BOUNDDED POINT EVALUATIONS AND LOCAL SPECTRAL THEORY

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Abstract

We study in this paper the concept of bounded point evaluations for cyclic operators. We give a negative answer to a question of L.R. Williams Dynamic Systems and Applications 3(1994) 103-112. Furthermore, we generalize some results of Williams and give a simple proof of theorem 2.5 of [?] that non normal hyponormal weighted shifts have fat local spectra.
To my loving father and mother

To my brothers and sisters.
Contents
Introduction

The main purpose of our dissertation is to study links between local spectral theory and the concept of bounded point evaluations for arbitrary cyclic operators on a Hilbert spaces.

In chapter 1, we first fix some terminology and recall some basic notions concerning the spectral theory of a bounded operator on complex Banach spaces and vector valued analytic functions. The concepts of spectrum and resolvent set are introduced, the various subdivisions of the spectrum are studied and the analyticity of resolvent map on the resolvent set is proved. The last part of this chapter is devoted to weighted shift operators which are widely used for solving many problems in operator theory. Namely, we describe the spectrum and its parts of weighted shift operators.

In chapter 2 two basic concepts of local spectral theory are discussed. In section one are given basic properties of the operators with the single valued extension property which have been considered by N. Dunford in 1952 those bounded operators $T$ on Banach space $X$ such that for every open set $U \subset \mathbb{C}$, the equation $(T - \lambda I)F(\lambda) = 0$ admits the zero function $F \equiv 0$ as a unique analytic solution. However, the generalization of the spectrum and resolvent set are introduced and studied. Finally in section two are studied the elementary properties of operators with Dunford’s Condition C, which will be used frequently throughout this dissertation.

Chapter 3 is devoted to two classes of the most interesting operators which are defined around the notion of normal operators, namely the subnormal and hyponormal operators. The concept of subnormal and hyponormal operators was introduced by Paul R. Halmos in 1950 in [?], and many basic properties of subnormal operators were proved by Bram in [?]. In this chapter, we present the basic tools for their understanding. Are shown that the minimal normal extension of a subnormal operator $S$ on a Hilbert space $\mathcal{H}$ is unique up to an invertible isometry and its spectrum is contained in the spectrum of $S$ and the hyponormal operators have the single valued extension property and the Dunford’s Condition C. Lastly, are given characterizations of normal, subnormal and hyponormal weighted shift operators.

Chapter 4 is inspired by the papers of L.R. Williams [?] and T.T. Trent [?]. It contains a study of the bounded point evaluations of cyclic operators and our results. We give a negative answer to the question A posed by Williams in [?], we compare the concept of analytic bounded point evaluation of a weighted shift operators given by A.L. Shields in [?] with the one defined by L.R. Williams, we prove that if $T$ is a cyclic bounded operator on a Hilbert space $\mathcal{H}$ with Dunford’s Condition C and without point spectrum, then the local spectra of $x$ with respect to $T$ is equal to the spectrum of $T$ for each $x$ in a dense subset of $\mathcal{H}$, finally, we give a simple proof of the Theorem 2.5 of [?] using the bounded point evaluation of weighted shifts and the fact that a non zero analytic function has isolate zeroes.
Chapter 1

Spectral Theory for Weighted Shift Operators

1.1 Preliminaries

In this section we assemble some background material from spectral theory of operators which will be needed in the sequel.

In what follows, $\mathcal{X}^*$ will denote the dual space of a complex Banach space $\mathcal{X}$ and $\mathcal{L}(\mathcal{X})$ will denote the algebra of all linear bounded operators on $\mathcal{X}$. The spectrum $\sigma(T)$ of an operator $T \in \mathcal{L}(\mathcal{X})$ is defined as follows

$$\sigma(T) = \{ \lambda \in \mathbb{C} : T - \lambda I \text{ is not invertible in } \mathcal{L}(\mathcal{H}) \}.$$ 

It is a nonempty compact subset of $\mathbb{C}$ contained in the ball $\{ \lambda \in \mathbb{C} / |\lambda| \leq \|T\| \}$ (see [2]). The spectral radius $r(T)$ of an operator $T \in \mathcal{L}(\mathcal{X})$ is defined by

$$r(T) = \sup\{ |\lambda| : \lambda \in \sigma(T) \};$$

and satisfies

$$r(T) = \lim_{n \to \infty} \|T^n\|^{\frac{1}{n}} \leq \|T\|$$

in the sense that the indicated limit always exist and has the indicated value (see [2]). If an operator $T \in \mathcal{L}(\mathcal{X})$ is invertible then

$$\sigma(T^{-1}) = \left\{ \frac{1}{\lambda} : \lambda \in \sigma(T) \right\}$$

and so

$$\frac{1}{r(T^{-1})} = \inf \{ |\lambda| : \lambda \in \sigma(T) \}$$

hence, in this case the spectrum $\sigma(T)$ of $T$ lies in the annulus $\{ \lambda \in \mathbb{C} : \frac{1}{r(T^{-1})} \leq |\lambda| \leq r(T) \}$. It is often useful to ask of a point in the spectrum of an operator $T \in \mathcal{L}(\mathcal{X})$ how it got there. The question reduces therefore to this: why is a non-invertible operator not invertible? There are several possible ways of answering the question; they have led to several classifications of spectra. Perhaps the simplest approach to the subject is to recall that an operator $T \in \mathcal{L}(\mathcal{X})$ is bounded from below (i.e: $\inf_{\|x\|=1} \|Tx\| > 0$) and has a dense range if and only if it is invertible in $\mathcal{L}(\mathcal{X})$. Consequently,

$$\sigma(T) = \sigma_{ap}(T) \cup \Gamma(T)$$
where $\sigma_{ap}(T)$ is the approximate point spectrum of $T$, the set of complex numbers $\lambda$ such that $T - \lambda I$ is not bounded from below, and $\Gamma(T)$ is the compression spectrum of $T$, the set of complex numbers $\lambda$ such that the range of $T - \lambda I$ is not dense in $\mathcal{X}$.

The following results summarizes the main properties of the approximate point spectrum.

**Definition 1.1.1** The lower bound of an operator $T \in \mathcal{L}(\mathcal{X})$ is the value given by

$$m(T) = \inf_{\|x\|=1} \|Tx\|.$$  

**Proposition 1.1.2** Let $T$ be an operator in $\mathcal{L}(\mathcal{X})$. If $\lambda \in \mathbb{C}$ such that $|\lambda| < m(T)$ then $T - \lambda I$ is bounded below.

**Proof.** It suffice to observe that

$$m(T) - |\lambda| \leq \|Tx\| - |\lambda| \leq \|(T - \lambda I)x\| \text{ for every } x \in \mathcal{X}, \|x\| = 1.$$  

**Remark 1.1.3** Let $T$ and $S$ be two operators in $\mathcal{L}(\mathcal{X})$. Then:

(i) A complex number $\lambda$ is in $\sigma_{ap}(T)$ if and only if $m(T - \lambda I) = 0$.

(ii) $m(T)\|x\| \leq \|Tx\|$ for every $x \in \mathcal{X}$.

(iii) If $T$ is invertible then $m(T) = \frac{1}{\|T^{-1}\|}$.

(iv) $m(T)m(S) \leq m(TS) \leq \|T\|m(S)$.

(v) $m(T) > 0$ if and only if $T$ is one to one and has closed range.

**Proposition 1.1.4** Let $T \in \mathcal{L}(\mathcal{X})$. Then the sequence $\left( \left( m(T^n) \right)^{\frac{1}{n}} \right)_n$ is convergent and has a limit equal to $\sup \left( m(T^n) \right)^{\frac{1}{n}}$. This limit will be denoted by $r_1(T)$.

**Proof.** It is clear that

$$\sup_{n \geq 1} \left( m(T^n) \right)^{\frac{1}{n}} \geq \limsup_{n \to \infty} \left( m(T^n) \right)^{\frac{1}{n}}.$$  

Now, fix $k$. Then for every $n$ there exists $p = p(n)$ and $r = r(n)$ such that $0 \leq r < k$ and $n = kp + r$. From Remark 1.1.3 it follows that

$$m(T^n) \geq m(T^k)^p m(T)^r.$$  

Hence

$$\liminf_n \left( m(T^n) \right)^{\frac{1}{n}} \geq \left( m(T^k) \right)^{\frac{1}{k}}.$$  

Since $k$ is arbitrary then

$$\liminf_n \left( m(T^n) \right)^{\frac{1}{n}} \geq \sup_n \left( m(T^n) \right)^{\frac{1}{n}}$$  

and the result follows from (*) and (**).
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Remark 1.1.5 The same argument of this proof can be used to prove that
\[ r(T) = \lim_{n \to \infty} \|T^n\|^{\frac{1}{n}} \]
for every operator \( T \) in \( \mathcal{L}(\mathcal{X}) \).

Proposition 1.1.6 For every operator \( T \) in \( \mathcal{L}(\mathcal{X}) \),
\[ \sigma_{\text{ap}}(T) \subset \{ \lambda \in \mathbb{C} : r_1(T) \leq |\lambda| \leq r(T) \} \]

Proof. Since \( \sigma_{\text{ap}}(T) \subset \sigma(T) \subset \{ \lambda \in \mathbb{C} : |\lambda| \leq r(T) \} \) then it suffice to prove that \( r_1(T) \leq |\lambda| \) for every \( \lambda \in \sigma_{\text{ap}}(T) \). So, let \( \lambda \in \mathbb{C} \) such that \( |\lambda| < r_1(T) \). Then there \( |\lambda|^n < m(T^n) \) for some integer \( n \). By Proposition 1.1.2 we see that
\[ T^n - \lambda^n I = (T^{n-1} + \lambda T^{n-2} + \ldots + \lambda^{n-1} I)(T - \lambda I) \]
is bounded below. Hence, it follows from the inequality (iv) in Remark 1.1.3 that \( (T - \lambda I) \) is bounded below; that is \( \lambda \notin \sigma_{\text{ap}}(T) \).

Theorem 1.1.7 Let \( T \in \mathcal{L}(\mathcal{X}) \).
(i) A complex number \( \lambda \) belongs to \( \sigma_{\text{ap}}(T) \) if and only if there exists a sequence \((x_n)\) of unit vectors in \( \mathcal{X} \) such that \( \lim_{n \to \infty} \| (T - \lambda I)x_n \| = 0 \). In particular, the approximate point spectrum of \( T \) contains the set of eigenvalues of \( T \), which is called the point spectrum of \( T \) and denoted by \( \sigma_p(T) \).
(ii) The approximate point spectrum of \( T \), \( \sigma_{\text{ap}}(T) \), is a closed non-empty subset of the spectrum of \( T \).
(iii) The boundary of the spectrum of \( T \) is contained in \( \sigma_{\text{ap}}(T) \).

Proof. The first property holds immediately from the definition of the approximate point spectrum of \( T \). Now we will show first that \( \sigma_{\text{ap}}(T) \subset \sigma(T) \). Suppose \( \lambda \notin \sigma(T) \), then \( (T - \lambda I) \) has bounded inverse and so
\[ M \leq \| (T - \lambda I)x \| \quad \text{for every } x \in \mathcal{X} \quad \|x\| = 1 \]
where \( M = \left\| \frac{1}{(T - \lambda I)^{-1}} \right\| \), so \( \lambda \notin \sigma_{\text{ap}}(T) \). Hence \( \sigma_{\text{ap}}(T) \subset \sigma(T) \).

