H} is a Hilbert space of dimension greater than one, and T is any operator on \mathfrak{H} , then there exist invertible operators P, A, Q on \mathfrak{H} satisfying $T = PAQ - QAP$.*

This theorem settles a question posed to us by Olga Taussky-Todd, to whom we are indebted for several interesting conversations.

The proof of the theorem splits naturally into cases depending on the dimension of \mathfrak{H} . In the finite-dimensional case, the proof depends on the following lemmas and (**4**, Theorem III).

LEMMA 1.1. *If $T \neq 0$ is an operator on an n -dimensional complex Hilbert space \mathfrak{H} ($1 < n < \infty$), then there exists an orthonormal basis for \mathfrak{H} relative to which the matrix $(\alpha_{ij})_{i,j=1}^n$ of T satisfies $\alpha_{11} \neq 0 \neq \alpha_{22}$.*

Proof. Consider the numerical range (or field of values) $W(T)$ of T . If $W(T)$ consists of a single point $\{\lambda\}$, then T is the scalar operator $T = \lambda I$ and

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the result is obvious. Otherwise, $W(T)$ contains at least two points, and thus the line segment joining them. Hence, $W(T)$ contains a number $\alpha_{11} \neq 0$ such that $\alpha_{11} \neq \text{trace } T$. Let x_1 be a unit vector such that $(Tx_1, x_1) = \alpha_{11}$, and extend $\{x_1\}$ to an orthonormal basis $\{x_1, \dots, x_n\}$ for \mathfrak{H} . Since $\alpha_{11} \neq \text{trace } T$, there must be some k ($2 \leq k \leq n$) such that $(Tx_k, x_k) = \beta \neq 0$. If we now interchange x_2 and x_k , then $(Tx_2, x_2) = \beta = \alpha_{22} \neq 0$, and the proof is complete.

LEMMA 1.2. *If T is an operator on an n -dimensional Hilbert space \mathfrak{H} ($1 < n < \infty$), then there exist invertible operators X and Y on \mathfrak{H} such that $\text{determinant } X = \text{determinant } Y$ and such that $T = X - Y$.*

Proof. If $T = 0$, the result is clear. We suppose that $T \neq 0$, and use the preceding lemma to pick a basis for \mathfrak{H} relative to which the matrix (α_{ij}) of T satisfies $\alpha_{11} \neq 0 \neq \alpha_{22}$. We write

$$\begin{pmatrix} \alpha_{11} & \alpha_{12} & \dots & \alpha_{1n} \\ \alpha_{21} & & & \cdot \\ \cdot & & & \cdot \\ \cdot & & & \cdot \\ \alpha_{n1} & \dots & \dots & \alpha_{nn} \end{pmatrix} = \begin{pmatrix} \alpha_{11} + d_1 & & & 0 \\ \alpha_{21} & \alpha_{22} + d_2 & & \\ \cdot & \cdot & \cdot & \\ \cdot & \cdot & \cdot & \\ \alpha_{n1} & & & \alpha_{nn} + d_n \end{pmatrix} - \begin{pmatrix} d_1 & -\alpha_{12} \dots -\alpha_{1n} \\ & d_2 & \cdot \\ & & \cdot & \cdot \\ 0 & & & d_n \end{pmatrix}$$

where the numbers d_1, \dots, d_n are to be determined. Let X and Y be the operators having these matrices (relative to the given basis), and note that to complete the proof, it suffices to show that the numbers d_i can be chosen so that

- (1) $d_1 d_2 \dots d_n \neq 0$ and
- (2) $(\alpha_{11} + d_1)(\alpha_{22} + d_2) \dots (\alpha_{nn} + d_n) = d_1 d_2 \dots d_n$.