Let \( \lambda \in \mathbb{C} \setminus \sigma_{\text{ap}}(T) \) then there exists \( M > 0 \) such that
\[ M \leq \| (T - \lambda I)x \| \quad \text{for every } x \in \mathcal{X} \quad \|x\| = 1 \]
Consequently, for every \( \mu \in \mathbb{C} \) such that \( |\lambda - \mu| < \frac{M}{2} \) we have
\[ \frac{M}{2} \leq \| (T - \lambda I)x \| - |\lambda - \mu| \leq \| (T - \mu I)x \| \quad \text{for every } x \in \mathcal{X} \quad \|x\| = 1, \]
so
\[ \{ \mu \in \mathbb{C} / \ |\lambda - \mu| < \frac{M}{2} \} \subset \mathbb{C} \setminus \sigma_{\text{ap}}(T). \]
Hence \( \sigma_{\text{ap}}(T) \) is closed.

Next, let \( \lambda \) be a point on the boundary of \( \sigma(T) \) and let \( \epsilon > 0 \). Since \( \sigma(T) \) and \( \mathbb{C} \setminus \sigma(T) \) have the same boundary then there is \( \mu \notin \sigma(T) \) such that \( |\lambda - \mu| < \frac{\epsilon}{2} \). We have
\[ \frac{\epsilon}{2} \leq \frac{1}{d(\mu)} \leq \| (T - \mu I)^{-1} \| \quad (\text{see next section Theorem 1.2.1}) \]
where \(d(\mu)\) denotes the distance from \(\mu\) to the spectrum of \(T\). Therefore there is \(x \in \mathcal{X}\) such that \(\|x\| = 1\) and
\[
\frac{1}{\epsilon} \leq \|(T - \mu I)^{-1}x\|.
\]
Let \(y = \frac{1}{\|(T - \mu I)^{-1}x\|}(T - \mu I)^{-1}x\) then \(\|y\| = 1\) and
\[
\|(T - \lambda I)y - (T - \mu I)y\| < \frac{\epsilon}{2}.
\]
It follows that
\[
\|(T - \lambda I)y\| \leq \|(T - \lambda I)y - (T - \mu I)y\| + \|(T - \mu I)y\|
< \frac{\epsilon}{2} + \frac{\|x\|}{\|(T - \mu I)^{-1}\|}
< \frac{3\epsilon}{2}
\]
Since \(\epsilon > 0\) is arbitrary it follows that \(\lambda \in \sigma_{ap}(T)\). Since a nonempty compact subset of \(\mathbb{C}\) has a nonempty boundary then \(\sigma_{ap}(T)\) is nonempty. The proof is complete. 

\textbf{Theorem 1.1.8} Let \(T^* \in \mathcal{L}(\mathcal{X}^*)\) be the adjoint operator of an bounded operator \(T \in \mathcal{L}(\mathcal{X})\). Then the following are equivalent:

(i) \(T\) is bounded below.

(ii) \(T^*\) is surjective.

For the proof, see [?\] Theorem 57.18 page 245.

\section{1.2 Analyticity of resolvents}

We first recall some basic notions and results from analyticity of functions whose domains are open sets in the complex plane and whose values are vectors of the Banach space \(\mathcal{X}\), the monographs [?], [?], and [?] contain further informations.

Let \(\Omega\) be a nonempty open subset of \(\mathbb{C}\). A function \(\rho\) from \(\Omega\) to \(\mathcal{X}\) is said to be a analytic in \(\Omega\) if for every \(\phi \in \mathcal{X}^*\), the complex value function \(\phi \circ \rho\) is analytic in the usual sense. In the case, when \(\mathcal{X}\) is a Hilbert space \(\mathcal{H}\) this definition is equivalent to that for every \(x \in \mathcal{H}\), the function
\[
\Omega \rightarrow \mathbb{C}
\lambda \mapsto \langle \rho(\lambda) , x \rangle
\]
is analytic in the usual sense. If \(\mathcal{X}\) is the algebra \(\mathcal{L}(\mathcal{H})\) of all bounded operators on a Hilbert space \(\mathcal{H}\), the analyticity of \(\rho\) can be interpreted as following: For every \(x, y \in \mathcal{H}\) the function
\[
\Omega \rightarrow \mathbb{C}
\lambda \mapsto \langle \rho(\lambda)x , y \rangle
\]
is analytic in the usual sense.

Using the fact that 0 is the unique element \(x \in \mathcal{X}\) satisfying \(\phi(x) = 0\) for every \(\phi \in \mathcal{X}^*\), we note that almost all the results concerning the analytic complex functions can be extended to vector valued functions. Such an useful result is Liouville's Theorem which say that every bounded entire vector valued function is a constant.

Let \(T \in \mathcal{L}(\mathcal{X})\) be a bounded operator on \(\mathcal{X}\). The complement set of the spectrum \(\sigma(T)\) of
1.2. ANALYTICITY OF RESOLVENTS

$T$, denoted by $\rho(T)$, is called the resolvent set of $T$. It is a nonempty open subset of $\mathbb{C}$. The mapping

$$R_T : \rho(T) \rightarrow \mathcal{L}(\mathcal{X})$$

$$\lambda \mapsto (T - \lambda I)^{-1}$$

is called the resolvent map of $T$.

**Theorem 1.2.1** Let $T \in \mathcal{L}(\mathcal{X})$. The resolvent map of $T$ is analytic on $\rho(T)$ vanishing at infinity and has the following property

$$\frac{1}{d(\lambda)} \leq \|R_T(\lambda)\| \text{ for every } \lambda \in \rho(T)$$

where $d(\lambda)$ is the distance from $\lambda$ to the spectrum $\sigma(T)$ of $T$. Therefore $\|R_T(\lambda)\| \rightarrow \infty$ as $d(\lambda) \rightarrow 0$.

In the proof of this Theorem we shall require the following Lemma:

**Lemma 1.2.2** Let $T \in \mathcal{L}(\mathcal{X})$ be a bounded operator on $\mathcal{X}$. If $\|T\| < 1$ then $I + T$ is invertible and

$$\|(I - T)^{-1} + I + T\| \leq \frac{\|T\|^2}{1 - \|T\|}.$$

In particular, if $\lambda$ is a fixed point in $\rho(T)$ then for every $\mu \in \mathbb{C}$ with $|\mu| < \|(T - \lambda I)^{-1}\|^{-1}$ we have $\lambda + \mu \in \rho(T)$.

**Proof.** Since $\|T\| < 1$ and $\|T^k\| \leq \|T\|^k$ for every $k \geq 0$ then this series $\sum_{n=0}^{+\infty} T^n$ converges in $\mathcal{L}(\mathcal{X})$ and,

$$(T - I) \left( \sum_{n=0}^{+\infty} T^n \right) = \left( \sum_{n=0}^{+\infty} T^n \right) (T - I) = I.$$

Hence, $T - I$ is invertible and

$$(T - I)^{-1} = \sum_{n=0}^{+\infty} T^n.$$ 

The rest of this proof follows easily.

**Proof of Theorem 1.2.1.** Let $\lambda$ be a fixed point in $\rho(T)$. Then there is $0 < \rho$ such that $\mu \in \rho(T)$ for every $\mu \in \mathbb{C}$ with $|\lambda - \mu| < \rho$. Since

$$\frac{(T - \lambda I)^{-1} - (T - \mu I)^{-1}}{\lambda - \mu} = (T - \lambda I)^{-1}(T - \mu I)^{-1},$$

then

$$\lim_{\mu \to \lambda} \frac{(T - \mu I)^{-1} - (T - \lambda I)^{-1}}{\mu - \lambda} = (T - \lambda I)^{-2}.$$

Therefore, for every $\phi \in \mathcal{L}(\mathcal{X})^*$

$$\lim_{\mu \to \lambda} \frac{\phi \circ R_T(\mu) - \phi \circ R_T(\lambda)}{\mu - \lambda} = \phi \circ (T - \lambda I)^{-2}.$$

Hence, $\phi \circ R_T$ is differentiable on $\rho(T)$. So, $R_T$ is analytic on $\rho(T)$. On the other hand, we have

$$\lambda R_T(\lambda) = \left( \frac{1}{\lambda} (T - I) \right)^{-1},$$

this shows that the resolvent $R_T$ of $T$ vanishing at infinity.

Since $\sigma(T)$ is a nonempty compact set then there is $\mu_0 \in \sigma(T)$ such that $d(\lambda) = |\mu_0 - \lambda|$. So, if $d(\lambda) = |\lambda_0 - \lambda| < \|R_T(\lambda_0)^{-1} = \|T - \lambda I \|^2$ then it follows from Lemma 1.2.2 that $\mu_0 = (\mu_0 - \lambda) + \lambda \in \rho(T)$ . Contradiction. This completes the proof of the Theorem. $\blacksquare$
1.3 Weighted Shift Operators

A weighted shift operator $T$ on complex Hilbert space $\mathcal{H}$ is an operator that maps each vector in some orthonormal basis $(e_n)_n$ into a scalar multiple of the next vector

$$T e_n = \omega_n e_{n+1}$$

for all $n$. The operator $T$ is called a unilateral weighted shift when the index $n$ runs the nonnegative integers and it is called a bilateral weighted shift when the $n$ runs over all integers. In what follows, $T$ will always denote a weighted shift operator with a weight sequence $(\omega_n)_n$. We shall sometimes omit the adjective ”weighted” and refer to $T$ simply as a shift.

An operator $A$ on a n-dimensional Hilbert space is called a finite-dimensional weighted shift if there are numbers $\{\alpha_1, ..., \alpha_n\}$ and an orthonormal basis $v_1, ..., v_n$ such that

$$Av_k = \alpha_k v_{k+1} \quad (k < n)$$
$$Av_n = 0.$$ Such operator is nilpotent (i.e: there is a positive integer $s$ such that $A^s = 0$).

Note that a weighted shift operator is injective if and only if none of the weights is zero. If finitely many weights are zero then $T$ is a direct sum of a finite-dimensional weighted shifts and an infinite-dimensional injective weighted shift. If infinitely many weights are zero then $T$ is the direct sum of an infinite family of finite-dimensional weighted shifts. This situation can be used to give an example of an operator with spectral radius 1, which is the direct sum of countable family of finite dimensional nilpotent operators each of which has spectrum $\{0\}$ (see [?]).

From now we shall assume that none of the weights is zero and let $\beta$ be the following sequence given by:

$$\beta_n = \begin{cases} 
\omega_0 \cdots \omega_{n-1} & \text{if } n > 0 \\
1 & \text{if } n = 0 \\
\frac{1}{\omega_n \cdots \omega_{-1}} & \text{if } n < 0
\end{cases}$$

An operator $U$ in $L(\mathcal{H})$ is called unitary if $UU^* = U^*U = I$ which is equivalent that $U$ is invertible and $U^{-1} = U^*$. Two bounded operators $A$ and $B$ on $\mathcal{H}$ are said to be unitarily equivalent if there is unitary operator $U$ in $L(\mathcal{H})$ such that $AU = UB$. Therefore, two unitarily equivalent bounded operators on $\mathcal{H}$ have the same spectrum, the same point spectrum, the same approximate point spectrum and the same compression spectrum.

1.3.1 Elementary properties

**Proposition 1.3.1** The shift $T$ is bounded if and only if the weight sequence is bounded. In this case,

$$\|T^n\| = \sup_k \|w_k w_{k+1} \cdots w_{n+k-1}\| = \sup_k \left|\frac{\beta_{n+k}}{\beta_k}\right|.$$ 

**Proof.** First suppose that the shift $T$ is bounded then for every integer $n$ we have

$$\|Te_n\| = |\omega_n| \leq \|T\|$$
so the weight \((ω_n)_n\) is obviously bounded. Conversely, suppose there is \(M > 0\) such that \(|ω_n| ≤ M\) for every \(n\) then for every \(x = \sum_n a_n e_n \in \mathcal{H}\) we have \(\|Tx\|^2 = \sum_n |a_n ω_n|^2\) therefore

\[
\|Tx\| ≤ M \left( \sum_n |a_n|^2 \right)^{\frac{1}{2}} = M \|x\| \text{ for every } x ∈ \mathcal{H}.
\]

Thus \(T\) is bounded. The equalities follow from the following relation

\[
T^n e_k = ω_k ω_{k+1} \cdots ω_{n+k-1} e_{n+k}.
\]

In the sequel we suppose that the weights \((ω_n)_n\) is bounded just to have that \(T\) is in \(L(\mathcal{H})\).

**Proposition 1.3.2** If \(T\) is bilateral weighted shift then

\[
T^* e_n = \overline{ω_{n-1}} e_{n-1} \text{ for every } n.
\]

If \(T\) is unilateral weighted shift then

\[
T^* e_n = \overline{ω_{n-1}} e_{n-1} \text{ for every } n > 1
\]

\[
T^* e_0 = 0
\]

**Proof.** For every integers \(n\) and \(k\) we have

\[
\langle T^* e_n , e_k \rangle = \langle e_n , Te_k \rangle = \langle e_n , ω_k e_{k+1} \rangle = \overline{ω_k}(e_n , e_{k+1})
\]

and the result follows. \(\blacksquare\)

**Remark 1.3.3** It follows from Proposition 1.3.2 that the unilateral shift \(T\) is never invertible because \(T^*\) is not, but the bilateral shift \(T\) can be invertible. It is invertible if and only if the weight \((\frac{1}{ω_n})_n\) is bounded if and only if \(\inf \frac{β_{n+1}}{β_n} > 0\). In this case, for \(n = 0, 1, 2, \ldots\)

\[
\|T^{-n}\| = \sup_k \frac{1}{|w_k w_{k+1} \cdots w_{n+k-1}|} = \sup_k \left| \frac{β_k}{β_{n+k}} \right| = \left[ \inf_k \frac{β_{n+k}}{β_k} \right]^{-1}.
\]

For this, it suffices to observe that \(T^{-1} v_k = \frac{1}{ω_k} v_{k+1}\) where \(v_k = e_{-k}\). Thus \(T^{-1}\) is a bilateral weighted shift and apply Proposition 1.3.1.

**Proposition 1.3.4** If \((λ_n)\) are complex numbers of modulus 1, then \(T\) is unitarily equivalent to the weighted shift operator with weight sequence \((λ_n + 1) ω_n)_n\).

**Proof.** Let \(U\) be the unitarily operator defined by \(U e_n = λ_n e_n\). Then the operator \(U^* TU\) is a weighted shift with the weight sequence indicated above. \(\blacksquare\)

**Corollary 1.3.5** \(T\) is unitarily equivalent to the weighted shift operator with weight sequence \((|ω_n|)_n\).
Proof. Choose \( \lambda_0 = 1 \); the choice of the remaining \( \lambda_n \) is then forced in both the unilateral and the bilateral cases.

Corollary 1.3.6

If \(|c| = 1\), then \( T \) and \( cT \) are unitarily equivalent.

Proof. Take \( \lambda_n = c^n \) for every \( n \) and apply Proposition 1.3.5.

Remark 1.3.7

In this section, our goal is to describe the spectrum, the point spectrum and the approximate point spectrum of the weighted shift operators. We have already mentioned that two unitarily equivalent bounded operators on \( \mathcal{H} \) have the same spectrum, the same point spectrum and the same approximate point spectrum. Therefore by the Corollary 1.3.5, we may assume that the weights are nonnegative and by the Corollary 1.3.6 we see that the spectrum, the point spectrum and the approximate point spectrum of \( T \) have circular symmetry about the origin. In particular, the entire circle \( \{ \lambda \in \mathbb{C} \mid |\lambda| = r(T) \} \) is in the spectrum \( \sigma(T) \) of \( T \). When \( T \) is bilateral invertible shift then by the same reason as before, the entire circle \( \{ \lambda \in \mathbb{C} \mid |\lambda| = \frac{1}{r(T-1)} \} \) is also in the spectrum \( \sigma(T) \) of \( T \).

1.3.2 Spectrum of Weighted Shifts

Theorem 1.3.8

If \( T \) is unilateral shift then its spectrum is the disc

\[
\{ \lambda \in \mathbb{C} \mid |\lambda| \leq r(T) \}.
\]

Proof.

It is known that

\[
\sigma(T) \subset \{ \lambda \in \mathbb{C} \mid |\lambda| \leq r(T) \}.
\]

Conversely, let \( \lambda \) be in the resolvent set \( \rho(T) \) of \( T \) and set \( x = (T - \lambda I)^{-1} e_0 = \sum_n a_n e_n \), then

\[
a_n = -\frac{1}{\lambda^{n+1} \beta_n} \text{ for every } n.
\]

On the other hand, we have

\[
\langle (T - \lambda I)^{-1} e_n, e_{n+k} \rangle = \frac{1}{\beta_n} \langle (T - \lambda I)^{-1} T^n e_0, e_{n+k} \rangle
\]

\[
= \frac{1}{\beta_n} \langle T^n (T - \lambda I)^{-1} e_0, e_{n+k} \rangle
\]

\[
= \frac{1}{\beta_n} \langle T^n x, e_{n+k} \rangle
\]

\[
= \frac{\beta_{n+k}}{\beta_n} a_k
\]

\[
= -\frac{\beta_{n+k}}{\beta_n} \frac{1}{\lambda^{k+1}}.
\]

Using Cauchy-Schwartz inequality we get,

\[
\frac{\beta_{n+k}}{\beta_n} \frac{1}{|\lambda|^{k+1}} \leq \|(T - \lambda I)^{-1}\|.
\]

By passing to the supremum on \( n \) we get

\[
\|T^k\| \leq |\lambda|^{k+1} \|(T - \lambda I)^{-1}\|.
\]

Taking \( k \)th roots and letting \( k \to \infty \) gives \( r(T) \leq |\lambda| \). Strict inequality must hold since the entire circle \( \{ \lambda \in \mathbb{C} \mid |\lambda| = r(T) \} \) is in the spectrum of \( T \) (see Remark 1.3.7). This proves the theorem.
1.3. WEIGHTED SHIFT OPERATORS

**Theorem 1.3.9** (i) *If* \( T \) *is a invertible bilateral weighted shift then its spectrum is the annulus*

\[ \{ \lambda \in \mathbb{C} : \frac{1}{r(T-1)} \leq |\lambda| \leq r(T) \}. \]

(ii) *If* \( T \) *is a bilateral weighted shift that is not invertible then its spectrum is the disc*

\[ \{ \lambda \in \mathbb{C} : |\lambda| \leq r(T) \}. \]

**Proof.** Let \( \lambda \neq 0 \) lie in the resolvent set \( \rho(T) \) of \( T \) then for \( x = \sum_{n \in \mathbb{Z}} a_ne_n = ((T-\lambda I)^{-1})e_0 \) we have

\[
\begin{aligned}
a_{n+1} - \lambda a_0 &= 1 \\
a_n\omega_n - \lambda a_{n+1} &= 0 \quad \text{for every } n \neq -1
\end{aligned}
\]

and so,

\[
a_n = \frac{1}{\lambda^n} a_0 \beta_n \quad \text{and} \quad a_{-n-1} = \lambda^n \beta_{-n-1} \omega_{-1} a_{-1} \quad \text{for every } n > 0.
\]

By the same calculation as in the proof of Theorem 1.3.8 we have

\[
|a_k| \leq \frac{\beta_n \beta_k}{\beta_{n+k}} \| (T - \lambda I)^{-1} \| \quad \text{for every } n \in \mathbb{N} \text{ and } k \in \mathbb{Z}.
\]

By (*) either \( a_0 \) or \( a_1 \) is non zero. Consider the following two cases.

**Case 1:** \( a_0 \neq 0 \). If we multiply \((***)\) by \( |\lambda|^k \) and apply \((**)\) we obtain

\[
|a_0| \frac{\beta_{n+k}}{\beta_n} \leq |\lambda|^k \| (T - \lambda I)^{-1} \| \quad \text{for every } k \in \mathbb{N}.
\]

Hence by the same argument of the proof of Theorem 1.3.8 we get \( r(T) < |\lambda| \).

**Case 2:** \( a_1 \neq 0 \). In \((***)\) we take \( k = -m \) for \( m \geq 1 \) and by applying the second equation in \((**)\) we get

\[
|\lambda|^{n+1} \frac{\beta_{n-m}}{\beta_n} \leq \frac{1}{\omega_{-1} a_{-1}} \| (T - \lambda I)^{-1} \|.
\]

In this case \( T \) is invertible (just take in the last inequality \( m = 1 \)). By passing to the infimum on \( n \) we get

\[
|\lambda|^{m-1} \| T^{-m} \| \leq \frac{1}{\omega_{-1} a_{-1}} \| (T - \lambda I)^{-1} \|.
\]

Taking \( m \)th roots and letting \( m \to \infty \) we obtain \( |\lambda| \leq \frac{1}{r(T-1)} \). The equality is excluded by circular symmetry about the origin. So, \( |\lambda| < \frac{1}{r(T-1)} \).

Since \( a_{-1}\omega_{-1} - \lambda a_0 = 1 \), then at least one of these two cases must hold for each given \( \lambda \). On the other hand they cannot both occur since the conclusions exclude one another. If \( T \) is invertible then the two cases together yield the desired conclusion. If \( T \) is not invertible then \( \inf_a \frac{\beta_{a+1}}{\beta_a} = 0 \). So, \( a_0 = \infty \) in the inequality \((***)\) and by passing to the infimum on \( n \) the second case will be excluded; hence the first case yields the desired conclusion. \( \blacksquare \)

**Theorem 1.3.10** *If* \( T \) *is unilateral weighted shift then*

\[ \sigma_{ap}(T) = \{ \lambda \in \mathbb{C} : r_1(T) \leq |\lambda| \leq r(T) \}. \]
Proof. It is known that \( \sigma_{ap}(T) \subset \{ \lambda \in \mathbb{C} : r_1(T) \leq |\lambda| \leq r(T) \} \) (see Proposition 1.1.6), so if \( r_1(T) = r(T) \) then by nonemptiness and circular symmetry of the approximate point spectrum of \( T \) we have \( \sigma_{ap}(T) = \{ \lambda \in \mathbb{C} : |\lambda| = r(T) \} \). Now suppose that \( r_1(T) < r(T) \). Since \( \sigma_{ap}(T) \) is closed and has circular symmetry it suffice to prove that every positive real \( c, r_1(T) < c < r(T) \) is lies in \( \sigma_{ap}(T) \). Let \( c \in (r_1(T), r(T)) \) and choose two reals numbers \( a, b \) such that \( r_1(T) < a < c < b < r(T) \). Let \( c > 0 \), since \( r(T) = \lim_{n \to \infty} \sup_{k} \beta_{m+k} \frac{1}{\beta_{m}} \) and \( \frac{a}{b} < 1 \) then there exists \( n \) and \( k \) such that

\[
\left[ \frac{c}{b} \right]^n < \epsilon \quad \text{and} \quad \left[ \frac{\beta_{m+k}}{\beta_{m}} \right]^{\frac{1}{n}} > b.
\]

Also, we have \( r_1(T) = \lim_{m \to \infty} \left[ \inf_{p} \beta_{m+p} \right]^{\frac{1}{m}} \) and \( \frac{a}{c} < 1 \) then there exist \( p \) and \( m > n + k \) such that

\[
\left[ \frac{a}{c} \right]^p < \epsilon \quad \text{and} \quad \left[ \frac{\beta_{m+p}}{\beta_{m}} \right]^{\frac{1}{p}} < a.
\]