This amounts to choosing each $d_i \neq 0$ so that (2) is satisfied. If $n = 2$, this is equivalent to choosing non-zero numbers d_1 and d_2 such that

$$\alpha_{11}d_2 + \alpha_{22}d_1 + \alpha_{11}\alpha_{22} = 0,$$

and this is a task that is easily accomplished since $\alpha_{11} \neq 0 \neq \alpha_{22}$. If $n > 2$ we first choose $d_i, 3 \leq i \leq n$, subject only to the conditions $d_i \neq 0 \neq d_i + \alpha_{ii}$. Next, we arrange things so that $\beta = d_2 \dots d_n$ is unequal to

$$\gamma = (\alpha_{22} + d_2) \dots (\alpha_{nn} + d_n).$$

To this end let $\xi = d_3 \dots d_n$ and $\eta = (\alpha_{33} + d_3) \dots (\alpha_{nn} + d_n)$. Then we need $d_2\xi \neq (\alpha_{22} + d_2)\eta$, and since ξ, η , and α_{22} are all non-zero, it is easy to see that we can choose d_2 so as to satisfy this requirement and also to satisfy $d_2 \neq 0 \neq d_2 + \alpha_{22}$. Assume this done; to complete the proof it then suffices to choose $d_1 \neq 0$ so that

$$d_1\beta = (\alpha_{11} + d_1)\gamma,$$

and since $\beta \neq \gamma \neq 0$ and $\alpha_{11} \neq 0$, this is possible.

COROLLARY 1.3. *The theorem is true if \mathfrak{S} is finite-dimensional.*

Proof. Let T be an operator on \mathfrak{S} . By Lemma 1.2 there exist invertible operators X and Y on \mathfrak{S} with equal determinants such that $T = X - Y$. According to (4, Theorem III), there exist operators $P, A,$ and Q such that $X = PAQ$ and $Y = QAP$; the invertibility of X and Y guarantees that $P, A,$ and Q are invertible. Thus $T = PAQ - QAP$, as desired.

2. The separable case. We turn now to the case in which \mathfrak{S} is a separable, infinite-dimensional, space. According to (1, Theorem 3), an operator T on \mathfrak{S} is a commutator if it is not of the form $\lambda + K$ for some non-zero scalar λ and compact operator K . For such a commutator T there exist operators P_1 and Q_1 such that $T = P_1Q_1 - Q_1P_1$. Since, as noted before, for any scalar μ we also have that

$$T = (P_1 + \mu)(Q_1 + \mu) - (Q_1 + \mu)(P_1 + \mu),$$

μ_0 can be chosen so that the operators $P = P_1 + \mu_0$ and $Q = Q_1 + \mu_0$ are invertible. If we then define $A = 1$, we have that

$$T = PAQ - QAP$$

with invertible $P, A,$ and Q . Thus, it suffices to prove the theorem for operators T of the form $T = \lambda + K$, where $\lambda \neq 0$ and K is compact.

We shall have occasion to write T as a matrix with operator entries, and in so doing, we observe the usual conventions. If \mathfrak{S} is written as the direct sum

$$\mathfrak{S} = \mathfrak{S}_1 \oplus \dots \oplus \mathfrak{S}_m,$$

and if E_i denotes the projection of \mathfrak{S} onto \mathfrak{S}_i , then we write $T = (T_{ij})_{i,j=1}^m$, where T_{ij} denotes the linear operator

$$T_{ij} = E_iTE_j|_{\mathfrak{S}_j}.$$

The following lemma begins our program.

LEMMA 2.1. *Let $T \in \mathfrak{L}(\mathfrak{S})$ be of the form $T = \lambda + K$ for $\lambda \neq 0$ and K compact, and let $\epsilon > 0$. Then there exists a finite-dimensional subspace \mathfrak{R} of \mathfrak{S} such that if \mathfrak{L} and \mathfrak{M} are subspaces satisfying $\mathfrak{L} \subset \mathfrak{R}^\perp$ and $\mathfrak{M} \subset \mathfrak{L}^\perp$, and if E and F denote the projections of \mathfrak{S} onto \mathfrak{L} and \mathfrak{M} , respectively, then $ETF, FTE,$ and $ETE - \lambda E$ all have norm less than ϵ .*

Proof. It is well known that there exist finite-dimensional projections P with the property that

$$\|K - PKP\| < \epsilon.$$

(Indeed, if $\{P_n\}$ is any sequence of projections converging strongly to the identity operator, then $\|K - P_nKP_n\| \rightarrow 0$.) Fix any one such projection P_0 , choose its range for the subspace \mathfrak{R} , and denote by L the operator $K - P_0KP_0$. Then with $\mathfrak{L}, \mathfrak{M}, E,$ and F as in the statement of the lemma, we have that $EK = EL$ and $KE = LE$. Hence, $EKF, FKE,$ and EKE all have norm less than ϵ , and the result follows.

PROPOSITION 2.2. *Let $T \in \mathfrak{K}(\mathfrak{H})$ be of the form $T = \lambda + K$ for $\lambda \neq 0$ and K compact. Then there exists a decomposition $\mathfrak{H} = \mathfrak{H}_1 \oplus \mathfrak{H}_2$ of \mathfrak{H} into the direct sum of two infinite-dimensional subspaces such that, if the corresponding matrix for T is*

$$T = \begin{pmatrix} A_1 & B_1 \\ C_1 & D_1 \end{pmatrix},$$

then both A_1 and D_1 are invertible.

Proof. To begin with, it is a simple matter to obtain via Lemma 2.1 a preliminary resolution $\mathfrak{H} = \mathfrak{R}_1 \oplus \mathfrak{R}_2$ with respect to which the matrix representation

$$(I) \quad T = \begin{pmatrix} A & B \\ C & D \end{pmatrix}$$

has the property that D is invertible. Indeed, we have only to choose for ϵ any positive number less than $|\lambda|$, and then choose \mathfrak{R}_2 to be any infinite-dimensional subspace whose orthocomplement \mathfrak{R}_1 is infinite-dimensional and contains the subspace \mathfrak{R} of Lemma 2.1. Next, note that if U is a unitary operator on \mathfrak{H} with $(U_{ij})_{i,j=1}^2$ as its matrix representation relative to the decomposition $\mathfrak{H} = \mathfrak{R}_1 \oplus \mathfrak{R}_2$, and if

$$(II) \quad \begin{pmatrix} U_{11} & U_{12} \\ U_{21} & U_{22} \end{pmatrix} \begin{pmatrix} A & B \\ C & D \end{pmatrix} \begin{pmatrix} U_{11}^* & U_{21}^* \\ U_{12}^* & U_{22}^* \end{pmatrix} = \begin{pmatrix} A_1 & B_1 \\ C_1 & D_1 \end{pmatrix},$$

then

$$\begin{pmatrix} U^*A_1U & U^*B_1U \\ U^*C_1U & U^*D_1U \end{pmatrix}$$

is the matrix representation for T relative to the decomposition $\mathfrak{H} = \mathfrak{H}_1 \oplus \mathfrak{H}_2$, where $\mathfrak{H}_i = U^*(\mathfrak{R}_i)$, $i = 1, 2$. Thus the theorem will be proved if we can find a unitary operator U such that in equation (II), both A_1 and D_1 are invertible. Now the operator A is a compression of T , and therefore is also of the form $A = \lambda + K_1$, where K_1 is a compact operator on \mathfrak{R}_1 . Hence, A is either invertible or has a non-trivial null space. In the former case, the proof is complete; in the latter case, the set of all those vectors $x \in \mathfrak{R}_1$ satisfying $A^kx = 0$ for some positive integer k form a non-trivial finite-dimensional subspace \mathfrak{N}_1 of \mathfrak{R}_1 . Let $\dim(\mathfrak{N}_1) = n$, and define $\mathfrak{N}_2 = \mathfrak{R}_1 \ominus \mathfrak{N}_1$, so that $\mathfrak{R}_1 = \mathfrak{N}_1 \oplus \mathfrak{N}_2$. The subspace \mathfrak{N}_1 is invariant under A , and if we write N for the nilpotent operator in $\mathfrak{K}(\mathfrak{N}_1)$ defined by $N = A|_{\mathfrak{N}_1}$, then the matrix representation for A relative to the decomposition $\mathfrak{R}_1 = \mathfrak{N}_1 \oplus \mathfrak{N}_2$ has the form