Let \( x = \sum_{s} a_s e_s \in \mathcal{H} \) where

\[
\begin{align*}
    a_k &= 1, \\
    a_s &= \frac{\beta_s}{\beta_k c^{s-k}} \quad \text{if} \quad k + 1 \leq s \leq m + p + 1, \\
    a_s &= 0 \quad \text{otherwise}.
\end{align*}
\]

And so,

\[
a_{s-1} \omega_{s-1} - ca_s = 0 \quad \text{for every} \quad k < s < m + p.
\]

\[
\left( * \right) \quad \frac{1}{a_{n+k}} = \frac{\beta_{s} \epsilon^{n}}{\beta_{n+k}} < \left[ \frac{c}{b} \right]^n < \epsilon,
\]

\[
\left( ** \right) \quad \frac{a_{m+k}}{a_m} = \frac{\beta_{m+p} \epsilon^p}{\beta_{m} c^p} < \left[ \frac{a}{c} \right]^p < \epsilon,
\]

Therefore from \( * \) it follows

\[
Tx - cx = a_{m+p} \omega_{m+p} e_{m+p+1} - ca_k e_k
\]

and hence

\[
\|Tx - cx\|^2 = |a_{m+p} \omega_{m+p}|^2 + c^2 \leq \|T\|^2 (|a_{m+p}|^2 + 1)
\]

because \( c < r(T) \leq \|T\| = \sup_{r} \omega_{r} \). On the other hand, we have

\[
\|x\|^2 = \sum_{r} |a_r|^2 \geq |a_m|^2 + |a_{n+k}|^2
\]

then

\[
\frac{\|Tx - cx\|^2}{\|x\|^2} \leq \frac{\|T\|^2 (|a_{m+p}|^2 + 1)}{|a_m|^2 + |a_{n+k}|^2} \leq \|T\|^2 \max \left( \frac{1}{|a_{n+k}|^2}, \frac{|a_{m+p}|^2}{|a_m|^2} \right) < \epsilon^2 \|T\|^2 \quad \text{by \( ** \)} \quad \text{and} \quad \left( ** \right)
\]
and so, \( c \) is in \( \sigma_{ap}(T) \).

In dealing with bilateral shifts we shall use the following notations.

\[
\begin{align*}
 r_{1}^+ &= \lim_{n \to +\infty} \left[ \inf_{j \geq 0} \frac{\beta_{n+j}}{\beta_{j}} \right]^{\frac{1}{n}}, \quad r_{1}^- &= \lim_{n \to +\infty} \left[ \inf_{j \geq 0} \frac{\beta_{j}}{\beta_{n+j}} \right]^{\frac{1}{n}} \\
 r_{1}^+ &= \lim_{n \to +\infty} \left[ \sup_{j \leq 0} \frac{\beta_{n+j}}{\beta_{j}} \right]^{\frac{1}{n}}, \quad r_{1}^- &= \lim_{n \to +\infty} \left[ \sup_{j \leq 0} \frac{\beta_{j}}{\beta_{n+j}} \right]^{\frac{1}{n}}
\end{align*}
\]

The proof of the following theorem require a lot of technical computations, we omit it here.

**Theorem 1.3.11 (RIDGE [?]).** If \( T \) is bilateral shift and if \( r^{-} < r_{1}^{+} \), then

\[
\sigma_{ap}(T) = \{ \lambda \in \mathbb{C} : r_{1}^{-} \leq |\lambda| \leq r_{1}^{+} \} \cup \{ \lambda \in \mathbb{C} : r_{1}^{+} \leq |\lambda| \leq r^{+} \}.
\]

Otherwise

\[
\sigma_{ap}(T) = \sigma(T) = \{ \lambda \in \mathbb{C} : \min(r_{1}^{-}, r_{1}^{+}) \leq |\lambda| \leq \max(r^{-}, r^{+}) \}.
\]

**Theorem 1.3.12** If \( T \) is unilateral shift then:

(i) \( \sigma_{p}(T) \) is empty.

(ii) \( \sigma_{p}(T^{*}) = \{0\} \) if \( r_{2}(T) = 0 \) otherwise

\[
\{ \lambda \in \mathbb{C} : |\lambda| < r_{2}(T) \} \subset \sigma_{p}(T^{*}) \subset \{ \lambda \in \mathbb{C} : |\lambda| \leq r_{2}(T) \}
\]

where \( r_{2}(T) = \lim_{n \to +\infty} \inf(\beta_{n})^{\frac{1}{n}} \). Furthermore, all eigenvalues of \( T^{*} \) are simples.

**Proof.** (i) let \( \lambda \in \mathbb{C} \) and let \( x = \sum_{n \in \mathbb{N}} a_{n}e_{n} \in \mathcal{H} \) such that \( Tx = \lambda x \). Then

\[
\lambda a_{0} = 0 \quad \text{and} \quad a_{n} \omega_{n} = \lambda a_{n+1} \quad \text{for every} \quad n \in \mathbb{N}.
\]

Since \( T \) is injective then \( \lambda \neq 0 \). Therefore, \( a_{n} = 0 \) for every \( n \in \mathbb{N} \). So, \( x = 0 \). Thus \( \sigma_{p}(T) \) is empty.

(ii) We have \( T^{*}e_{0} = 0 \) then \( 0 \in \sigma_{p}(T^{*}) \). Let \( \lambda \in \sigma_{p}(T^{*}) \) and let \( x = \sum_{n \in \mathbb{N}} a_{n}e_{n} \) be a corresponding eigenvector that is \( T^{*}x = \lambda x \). And so, \( \lambda a_{n} = \omega_{n}a_{n+1} \) for every \( n \in \mathbb{N} \). Therefore,

\[
a_{n} = \frac{a_{0}\lambda^{n}}{\beta_{n}} \quad \text{for every} \quad n \geq 1.
\]

Hence \( a_{0} \neq 0 \) and \( x = a_{0}(e_{0} + \sum_{n \geq 1} \frac{\lambda^{n}}{\beta_{n}}e_{n}) \). Therefore the eigenvalues of \( T^{*} \) are simple and

\[
\|x\|^{2} = |a_{0}|^{2}(1 + \sum_{n \geq 1} \frac{|\lambda|^{2n}}{\beta_{n}^{2}}).
\]

By the Cauchy-Hadamard formula for the radius of convergence we get that \( |\lambda| \leq r_{2}(T) \). So,

\[
\sigma_{p}(T^{*}) \subset \{ \lambda \in \mathbb{C} : |\lambda| \leq r_{2}(T) \}.
\]

On the other hand, if \( \lambda < r_{2}(T) \) then the series \( \sum_{n \geq 1} \frac{|\lambda|^{2n}}{\beta_{n}^{2}} \) is convergent and so, \( x = e_{0} + \sum_{n \geq 1} \frac{\lambda^{n}}{\beta_{n}}e_{n} \) is an eigenvector corresponding to \( \lambda \). Thus

\[
\{ \lambda \in \mathbb{C} : |\lambda| < r_{2}(T) \} \subset \sigma_{p}(T^{*}).
\]

This proves the desired result. \( \blacksquare \)
Remark 1.3.13 By circular symmetry, one of the two containments in (ii) must be equality. If the unilateral weighted shift $T$ is not injective then $\sigma_p(T) = \{0\}$.

If $T$ is bilateral shift, set:

$$r_2^+ = \liminf_{n \to +\infty} (\beta_n)^{\frac{1}{n}}, \quad r_3^+ = \limsup_{n \to +\infty} (\beta_n)^{\frac{1}{n}},$$

$$r_2^- = \liminf_{n \to +\infty} \left(\frac{1}{\beta_n}\right)^{\frac{1}{n}}, \quad r_3^- = \limsup_{n \to +\infty} \left(\frac{1}{\beta_n}\right)^{\frac{1}{n}}.$$ 

Clearly we have

$$r_1^- \leq r_2^- \leq r_3^- \leq r^-, \quad r_1^+ \leq r_2^+ \leq r_3^+ \leq r^+.$$ 

Theorem 1.3.14 If $T$ is bilateral shift then:

(i) All eigenvalues of $T$ and $T^*$ are simple.

(ii) \(\{\lambda \in \mathbb{C} : r_2^- < |\lambda| < r_3^-\} \subset \sigma_p(T) \subset \{\lambda \in \mathbb{C} : r_3^- \leq |\lambda| \leq r_2^-\}\).

(iii) \(\{\lambda \in \mathbb{C} : r_2^+ < |\lambda| < r_3^+\} \subset \sigma_p(T^*) \subset \{\lambda \in \mathbb{C} : r_3^+ \leq |\lambda| \leq r_2^+\}\).

(iv) At least one of $\sigma_p(T)$, $\sigma_p(T^*)$ is empty.

Proof. Let $\lambda \in \sigma_p(T)$ and let $x = \sum_{n \in \mathbb{Z}} a_ne_n$ be an eigenvector corresponding to $\lambda$. Then

$$a_{n-1}\omega_{n-1} = \lambda a_n$$

for every $n \in \mathbb{Z}$.

And so,

$$a_n = \frac{a_0}{\beta_n} \lambda^n, \quad a_{-n} = a_0\beta_{-n} \lambda^n$$

for every $n \geq 1$.

From this we see that the eigenvalues are simple. Further, $\lambda$ is eigenvalue for $T$ if and only if the sequence $(a_n)_{n \in \mathbb{Z}}$ defined above is square summable. This leads to two power series, one in $\lambda$ and the other in $\beta$, and the result follows from the formula for the radius of convergence. The case of $T^*$ is similar.

Finally, assume that $\lambda$ and $\mu$ are eigenvalues of $T$ and $T^*$ respectively. We wish to show that this is impossible. By what has gone before we must have

$$r_3^- \leq |\lambda| \leq r_2^- \quad \text{and} \quad r_3^+ \leq |\mu| \leq r_2^+.$$ 

Since $r_2^- \leq r_3^-$ and $r_2^+ \leq r_3^+$, then $|\lambda| = |\mu|$. Also, an examination of the series which must converge shows that

$$\sum_{n \geq 1} |\mu|^{2n} < \infty \quad \text{and} \quad \sum_{n \geq 1} |\beta_n|^{2n} < \infty$$

which is impossible since $|\lambda| = |\mu|$.

Remark 1.3.15

(i) If $r_2^- < r_3^-$ then $\sigma_p(T)$ is empty; if $r_2^+ < r_3^-$ then $\sigma_p(T^*)$ is empty.

(ii) By circular symmetry, one of the containments in (ii), and one of the containments in (iii) must be equality.

(iii) Let $\omega_n = 1$ for every $n \in \mathbb{Z}$. Then $T$ is unitary bilateral shift and its adjoint which is its inverse is given by $T^*e_n = T^{-1}e_n = e_{n-1}$ for every $n$. On the other hand $T$ and $T^*$ are unitarily equivalent since $R^{-1}TR = T^*$ where $R$ is the unitary operator determined by $R e_n = e_{-n}$ for every $n \in \mathbb{Z}$. So, $T$ and $T^*$ have the same spectrum, the same point spectrum and the same approximate point spectrum. Hence from the property (iv) it follows that $\sigma_p(T) = \sigma_p(T^*) = \emptyset$.