$$A = \begin{pmatrix} N & A_{12} \\ 0 & A_{22} \end{pmatrix}.$$

The advantage of this particular dissection of A is that the diagonal entry A_{22} is invertible. To see this, note that A_{22} is of the form $\lambda + K_2$, where $K_2 \in \mathfrak{K}(\mathfrak{N}_2)$ is compact. Thus, it suffices to show that A_{22} has trivial null space.

Suppose, accordingly, that $A_{22}y = 0$ with $y \in \mathfrak{N}_2$. Then $Ay \in \mathfrak{N}_1$, so that $A^k(Ay) = A^{k+1}y = 0$ for some $k > 0$. But then $y \in \mathfrak{N}_1$, and therefore $y = 0$.

Let now \mathfrak{M}_1 be an n -dimensional subspace of \mathfrak{R}_2 , the precise determination of which will be made later, and write $\mathfrak{M}_2 = \mathfrak{R}_2 \ominus \mathfrak{M}_1$, so that

$$\mathfrak{S} = \mathfrak{N}_1 \oplus \mathfrak{N}_2 \oplus \mathfrak{M}_1 \oplus \mathfrak{M}_2.$$

The matrix representation for T corresponding to this decomposition may be written as

$$(III) \quad T = \begin{pmatrix} N & A_{12} & B_{11} & B_{12} \\ 0 & A_{22} & B_{21} & B_{22} \\ C_{11} & C_{12} & D_{11} & D_{12} \\ C_{21} & C_{22} & D_{21} & D_{22} \end{pmatrix}.$$

We next consider unitary operators $U(\theta)$ on \mathfrak{S} ($0 < \theta < \pi/2$) whose matrices relative to this same decomposition of \mathfrak{S} have the form

$$U(\theta) = \begin{pmatrix} \cos \theta & 0 & \sin \theta V & 0 \\ 0 & 1 & 0 & 0 \\ -\sin \theta V^* & 0 & \cos \theta & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix},$$

where V is some arbitrary isometry mapping \mathfrak{M}_1 onto \mathfrak{N}_1 . A brief calculation shows that in the representation of $U(\theta)TU^*(\theta)$ as a 2×2 matrix corresponding to the splitting $\mathfrak{S} = \mathfrak{R}_1 \oplus \mathfrak{R}_2$ (see (II) above), the entries A_1 and D_1 are given by

$$A_1(\theta) = \begin{pmatrix} \cos^2 \theta N + \sin^2 \theta VD_{11}V^* + \sin \theta \cos \theta (B_{11}V^* + VC_{11}) & \cos \theta A_{12} + \sin \theta VC_{12} \\ \sin \theta B_{21}V^* & A_{22} \end{pmatrix}$$

and

$$D_1(\theta) = \begin{pmatrix} \sin^2 \theta V^*NV + \cos^2 \theta D_{11} - \sin \theta \cos \theta (C_{11}V + V^*B_{11}) & -\sin \theta V^*B_{12} + \cos \theta D_{12} \\ -\sin \theta C_{21}V + \cos \theta D_{21} & D_{22} \end{pmatrix}.$$

Thus our task reduces to choosing the subspace \mathfrak{M}_1 and the angle θ in such a way that these operators are invertible. To this end, note that the entries of the matrix (III) are all bounded in norm by $\|T\|$, independently of how the subspace \mathfrak{M}_1 is selected. It follows that $\|D_1(\theta) - D\| \rightarrow 0$ as $\theta \rightarrow 0$, and that this convergence is uniform with respect to \mathfrak{M}_1 . Since D is invertible, there exist angles $\theta > 0$ so small that $D(\theta)$ is invertible no matter how \mathfrak{M}_1 is chosen. We choose one such θ_0 , hold it fixed, and proceed to adjust \mathfrak{M}_1 so as to make $A_1(\theta_0)$ invertible. That such a choice is possible may be seen most clearly as follows. Let $D_{11} = \lambda + K_3$, with λ and K_3 in $\mathfrak{L}(\mathfrak{M}_1)$. (The operator K_3 depends,

of course, on the choice of \mathfrak{M}_1 .) Also, write $A_1(\theta_0) = A_0 + \delta(\mathfrak{M}_1)$, where

$$A_0 = \begin{pmatrix} \cos^2\theta_0 N + \sin^2\theta_0 \lambda & \cos\theta_0 A_{12} \\ 0 & A_{22} \end{pmatrix}$$

and

$$\delta(\mathfrak{M}_1) = \begin{pmatrix} \sin^2\theta_0 V K_3 V^* + \sin\theta_0 \cos\theta_0 (B_{11} V^* + V C_{11}) & \sin\theta_0 V C_{12} \\ \sin\theta_0 B_{21} V^* & 0 \end{pmatrix}$$

so that A_0 is independent of the choice of \mathfrak{M}_1 . Since N is nilpotent and $\sin^2\theta_0 \lambda$ is a non-zero scalar, the entry $\cos^2\theta_0 N + \sin^2\theta_0 \lambda$ of A_0 is invertible; since A_{22} is also known to be invertible, it follows that A_0 is invertible. On the other hand, according to Lemma 2.1, it is possible to choose \mathfrak{M}_1 in such a way so as to make B_{11} , B_{21} , C_{11} , C_{12} , and K_3 as small in norm as desired. Since $\|V\| = 1$, it follows that by appropriate choice of \mathfrak{M}_1 , $\|\delta(\mathfrak{M}_1)\|$ can be made arbitrarily small. Hence $A_1(\theta_0)$ can be made arbitrarily close to A_0 , and the result follows.

Summary. We have shown that if T is an arbitrary operator of the form $\lambda + K$ with $\lambda \neq 0$ and K compact, then T can be viewed, relative to some decomposition $\mathfrak{S} = \mathfrak{S}_1 \oplus \mathfrak{S}_2$ of \mathfrak{S} , as a 2×2 matrix whose diagonal entries are invertible.

If we now identify \mathfrak{S}_2 with \mathfrak{S}_1 via a unitary isomorphism, then \mathfrak{S} is identified with $\mathfrak{S}_1 \oplus \mathfrak{S}_1$, and T is identified with (is unitarily equivalent to) an operator $T_1 \in \mathfrak{L}(\mathfrak{S}_1 \oplus \mathfrak{S}_1)$. The advantage of this identification is that T_1 can be regarded as a 2×2 matrix all of whose entries act on the same space \mathfrak{S}_1 ; of course, the diagonal entries of T_1 remain invertible. The following lemma thus concludes the proof of our theorem in the separable case.

LEMMA 2.3. *If T is an operator on $\mathfrak{L}(\mathfrak{S} \oplus \mathfrak{S})$ whose 2×2 matrix over $\mathfrak{L}(\mathfrak{S})$ is*

$$\begin{pmatrix} T_1 & T_2 \\ T_3 & T_4 \end{pmatrix},$$

where T_1 and T_4 are invertible operators, then there exist invertible operators P , A , and Q on $\mathfrak{S} \oplus \mathfrak{S}$ such that $T = PAQ - QAP$.