Also, in this case we have $r(T) = r_1(T) = r_1^+ = r_1^- = 1$, $i = 1, 2, 3$. So,

$$\sigma(T) = \sigma_{ap}(T) = \sigma(T^{-1}) = \sigma_{ap}(T^{-1}) = \{\lambda \in \mathbb{C} : |\lambda| = 1\}.$$
Chapter 2

Analytic Extension and Spectral Theory

2.1 Operators with the single-valued Extension Property

In this section, we introduce a generalization of the spectrum and the resolvent set due to Dunford (see [?]). Let \( \mathcal{X} \) be a complex Banach space and \( \mathcal{L}(\mathcal{X}) \) be the algebra of all bounded operators on \( \mathcal{X} \). For an operator \( T \in \mathcal{L}(\mathcal{X}) \), the resolvent map \( R_T \) is an operator valued analytic function on the resolvent set \( \rho(T) \) of \( T \). For \( x \in \mathcal{X} \) the vector valued function \( R_{T,x} \) defined on \( \rho(T) \) by

\[
R_{T,x}(\lambda) = R_T(\lambda)x = (T - \lambda I)^{-1}x
\]

is analytic on \( \rho(T) \) and satisfying the following equation

\[
(T - \lambda I)R_{T,x}(\lambda) = x \quad \text{for every } \lambda \in \rho(T).
\]

In many cases the function \( R_{T,x} \) can be extended to be analytic on an open set properly containing \( \rho(T) \). We will call a vector valued analytic function \( F \) an analytic extension of \( R_{T,x} \) if the domain \( D(F) \) of \( F \) is containing \( \rho(T) \) and

\[
(T - \lambda I)F(\lambda) = x \quad \text{for every } \lambda \in D(F).
\]

We now encounter the possibility that there may be many extensions of \( R_{T,x} \) and they may not agree on their common domain. However if all extensions do agree on their common domain for each \( x \in \mathcal{X} \) we say that the operator \( T \) has the single valued extension property (s.v.e.p). Therefore, for such operator \( T \) and for every \( x \in \mathcal{X} \), the function \( R_{T,x} \) has a maximal extension called the maximal single valued extension of \( R_{T,x} \) and denoted by \( x(\cdot) \). Thus \( x(\cdot) \) is an analytic vector valued function such that

\[
(T - \lambda I)x(\lambda) = x \quad \text{for every } \lambda \in D(x);
\]

the open domain \( D(x) \) of \( x \) will be called the local resolvent set of \( x \) and is denoted by \( \rho_T(x) \) and its complement, denoted by \( \sigma_T(x) \), will be called the local spectrum of \( x \); it is a closed subset of the spectrum \( \sigma(T) \) of \( T \). We will see by using Liouville’s Theorem that the local spectrum of \( x \) is empty if and only if \( x \) is the zero vector.

Before proceeding further it would be well to note that, an operator \( T \in \mathcal{L}(\mathcal{X}) \) has the single valued extension property (s.v.e.p) if and only if for every open set \( U \subset \mathbb{C} \), the only analytic solution of the equation \((T - \lambda I)F(\lambda) = 0 \) for \( \lambda \in U \) is the zero function \( F \equiv 0 \). Therefore, if \( T \in \mathcal{L}(\mathcal{X}) \) has (s.v.e.p) then the local resolvent of every vector \( x \in \mathcal{X} \) is the set of complex numbers \( \lambda \) such that there exists an vector valued analytic function \( F \) defined in a
open neighborhood $V$ of $\lambda$ which verifies $(T - \mu I)F(\mu) = x$ for every $\mu \in V$.

**Example of an operator without s.v.e.p.** Let $T$ be the adjoint of the unilateral weighted shift operator with the weight $\omega_n = 1$ for every $n \in \mathbb{N}$. I.e:

$$T \epsilon_n = \begin{cases} 
0 & \text{if } n = 0 \\
e_n & \text{if } n > 0 
\end{cases}$$

Clearly, we have $\|T\| = 1$ and $\sigma(T) = \{\lambda \in \mathbb{C} : |\lambda| \leq 1\}$ then $\rho(T) = \{\lambda \in \mathbb{C} : |\lambda| > 1\}$. Therefore, for every $\lambda \in \rho(T)$ we have $\|\frac{T}{\lambda}\| < 1$ and so,

$$(T - \lambda I)^{-1} = \frac{1}{\lambda}(\frac{T}{\lambda} - I)^{-1}$$

$$= \frac{1}{\lambda} \sum_{n=0}^{+\infty} \left(\frac{T}{\lambda}\right)^n$$

$$= -\sum_{n=0}^{+\infty} \frac{T^n}{\lambda^{n+1}}.$$ 

Hence for every $\lambda \in \rho(T)$,

$$R_{T,e_0} = (T - \lambda I)^{-1}e_0$$

$$= -\sum_{n=0}^{+\infty} \frac{T^n e_0}{\lambda^{n+1}}$$

$$= -\frac{e_0}{\lambda} \text{ because } T^n e_0 = 0 \text{ for every } n \geq 0.$$ 

We will now exhibit two different analytic extensions of $R_{T,e_0}$. First let $F(\lambda) = -\frac{e_0}{\lambda}$ for $\lambda \neq 0$; then $(T - \lambda I)F(\lambda) = e_0$ for every $\lambda \neq 0$. And so, $F$ is clearly an analytic extension of $R_{T,e_0}$. Second, set

$$G(\lambda) = \begin{cases} 
-\frac{e_0}{\lambda} & \text{if } |\lambda| > 1 \\
+\infty \sum_{n=0}^{+\infty} \lambda^n e_{n+1} & \text{if } |\lambda| < 1.
\end{cases}$$

Clearly, that $G$ has the right form for $\lambda \in \mathbb{C}$, $|\lambda| > 1$. For $\lambda \in \mathbb{C}$, $|\lambda| < 1$ we have

$$(T - \lambda I)G(\lambda) = (T - \lambda I) \sum_{n=0}^{+\infty} \lambda^n e_{n+1}$$

$$= \sum_{n=0}^{+\infty} \lambda^n T e_{n+1} - \sum_{n=0}^{+\infty} \lambda^{n+1} e_{n+1}$$

$$= \sum_{n=0}^{+\infty} \lambda^n e_n - \sum_{n=0}^{+\infty} \lambda^{n+1} e_{n+1}$$

$$= e_0.$$ 

Thus $G$ is also an analytic extension of $R_{T,e_0}$. However the two analytic extensions $F$ and $G$ do not coincide on $\{\lambda \in \mathbb{C} : 0 < |\lambda| < 1\}$.

**Note.** The operator $T$ considered in this example is the adjoint of the of the unilateral weighted shift operator with the weight $\omega_n = 1$ for every $n \in \mathbb{N}$ which is an isometric non unitary operator. More than that one can prove that the adjoint operator of every isometric non unitary operator
on Hilbert space do not has the s.v.e.p (see [?]).

**Example of operators having the s.v.e.p.** Every operator \( T \in \mathcal{L}(\mathcal{X}) \) which has empty interior of its point spectrum has the s.v.e.p. Let \( F \) be a vector valued analytic function such that

\[
(T - \lambda I)F(\lambda) = 0 \quad \text{for every } \lambda \in D(F).
\]

Then

\[
TF(\lambda) = \lambda F(\lambda) \quad \text{for every } \lambda \in D(F).
\]

And so, \( F \) must be identically zero otherwise there is \( \lambda \in D(F) \) and \( r > 0 \) such that \( B_r = \{ \mu \in \mathbb{C} : |\lambda - \mu| < r \} \subset D(F) \) and \( F(\mu) \neq 0 \) for every \( \mu \in B_r \). Thus, \( B_r \subset \sigma_p(T) \). Contradiction. Hence, \( F \equiv 0 \) and so, \( T \) has the s.v.e.p.

**Note.** From Theorem 1.3.12 and Remark 1.3.13 follow that every unilateral weighted shift operator has the s.v.e.p.

We will see that all operators in which we are interested have the s.v.e.p. So, from now on we will be dealing with operators which have only single-valued extension property (s.v.e.p).

**Proposition 2.1.1** Let \( T \in \mathcal{L}(\mathcal{X}) \) be an operator which has the s.v.e.p and let \( x, y \) be two vectors of \( \mathcal{X} \). Then:

(i) \( \sigma_T(\alpha x) = \sigma_T(x) \) for every \( \alpha \neq 0 \).

(ii) \( \sigma_T(x + y) \subset \sigma_T(x) \cup \sigma_T(y) \).

(iii) \( \sigma_T(x) \) is empty if and only if \( x = 0 \).

(iv) \( \sigma_T(\tilde{x}(\lambda)) = \sigma_T(x) \) for every \( \lambda \in \rho_T(x) \) fixed point.

(v) For every operator \( S \in \mathcal{L}(\mathcal{X}) \) which commutes with \( T \), \( \sigma_T(Sx) \subset \sigma_T(x) \).

**Proof.**

(ii) The result follows immediately since

\[
(T - \lambda I)(\tilde{x}(\lambda) + \tilde{y}(\lambda)) = x + y \quad \text{for every } \lambda \in \rho_T(x) \cap \rho_T(y).
\]

(iii) If \( x = 0 \) then there noting to prove. Now, suppose that \( \sigma_T(x) = \emptyset \) then the maximal single valued extension \( \tilde{x} \) is an entire function which coincide with \( R_{T,x} \) on \( \rho(T) \). Since \( \{ \lambda \in \mathbb{C} : |\lambda| > \|T\| \} \subset \rho(T) \) and \( \tilde{x}(\lambda) = R_{T,x}(\lambda)x \) for \( |\lambda| > \|T\| \) then from the Theorem 1.2.1 it follows that \( \lim_{|\lambda| \to +\infty} \tilde{x}(\lambda) = 0 \). Therefore by Liouville’s Theorem

\[
\tilde{x} \equiv 0.
\]

In particular, \( R_{T,x} \equiv 0 \) on \( \rho(T) \). Hence \( x = 0 \).

(iv) Fix \( \lambda \in \rho_T(x) \) and let \( y = \tilde{x}(\lambda) \). Let we first prove that \( \sigma_T(x) \subset \sigma_T(y) \). We have

\[
(T - \mu I)\tilde{y}(\mu) = y \quad \text{for every } \mu \in \rho_T(y)
\]

then for every \( \mu \in \rho_T(y) \) we have

\[
(T - \lambda I)(T - \mu I)\tilde{y}(\mu) = (T - \mu I)(T - \lambda I)\tilde{y}(\mu) = (T - \lambda I)y = (T - \lambda I)\tilde{x}(\lambda) \quad \text{because } y = \tilde{x}(\lambda) = x
\]

Since \( (T - \lambda I)\tilde{y} \) is a vector valued analytic function on \( \rho_T(y) \) then it is analytic extension of \( R_{T,x} \). Therefore, \( \rho_T(y) \subset \rho_T(x) \). Thus \( \sigma_T(x) \subset \sigma_T(y) = \sigma_T(\tilde{x}(\lambda)) \).

Conversely, let use to prove that \( \sigma_T(\tilde{x}(\lambda)) \subset \sigma_T(x) \). Clearly that the vector valued function \( F \) defined in \( \rho_T(x) \) by

\[
F(\lambda) = \begin{cases} \frac{\tilde{x}(\mu) - \tilde{x}(\lambda)}{\mu - \lambda} & \text{if } \mu \neq \lambda \\ \tilde{x}(\lambda) & \text{if } \mu = \lambda \end{cases}
\]
is an analytic function on \( \rho_T(x) \) and for \( \mu \neq \lambda \) we have

\[
(T - \mu I) F(\mu) = (T - \mu I) \frac{\bar{x}(\mu) - \bar{x}(\lambda)}{\mu - \lambda}
\]

\[
= \frac{1}{\mu - \lambda} \left[ (T - \mu I)\bar{x}(\mu) - (T - \mu I)\bar{x}(\lambda) \right]
\]

\[
= \frac{1}{\mu - \lambda} \left[ x - (T - \lambda I)\bar{x}(\lambda) + (\mu - \lambda)\bar{x}(\lambda) \right]
\]

\[
= \frac{1}{\mu - \lambda} \left[ x - x + (\mu - \lambda)\bar{x}(\lambda) \right]
\]

\[
= \bar{x}(\lambda)
\]

\[
= y
\]

This equality holds also for \( \mu = \lambda \) by making \( \mu \to \lambda \) in the last equality, and so

\[
(T - \mu I) F(\mu) = \bar{x}(\lambda) \text{ for every } \lambda \in \rho_T(x).
\]

Hence, \( \rho_T(x) \subset \rho_T(\bar{x}(\lambda)) \). Thus \( \sigma_T(\bar{x}(\lambda)) \subset \sigma_T(x) \). And so, \( \sigma_T(x) = \sigma_T(\bar{x}(\lambda)) \) for every \( \lambda \in \rho_T(x) \).

(v) We have \( (T - \lambda I)\bar{x}(\lambda) = x \) for every \( \lambda \in \rho_T(x) \) then \( S(T - \lambda I)\bar{x}(\lambda) = Sx \) for every \( \lambda \in \rho_T(x) \) and so,

\[
(T - \lambda I)S\bar{x}(\lambda) = Sx \text{ for every } \lambda \in \rho_T(x).
\]

Hence the vector valued function \( S\bar{x} \) is an analytic extension of \( R_{T,Sx} \). Therefore, \( \rho_T(x) \subset \rho_T(Sx) \). Thus, \( \sigma_T(Sx) \subset \sigma_T(x) \).

\[\text{Theorem 2.1.2} \quad \text{For every operator } T \in \mathcal{L}(\mathcal{X}) \text{ which has the s.v.e.p} \]

\[\sigma(T) = \bigcup_{x \in \mathcal{X}} \sigma_T(x).\]

\[\text{Proof.} \quad \text{Assume that there is } \lambda \in \sigma(T) \text{ such that } \lambda \notin \sigma_T(x) \text{ for every } x \in \mathcal{X}. \]

Then \( \lambda \notin \bigcap_{x \in \mathcal{X}} \rho_T(x) \) and so, \( (T - \lambda I)\bar{x}(\lambda) = x \) for every \( x \in \mathcal{X} \), in particular the operator \( T - \lambda I \) is surjective. Therefore it follows from the open mapping Theorem (see [?] and [?]) that \( T - \lambda I \) must be not injective since \( \lambda \in \sigma(T) \). So, there is \( x_1 \neq 0 \in \mathcal{X} \) such that \( (T - \lambda I)x_1 = 0 \). Since \( T - \lambda I \) is bounded surjective operator then it follows from the Open Mapping Theorem that there is a constant \( c > 0 \) such that for every \( x \in \mathcal{X} \) there exists \( y \in \mathcal{X} \) with

\[\|y\| < c\|x\| \text{ and } (T - \lambda I)y = x.\]

Therefore, there is a sequence \( (x_n) \) of elements of \( \mathcal{X} \) such that

\[(T - \lambda I)x_{n+1} = x_n \text{ and } \|x_{n+1}\| \leq c\|x_n\| \text{ for every } n \geq 1.\]

And so, the function \( F \) defined on the open disc \( D = \{ \mu \in \mathbb{C} : |\mu - \lambda| < \frac{1}{c} \} \) by

\[F(\mu) = \sum_{n=0}^{\infty} x_{n+1}(\mu - \lambda)^n\]
2.2 OPERATORS WITH THE DUNFORD’S CONDITION C (DCC)

is a non identically zero analytic function on $D$ and for every $\mu \in D$ we have

$$(T - \mu I)F(\mu) = (T - \lambda)F(\mu) + (\lambda - \mu)F(\mu)$$

$$= \sum_{n=0}^{+\infty} (T - \lambda)\left(x_{n+1}(\mu - \lambda)^n\right) - (\lambda - \mu)\sum_{n=0}^{+\infty} x_{n+1}(\mu - \lambda)^n$$

$$= \sum_{n=1}^{+\infty} x_n(\mu - \lambda)^n - \sum_{n=0}^{+\infty} x_{n+1}(\mu - \lambda)^{n+1}$$

$$= 0$$

Contradiction with the assumption that $T$ has the s.v.e.p, and so the desired result holds.

**Lemma 2.1.3** Let $T \in \mathcal{L}(X)$ be a surjective bounded operator which has the s.v.e.p. Then $T$ is invertible in $\mathcal{L}(X)$. Therefore,

$$\sigma(T) = \{ \lambda \in \mathbb{C} : T - \lambda I \text{ is not onto} \}.$$

**Proof.** Suppose that $T$ is not invertible in $\mathcal{L}(X)$ then it is not injective. So, there is an element $x_1 \neq 0 \in X$ such that $Tx_1 = 0$. And so, by the same argument of the proof of Theorem 2.1.2 it follows that $T$ has not the s.v.e.p.

**Theorem 2.1.4** Let $T \in \mathcal{L}(\mathcal{H})$ be an bounded operator on a Hilbert space $\mathcal{H}$. If its adjoint $T^*$ has s.e.v.p then

$$\sigma(T) = \sigma_{ap}(T).$$

**Proof.** It is known that $\sigma_{ap}(T) \subset \sigma(T)$. Conversely, let $\lambda \notin \sigma_{ap}(T)$ then $T - \lambda I$ is bounded below. So, it follows from Theorem 1.1.8 that $T^* - \overline{\lambda} I$ is surjective. Hence by the Lemma 2.1.3 $\overline{\lambda} \notin \sigma(T^*)$. Thus $\lambda \notin \sigma(T)$ since $\sigma(T) = \overline{\sigma(T^*)}$. This complete the proof of the Theorem.

2.2 Operators with the Dunford’s Condition C (DCC)

Let $T \in \mathcal{L}(\mathcal{X})$ be an operator which has s.v.e.p. So, from the properties (i) and (ii) of Proposition 2.1.1 it follows that for every closed subset $F$ of $\mathbb{C}$ the following subset of $\mathcal{X}$ given by

$$\mathcal{X}_T(F) = \{ x \in \mathcal{X} : \sigma_T(x) \subset F \}$$

is a linear subspace of $\mathcal{X}$ not necessarily closed as we will see in the example below. The operator $T$ is said to satisfy the **Dunford’s Condition C** (DCC) if for every closed subset $F$ of $\mathbb{C}$ the linear subspace $\mathcal{X}_T(F)$ is closed. We will see that all the operators in which we are interested have the Dunford’s Condition C (DCC). Also, from the property (v) of the same proposition we see that the linear subspace $\mathcal{X}_T(F)$ is invariant linear subspace for every operator $S \in \mathcal{L}(\mathcal{X})$ which commutes with $T$. i.e:

$$S\left(\mathcal{X}_T(F)\right) \subset \mathcal{X}_T(F) \text{ for every } S \in \mathcal{L}(\mathcal{X}), \text{ with } ST = TS.$$

If the linear subspace $\mathcal{X}_T(F)$ is a closed and not trivial (i.e: $\mathcal{X}_T(F) \neq \{0\}$ and $\mathcal{X}_T(F) \neq \mathcal{X}$) then it is a non trivial hyperinvariant linear subspace for $T$ and this is one method to construct the non trivial hyperinvariant subspaces for a given operator in $\mathcal{L}(\mathcal{X})$. (More information about the existence of hyperinvariant subspaces of operators can be found in [2], [3] and [4]).

Since for every $x \in \mathcal{X}$ we have $\sigma_T(x) \subset \sigma(T)$ then obviously we have $\mathcal{X}_T(F) = \mathcal{X}_T(F \cap \sigma(T))$, hence $F$ can be supposed to be a closed subset of $\sigma(T)$.

Before proceeding further it would be well to point out that the Dunford’s Condition C does
not hold for every operator which has the single valued extension property. Let us consider an example. Let \((e_n)_n\) be an orthonormal basis for a Hilbert space \(H\), and let \(T\) be the unilateral weighted shift operator given by
\[
T e_n = \begin{cases} 
e_n+1 & \text{if } n \text{ is not square integer}, \\ 0 & \text{if } n \text{ is square integer}. \end{cases}
\]

It is easy to see that \(\|T^n\| = 1\) for every \(n \geq 1\), then spectral radius \(r(T)\) of \(T\) is 1. On the other hand, \(T\) has the s.v.e.p property since \(\sigma_p(T) = \{0\}\) (see Remark 1.3.13) Now let us suppose that \(X_T(\{0\})\) is closed. For every non-negative integer \(n\) there is an integer \(k\) such that \(k^2 \leq n < (k+1)^2\) and so,

- If \(n = k^2\) then
  \[
  (T - \lambda I) \left( - \frac{e_n}{\lambda} \right) = e_n \quad \text{for } \lambda \neq 0.
  \]

- If \(n = k^2 + s\) for some \(0 < s < 2k + 1\) then
  \[
  (T - \lambda I) \left( - \frac{e_n}{\lambda} - \frac{e_{n+1}}{\lambda^2} - \ldots - \frac{e_{(k+1)^2-s}}{\lambda^{2(k+1)-s}} \right) = e_n \quad \text{for } \lambda \neq 0.
  \]

Then for every integer \(n \geq 0\), \(\sigma_T(e_n) = \{0\}\). And so, \(X_T(\{0\}) = H\). Hence it follows from Theorem 2.1.2 that \(\sigma(T) = \{0\}\). This is impossible since the spectral radius \(r(T)\) of \(T\) is 1.

**Theorem 2.2.1** For every closed subset \(F\) of \(\sigma(T)\) such that the linear subspace \(X_T(F)\) is closed,
\[
\sigma(T|_{X_T(F)}) \subset F.
\]

**Proof.** Let \(\lambda \notin F\). We will prove that \(\lambda \notin \rho(T|_{X_T(F)})\), that is \(T|_{X_T(F)} - \lambda I\) is invertible in \(\mathcal{L}(X_T(F))\). Since for every \(x \in X_T(F)\) we have \(\mathbb{C}\setminus F \subset \rho_T(x)\). Therefore \(\tilde{x}(\lambda)\) makes sense and from the equality (iv) of Proposition 2.1.1 it follows that \(\rho_T(x) = \rho_T(\tilde{x}(\lambda))\). Thus \(\tilde{x}(\lambda) \in X_T(F)\).

Let \(A\) be the map from \(X_T(F)\) to itself defined by \(Ax = \tilde{x}(\lambda)\). It is evident that \(A\) is linear map. We will show by using the Closed Graph Theorem (see [?] and [?]) that the linear map \(A\) is bounded. Let \((x_n)_n\) be a sequence of elements of \(X_T(F)\) such that \(\lim x_n = x \in X\) and \(\lim Ax_n = y \in X\). Since \(X_T(F)\) is closed linear subspace of \(X\) then \(x \in X_T(F)\) and \(y = \lim Ax_n = \lim \tilde{x}(\lambda) \in X_T(F)\). We have
\[
(T - \lambda I)Ax_n = (T - \lambda I)\tilde{x}(\lambda) = x_n,
\]

from which it follows by the continuity of the linear operator \((T - \lambda I)\) that \((T - \lambda I)y = x\). On the other hand, \((T - \lambda I)\tilde{x}(\lambda) = x\), therefore, \((T - \lambda I)\left( \tilde{x}(\lambda) - y \right) = 0\). Since \(X_T(F)\) is a linear subspace and \(y, \tilde{x}(\lambda) \in X_T(F)\) then \(z = \tilde{x}(\lambda) - y \in X_T(F)\). Hence
\[
\sigma_T(z) \subset F.
\]

Let \(G\) be the vector valued analytic function defined on \(\mathbb{C}\setminus \{\lambda\}\) by \(G(\mu) = \frac{1}{\lambda - \mu} z\). Then for every \(\mu \in \mathbb{C}\setminus \{\lambda\}\) we have
\[
(T - \mu I)G(\mu) = \frac{1}{\lambda - \mu}(T - \mu I)z
\]
\[
= \frac{1}{\lambda - \mu} \left( (T - \lambda I)z + (\lambda - \mu)z \right)
\]
\[
= \frac{1}{\lambda - \mu} \left( 0 + (\lambda - \mu)z \right)
\]
\[
= z.
\]
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Hence \( \mathbb{C} \setminus \{\lambda\} \subset \rho_T(z) \). Thus

\[
\sigma_T(z) \subset \{\lambda\}.
\]

Since \( \lambda \notin F \) then from \((*)\) and \((***)\) it follows that \( \sigma_T(z) \subset F \cap \{\lambda\} = \emptyset \). Hence \( \sigma_T(z) = \emptyset \) and so, \( z = 0 \) by the property (iii) of Proposition 2.1.1. Therefore,

\[
Ax = \bar{x}(\lambda) = y.
\]

And so, by the Closed Graph Theorem the linear operator \( A \) is bounded on \( X_T(F) \).