Proof. We define P , A , and Q by writing

$$P = \begin{pmatrix} -T_1 & 0 \\ 0 & T_4 \end{pmatrix}, \quad A = \begin{pmatrix} A_1 & 0 \\ 1 & A_2 \end{pmatrix}, \quad \text{and} \quad Q = \begin{pmatrix} Q_1 & 1 \\ 0 & Q_2 \end{pmatrix},$$

where the entries A_i , Q_i are to be determined. Note that if A_i and Q_i , $i = 1, 2$, are all invertible, then P , A , and Q are invertible also. A brief calculation reduces the matrix equation $PAQ - QAP = T$ to the system of equations

$$(IV) \quad \begin{cases} Q_1 A_1 T_1 - T_1 A_1 Q_1 = 0 \\ T_1 A_1 + A_2 T_4 = -T_2 \\ T_4 Q_1 + Q_2 T_1 = T_3 \\ T_4 A_2 Q_2 - Q_2 A_2 T_4 = 0. \end{cases}$$

That this system possesses invertible solutions A_1, A_2, Q_1, Q_2 when T_1 and T_4 are themselves both invertible may be seen as follows. If we agree to write

$$A_2 = \alpha Q_2^{-1} \quad \text{and} \quad Q_1 = \beta T_1$$

(where α and β denote positive parameters to be determined), then the first and last equations will be automatically satisfied, so that the problem reduces to solving the third equation

$$\beta T_4 T_1 + Q_2 T_1 = T_3$$

for Q_2 in such a way as to make Q_2 invertible, and then solving the second equation

$$T_1 A_1 + \alpha Q_2^{-1} T_4 = -T_2$$

for A_1 in such a way as to make it invertible. Clearly these requirements will be met if β is first chosen large enough to ensure the invertibility of $T_3 - \beta T_4 T_1$ and if α is then chosen large enough to make $T_2 + \alpha Q_2^{-1} T_4$ invertible.

3. The non-separable case. In this section we sketch a proof of the theorem in the case that $\dim(\mathfrak{S}) = \aleph > \aleph_0$. Let (K) denote the maximal proper norm-closed ideal in $\mathfrak{L}(\mathfrak{S})$. According to (1, Theorem 4), the non-commutators on \mathfrak{S} are exactly the operators of the form $\lambda + K$, where $\lambda \neq 0$ and $K \in (K)$. Furthermore, just as above, it suffices to treat the non-commutators. Let $T = \lambda + K$ be such an operator. Then the lemma obtained from Lemma 2.1 above by replacing the phrase “finite-dimensional subspace \mathfrak{R} ” by “subspace \mathfrak{R} of dimension less than \aleph ” is valid for T and is essentially contained in (1, Lemma 6.1) and (2, Lemma 4.1). Accordingly, in the notation of Lemma 2.1, let $\epsilon = |\lambda|/2$, let \mathfrak{R} be the corresponding subspace of dimension less than \aleph , and let \mathfrak{M} denote the smallest invariant subspace of T that contains \mathfrak{R} . An easy cardinality argument shows that \mathfrak{M} has dimension equal to that of \mathfrak{R} . Since \mathfrak{M}^\perp is orthogonal to \mathfrak{R} , the compression Z of $T - \lambda$ to \mathfrak{M}^\perp has norm less than $\epsilon = |\lambda|/2$, and it follows that the matrix for T relative to the decomposition $\mathfrak{S} = \mathfrak{M} \oplus \mathfrak{M}^\perp$ has the form

$$T = \begin{pmatrix} X & Y \\ 0 & Z + \lambda \end{pmatrix}.$$

Since $X \in \mathfrak{L}(\mathfrak{M})$ and $\dim \mathfrak{M} < \aleph$, we may assume by transfinite induction that the conclusion of the theorem holds for X . To see that the conclusion of the theorem also holds for $Z + \lambda$, write $\mathfrak{M}^\perp = \mathfrak{N}_1 \oplus \mathfrak{N}_2$, where $\dim \mathfrak{N}_1 = \dim \mathfrak{N}_2 = \dim \mathfrak{M}^\perp$. Then the matrix for $Z + \lambda$ relative to this resolution has the form