Now, let us prove that the operator \( A \) is exactly the inverse of \( T|X_T(F) - \lambda I \). For every \( x \in X_T(F) \) we have

\[
(T|X_T(F) - \lambda I)Ax = (T - \lambda I)\bar{x}(\lambda) = x.
\]

On the other hand for every \( x \in X_T(F) \), we have \( (T - \lambda I)x(\lambda) = (T - \lambda I)\bar{x}(\lambda) \), for which, according to the definition of \( A \), it follows that

\[
A(T|X_T(F) - \lambda I)x = (T - \lambda I)x(\lambda) = (T - \lambda I)\bar{x}(\lambda) = x.
\]

Therefore \( \lambda \in \rho(T|X_T(F)) \).

We shall see an application of this Theorem. The operator \( T \) is said to be cyclic if there is a vector \( x \in \mathcal{X} \) such that the linear subspace generated by \( \{T^n x : n \in \mathbb{N}\} \) is dense in \( \mathcal{X} \). The vector \( x \) is called a cyclic vector for \( T \). If \( T \) is an injective unilateral weighted shift with weights \( (\omega_n)_n \) that is \( Te_n = \omega_ne_{n+1}, n \geq 0 \) where \( (e_n)_n \) is a orthonormal basis of \( H \), then \( T \) is cyclic with cyclic vector \( e_0 \) since \( T^n e_0 = \omega_0...\omega_{n-1}e_n \) for every \( n \geq 1 \). (For more information about cyclic shifts see [2] and [3]).

**Proposition 2.2.2** Suppose that the operator \( T \) is cyclic with a cyclic vector \( x \in \mathcal{X} \). If the operator \( T \) satisfies DCC then \( \sigma_T(x) = \sigma(T) \).

**Proof.** Let \( F = \sigma_T(x) \). Then \( x \in X_T(F) \). Since \( X_T(F) \) is invariant subspace for \( T \) then \( T^n x \in X_T(F) \) for every \( n \in \mathbb{N} \). Hence, the linear subspace generated by \( \{T^n x : n \in \mathbb{N}\} \) is contained in \( X_T(F) \). Since the operator \( T \) satisfies DCC property then \( X_T(F) \) is closed subspace and so it follows from the density in \( \mathcal{X} \) of the linear subspace generated by \( \{T^n x : n \in \mathbb{N}\} \) that \( \mathcal{X} = X_T(F) \). And so, from Theorem 2.2.1 it follows that

\[
\sigma(T) = \sigma(T|X_T(F)) \subset F = \sigma_T(x).
\]

Therefore, \( \sigma(T) = \sigma_T(x) \).

**Remark 2.2.3** In this proposition we did not use the DDC property of the operator \( T \). We used only the fact that \( X_T(F) \) is closed linear subspace where \( F = \sigma_T(x) \).
Chapter 3

Subnormal and Hyponormal Operators

Several classes of Hilbert spaces operators are defined around the notion of normal operator. In 1950, Paul R. Halmos, motivated by the successful development of the theory of normal operators, introduced the notions of subnormality and hyponormality for bounded Hilbert space operators, in an attempt to extend the basic facts of the spectral theory of normal operators. Those classes of operators have been the subject of much investigation during the last fifty years and many important developments in the operator theory have dealt with them e.g: S.Brown’s proof of the existence of non trivial subspaces, J.Conway and R.Olin’s construction of the functional calculus and J.Thomson’s description of the spectral picture in the cyclic case for subnormal operators, etc. Our goal in this chapter is to study the structure of those classes of operators and to collect sufficient material to study the local spectra of cyclic hyponormal operators in the next chapter.

3.1 Subnormal and Hyponormal Operators

In this section, we recall first the general and fundamental properties of normal bounded Hilbert space operators. Let $\mathcal{L}(\mathcal{H})$ be the algebra of all linear bounded operators on a Hilbert space $\mathcal{H}$. An operator $T \in \mathcal{L}(\mathcal{H})$ is said to be normal if it commutes with its adjoint $T^*$ i.e: $T^*T = TT^*$. We first note that for every operator $T \in \mathcal{L}(\mathcal{H})$,

$$\langle (T^*T - TT^*)x, x \rangle = \|Tx\|^2 - \|T^*x\|^2$$

for every $x \in \mathcal{H}$. (3.1)

**Proposition 3.1.1** An operator $T \in \mathcal{L}(\mathcal{H})$ is normal if and only if $\|Tx\| = \|T^*x\|$ for every $x \in \mathcal{H}$.

**Proof.** Since the zero operator is the unique operator $S \in \mathcal{L}(\mathcal{H})$ which has the property that

$$\langle Sx, x \rangle = 0$$

for every $x \in \mathcal{H}$, then the proof follows immediately from the identity (3.1).

Using Liouville’s Theorem, Fugled Putman in 1950 give a nice characterization of normal operators.

**Theorem 3.1.2 (Fugled)** An operator $T \in \mathcal{L}(\mathcal{H})$ is normal if and only if $T^*S = ST^*$ for every operator $S \in \mathcal{L}(\mathcal{H})$ such that $TS = ST$. 24
In the proof of this theorem, we shall need the following:

Let $T$ be an operator in $\mathcal{L}(\mathcal{H})$. For every positive integer $n$, let $T_n = I + \frac{T}{1!} + \frac{T^2}{2!} + \cdots + \frac{T^n}{n!}$. $T_n$ is a bounded operator in $\mathcal{L}(\mathcal{H})$; and the sequence $(T_n)_n$ converges in $\mathcal{L}(\mathcal{H})$ to an invertible operator denoted by $e^T$ and its inverse is $e^{-T}$ i.e:

$$
(e^T)^{-1} = e^{-T};
$$

(3.2)

in addition, if $S$ is an operator in $\mathcal{L}(\mathcal{H})$ which commutes with $T$ then:

$$
e^T S = S e^T,
$$

(3.3)

and

$$
e^{T+S} = e^T e^S = e^S e^T.
$$

(3.4)

**Proof of Theorem 3.1.2.** Let $S$ be an operator in $\mathcal{L}(\mathcal{H})$ such that $TS = ST$. Then it follow from (3.2) and (3.3) that $S = e^{-izT} S e^{izT}$ for every $z \in \mathbb{C}$. And so, for every $z \in \mathbb{C}$,

$$
e^{-izT} S e^{izT} = e^{-izT} e^{-izT} S e^{izT} e^{izT} = e^{-i(z^2 + iz)} S e^{i(z^2 + iz)}
$$

Since the operator $e^{i(z^2 + iz)}$ is unitary and its adjoint is $e^{-i(z^2 + iz)}$ then

$$
e^{-izT} S e^{izT} = \left(e^{i(z^2 + iz)}\right)^* S e^{i(z^2 + iz)};
$$

therefore the following operator valued function defined on $\mathbb{C}$ by $\phi(z) = e^{-izT} S e^{izT}$ is bounded analytic function, by Liouville’s Theorem it is a constant function; in particular its derivative $\phi'$ is zero. So,

$$
\phi'(z) = -iT* e^{-izT} S e^{izT} + e^{-izT} S (iT*) e^{izT} = -iT* \phi(z) + e^{-izT} S e^{izT} (iT*) \quad \text{by (1.5)}
$$

$$
= -iT^* \phi(z) + \phi(z) (iT*) = 0.
$$

Hence,

$$
T^* \phi(z) = \phi(z) T^* \quad \text{for every } z \in \mathbb{C}.
$$

Thus $T^* S = ST^*$ because $\phi(0) = S$. 

Let $T$ be the unweighted bilateral shift i.e: $Te_n = e_{n+1}$ for every $n \in \mathbb{Z}$ where $(e_n)_n$ is an orthonormal basis for $\mathcal{H}$. Then $T$ is unitary operator since by Proposition 1.3.2, we have $T^* e_n = e_{n-1}$ for every $n \in \mathbb{Z}$. So, in particular, it is normal operator. In fact, one can see that a bilateral weighted shift operator $T$ with a nonnegative weights $(\omega_n)_{n \in \mathbb{Z}}$ is normal if and only if the weights $(\omega_n)_{n \in \mathbb{Z}}$ are constant; and a non zero unilateral shift is never normal. We note that the restriction of a normal operator $T \in \mathcal{L}(\mathcal{H})$ on a proper closed invariant subspace is not in general normal operator. Let $S$ be the restriction of the unweighted bilateral shift $T$ on the closed linear subspace $K$ generated by $\{e_n : n \geq 0\}$ which is proper closed invariant subspace for $T$. So, the restriction $S$ is exactly the unweighted unilateral shift on the Hilbert space $K$ which is not normal but it has a normal extension. Therefore it is natural to introduce and study the theory of subnormal operators which constitute a considerably more useful and deeper generalization of the theory of the normal operators. An operator $T \in \mathcal{L}(\mathcal{H})$ is called a **subnormal operator** if it has a normal extension i.e: if there exists a normal operator $S$
Lemma 3.1.3 Let $T \in \mathcal{L}(\mathcal{H})$ be a subnormal operator. A normal extension $S$ of $T$ on an Hilbert space $K$ is a minimal normal extension if and only if $K$ coincides with the closure of the linear subspace generated by $\{S^nx : n \in \mathbb{N}\}$.

Proof. It suffices to observe that the closed linear subspace $K_0$ generated by $\{S^nx : n \in \mathbb{N}\}$ is containing $\mathcal{H}$, invariant by $S$ and $S^*$ and the restriction of $S$ over $K_0$ is normal. 

Lemma 3.1.4 Let $S$ be a normal extension on a Hilbert space $K$ of an subnormal operator $T \in \mathcal{L}(\mathcal{H})$. Then for every finite family of elements $x_1, x_2, ..., x_n \in \mathcal{H}$ we have,

$$\| \sum_{i=1}^{n} S^{*i}x_i \| = \| \sum_{i=1}^{n} T^{*i}x_i \|.$$

Proof. We have

$$\| \sum_{i=1}^{n} S^{*i}x_i \|^2 = \langle \sum_{i=1}^{n} S^{*i}x_i, \sum_{j=1}^{n} S^{*j}x_j \rangle = \sum_{1 \leq i, j \leq n} \langle S^{*i}x_i, S^{*j}x_j \rangle = \sum_{1 \leq i, j \leq n} \langle S^ix_i, S^jx_j \rangle = \sum_{1 \leq i, j \leq n} \langle T^ix_i, T^jx_j \rangle \quad \text{because } S|_{\mathcal{H}} = T = \| \sum_{i=1}^{n} T^{*i}x_i \|^2.$$

Theorem 3.1.5 If $S_1$ and $S_2$ are two minimal normal extensions, on Hilbert spaces respectively $K_1$ and $K_2$, of a subnormal operator $T \in \mathcal{H}$ then there exists an invertible isometry $U$ from $K_1$ onto $K_2$ such that $US_1 = S_2U$ and $U(x) = x$ for every $x \in \mathcal{H}$.