$$Z + \lambda = \begin{pmatrix} Z_1 + \lambda & Z_2 \\ Z_3 & Z_4 + \lambda \end{pmatrix},$$

and since $\|Z\| < \epsilon = |\lambda|/2$, $\|Z_1\|, \|Z_4\| < |\lambda|/2$, from which it follows that $Z_1 + \lambda$ and $Z_4 + \lambda$ are invertible. Thus Lemma 2.3, which is easily seen to be

independent of the dimension of \mathfrak{S} , can be applied to yield the desired conclusion for $Z + \lambda$.

The proof in the non-separable case is completed by the following lemma.

LEMMA 3.1. *Suppose that the conclusion of the theorem holds for operators X and Z on Hilbert spaces \mathfrak{S} and \mathfrak{R} , respectively, and let Y be any operator from \mathfrak{R} to \mathfrak{S} . Then the conclusion of the theorem also holds for the operator*

$$\begin{pmatrix} X & Y \\ 0 & Z \end{pmatrix}$$

on the space $\mathfrak{S} \oplus \mathfrak{R}$.

Proof. Choose invertible operators P_i, A_i , and Q_i ($i = 1, 2$) such that

$$P_1A_1Q_1 - Q_1A_1P_1 = X \quad \text{and} \quad P_2A_2Q_2 - Q_2A_2P_2 = Z.$$

Let $P \in \mathfrak{L}(\mathfrak{S} \oplus \mathfrak{R})$ denote the operator

$$P = \begin{pmatrix} P_1 & 0 \\ 0 & P_2 \end{pmatrix},$$

let

$$A_s = \begin{pmatrix} A_1 & 0 \\ 0 & sA_2 \end{pmatrix},$$

where s is a positive real parameter to be determined, and finally, let

$$Q_s(W) = \begin{pmatrix} Q_1 & W \\ 0 & s^{-1}Q_2 \end{pmatrix},$$

where W is an operator from \mathfrak{R} to \mathfrak{S} which also is to be determined. A simple calculation shows that

$$PA_sQ_s(W) - Q_s(W)A_sP = \begin{pmatrix} X & (P_1A_1)W - W(sA_2P_2) \\ 0 & Z \end{pmatrix}$$

so that, to complete the proof, it suffices to solve the equation

$$(V) \quad (P_1A_1)W - W(sA_2P_2) = Y$$

for s and W . Now for fixed s , it is well known that this equation possesses a unique solution W provided only that the spectra of P_1A_1 and sA_2P_2 are disjoint. Furthermore, since A_2P_2 is invertible, it is obviously possible to make these spectra disjoint by choosing s sufficiently large.

REMARK 3.2. The complete story concerning (V) is as follows: the spectrum of the linear transformation

$$W \rightarrow BW - WC$$

is precisely the set of differences $\beta - \gamma$, where β and γ run over the spectra of B and C , respectively. The usual proof of this fact (see 3) assumes that B, C , and W are all operators on the same Hilbert space, but the argument can easily

be modified so as to apply to the case in which B and C act on different Hilbert spaces and W is a linear transformation from one Hilbert space to the other.

REMARK 3.3. A very short construction due to Paul Federbush shows that every operator T on an infinite-dimensional space can be written as $T = PAQ - QAP$ for P, A, Q not invertible. The argument goes as follows. Write $\mathfrak{S} = \mathfrak{M} \oplus \mathfrak{M}^+$, where \mathfrak{M} and \mathfrak{M}^+ are of the same dimension, and let P (Q) be an isometry with range \mathfrak{M} (\mathfrak{M}^+). If X is an arbitrary operator, then $X = PAQ - QAP$, where $A = P^*XQ^* - Q^*XP^*$. We are also indebted to Federbush for bringing (4, Theorem 3) to our attention.

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