Proof. Let $M_1$ and $M_2$ be the linear subspaces respectively of $K_1$ and $K_2$ generated respectively by the sets $\{S_1^ix : i \in \mathbb{N} \text{ and } x \in \mathcal{H}\}$ and $\{S_2^ix : i \in \mathbb{N} \text{ and } x \in \mathcal{H}\}$. It follows from the Lemma 3.1.3 that $K_1 = M_1$ and $K_2 = M_2$. On the other, it follows from the Lemma 3.1.4 that for every finite family of elements $x_1, x_2, ..., x_n \in \mathcal{H}$

$$\| \sum_{i} S_1^{*i}x_i \| = \| \sum_{i} S_2^{*i}x_i \|.$$

Hence the correspondence defined from $M_1$ onto $M_2$ by $\sum_{i} S_1^{*i}x_i \rightarrow \sum_{i} S_2^{*i}x_i$ is an isometry, we denote it by $U$. Therefore $U$ has a unique isometric extension that maps $K_1$ onto $K_2$ and it is the identity on $\mathcal{H}$. To prove that $US_1 = S_2U$, it suffices to verify that $US_1$ agrees with $S_2U$ on
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and this is implied by

\[ M_1 \text{ and this is implied by} \]

\[ \begin{align*}
    U S_1 \left( \sum_i S_1^i x_i \right) &= U \left( \sum_i S_1 S_1^i x_i \right) = U \left( \sum_i S_1^i S_1 x_i \right) \\
    &\quad = U \left( \sum_i S_1^i T x_i \right) \quad = \left( \sum_i S_2^i T x_i \right) \\
    &\quad = \left( \sum_i S_2^i S_2 x_i \right) = S_2 \left( \sum_i S_2^i x_i \right) \\
    &\quad = S_2 U \left( \sum_i S_1^i x_i \right).
\end{align*} \]

This complete the proof of the Theorem.

\[ \square \]

**Theorem 3.1.6** If \( T \in \mathcal{L}(\mathcal{H}) \) is a subnormal operator and \( S \) is its minimal normal extension on a Hilbert space \( \mathcal{K} \) then \( \sigma(S) \subset \sigma(T) \). In particular, the minimal normal extension of an invertible subnormal operator is invertible.

**Proof.** To prove that \( \sigma(S) \subset \sigma(T) \) it suffices to prove that \( S - \lambda I \) is invertible for every \( \lambda \in \rho(T) \). On the other hand, \( S - \lambda I \) is a minimal normal extension of the operator \( T - \lambda I \) so, the assertion reduces to prove that if \( T \) is invertible then its minimal normal extension \( S \) is invertible. It is clear that the closure of the range of \( S \) is a closed invariant subspace \( M \) for \( S \) and for its adjoint \( S^* \). Therefore the restriction of \( S \) on \( M \) is normal. On the other hand, \( \mathcal{H} = T \mathcal{H} = S \mathcal{H} \) is contained in the range of \( S \). Hence it follows from the minimality of \( S \) that \( M = \mathcal{K} \). So, \( S \) has dense range, therefore it suffices to prove that \( S \) is bounded from below. Let \( L \) be the linear subspace generated by the set \( \{ S^* x : i \in \mathbb{N} \text{ and } x \in \mathcal{H} \} \) then by the Lemma 3.1.3 it follows that \( L = \mathcal{K} \). For every finite sum \( \sum_i S^i x_i \in L \) we have

\[ \| S \left( \sum_i S^i x_i \right) \| = \| S^* \left( \sum_i S^i x_i \right) \| \quad \text{by proposition 3.1.1} \]

\[ = \left\| \sum_i S^i S^i x_i \right\| \]

\[ = \left\| \sum_i T^i S^i x_i \right\| \quad \text{Lemma (3.1.4)} \]

\[ = \left\| T^* \left( \sum_i T^i x_i \right) \right\| \]

\[ \geq \frac{1}{\| T^{-1} \|} \left\| \sum_i T^i x_i \right\| = \frac{1}{\| T^{-1} \|} \left\| \sum_i S^i x_i \right\| \quad \text{Lemma (3.1.4)}. \]

Thus \( \frac{1}{\| T^{-1} \|} \| x \| \leq \| S x \| \) for every \( x \in \mathcal{K} \); and the desired result follows.

\[ \square \]

**Theorem 3.1.7** For every subnormal operator \( T \in \mathcal{L}(\mathcal{H}) \), \( \| T^* x \| \leq \| T x \| \) for every \( x \in \mathcal{H} \).

**Proof.** First let us prove that \( T^* x = P S^* x \) for every \( x \in \mathcal{H} \) where \( S \) a minimal normal extension on a Hilbert space \( \mathcal{K} \) and \( P \) is the linear projection from \( \mathcal{K} = \mathcal{H} \oplus \mathcal{H}^\perp \) onto \( \mathcal{H} \). For every \( x, y \in \mathcal{H} \) we have

\[ \langle T^* x, y \rangle = \langle x, Ty \rangle = \langle x, S y \rangle \quad \text{because } S|_{\mathcal{H}} = T \]

\[ = \langle S^* x, y \rangle = \langle S^* x, P y \rangle \quad \text{because } P|_{\mathcal{H}} = I \]

\[ = \langle P^* S^* x, y \rangle = \langle P S^* x, y \rangle. \]
Therefore $T^*x = PS^*x$ for every $x \in \mathcal{H}$. And so, for every $x \in \mathcal{H}$ we have
\[
\|T^*x\| = \|PS^*x\| \leq \|S^*x\| = \|Sx\| = \|T_x\| \quad \text{because } \|P\| = 1.
\]

**Question.** Does the converse of Theorem 3.1.7 hold?

The answer is negative. Let us consider the following operator $T = S^* + 2S$ where $S$ is the unweighted shift in $\mathcal{H}$ i.e. $Se_n = e_{n+1}$ for every $n \in \mathbb{N}$ with $(e_n)_n$ is an orthonormal basis of $\mathcal{H}$.

A simple computation show that for every $x = \sum \alpha_n e_n \in \mathcal{H}$ we have
\[
\langle (T^*T - TT^*)x, x \rangle = 3\alpha_0^2 \geq 0.
\]

And so, from the equality (3.1) it follows that $\|T^*x\| \leq \|T_x\|$ for every $x \in \mathcal{H}$. On the other hand, for $x = e_0 - 2e_2$ we have
\[
\|T^*x\| = \|T^2x\| = \sqrt{89} > 80 = \|T^2a;\|.
\]

Hence by Theorem 3.1.7, $T^2$ is not subnormal operator. Therefore $T$ is not subnormal operator since every power of a subnormal operator is subnormal.

So, it is natural to study this class of operators $T \in \mathcal{L}(\mathcal{H})$ which have the property
\[
\|T^*x\| \leq \|T_x\| \quad \text{for every } x \in \mathcal{H}.
\]

A such operator is called hyponormal.

**Proposition 3.1.8** If $T$ is hyponormal operator on $\mathcal{H}$ then:

(i) $Tx = \lambda x$ implies $T^*x = \overline{\lambda}x$.

(ii) If $Tx = \lambda x$ and $Ty = \mu y$ for $\lambda \neq \mu$ then $\langle x, y \rangle = 0$.

**Proof.** The first property hold immediately since $T - \lambda I$ is also hyponormal operator. Now, for $\lambda \neq \mu$ we have
\[
\lambda \langle x, y \rangle = \langle \lambda x, y \rangle = \langle Tx, y \rangle = \langle x, T^*y \rangle = \langle x, \overline{\mu}y \rangle \quad \text{by the last property}
\]

Since $\lambda \neq \mu$ then $\langle x, y \rangle = 0$. 

**Remark 3.1.9** If $T \in \mathcal{L}(\mathcal{H})$ is a hyponormal operator then for every $\lambda \in \mathbb{C}$, $\ker(T - \lambda I)$ is a closed invariant subspace for $T$ and $T^*$.

**Theorem 3.1.10** For every hyponormal operator $T \in \mathcal{L}(\mathcal{H})$,
\[
\|T^n\| = \|T\|^n \quad \text{for every } n \geq 1.
\]

**Proof.** For $n = 1$, the equality is trivial; proceed by induction. For every vector $x \in \mathcal{H}$ we have
\[
\|T^n x\|^2 = \langle T^n x, T^n x \rangle = \langle T^*T^n x, T^{n-1} x \rangle \leq \|T^*T^n x\| \|T^{n-1} x\| \leq \|T^{n+1} x\| \|T^{n-1} x\| \leq \|T^{n+1}\| \|x\| \|T^{n-1}\| \|x\|^2.
\]

Since the vector $x$ is arbitrary, it follows that
\[
\|T^n\|^2 \leq \|T^{n+1}\| \|T^{n-1}\|.
\]
In view of the induction hypothesis \((\|T^k\| = \|T\|^k \text{ whenever } 1 \leq k \leq n)\), this can be rewritten as

\[\|T\|^{2\alpha} \leq \|T^{n+1}\| \cdot \|T\|^{n-1}.\]

Therefore,

\[\|T\|^{n+1} \leq \|T^{n+1}\|.\]

Since the reverse inequality is universal, the induction step is accomplished. \(\blacksquare\)

**Corollary 3.1.11** If \(T \in \mathcal{L}(\mathcal{H})\) is a subnormal operator and \(S\) is its minimal normal extension on a Hilbert space \(\mathcal{K}\) then \(\|T\| = \|S\|\).

**Proof.** Since every subnormal operator is hyponormal (see Theorem 3.1.7), it follows from the last theorem that \(r(T) = \|T\|\). On the other hand, we have \(\sigma(S) \subset \sigma(T)\) (see Theorem 3.1.6) then

\[\|S\| = r(S) \leq r(T) = \|T\|\]

The reverse inequality is trivial since \(T\) is the restriction of \(S\) on \(\mathcal{H}\). \(\blacksquare\)

**Theorem 3.1.12** Every hyponormal operator \(T\) on \(\mathcal{H}\) has the single valued extension property.

**Proof.** Let \(F\) be a vector valued analytic function such that

\[(T - \lambda I)F(\lambda) = 0 \text{ for every } \lambda \in D(F).\]

Then

\[TF(\lambda) = \lambda F(\lambda) \text{ for every } \lambda \in D(F).\]

Fix \(\lambda \in D(F)\), then for every \(\mu \neq \lambda \in D(F)\) we have by the last proposition

\[\langle F(\lambda), F(\mu) \rangle = 0.\]

By Pythagorean Theorem it follows that

\[\|F(\lambda) - F(\mu)\|^2 = \|F(\lambda)\|^2 + \|F(\mu)\|^2.\]

Letting \(\mu \to \lambda\) gives \(F(\lambda) = 0\). Therefore \(F\) is identically zero since \(\lambda\) is arbitrary element in \(D(F)\). This complete the proof. \(\blacksquare\)

**Theorem 3.1.13** (Stampfli and Radjabalipour) Every hyponormal operator \(T\) on \(\mathcal{H}\) has the Dunford’s Condition C (DCC).

For the proof see [?] and [?]. \(\blacksquare\)

### 3.2 Characterization of Subnormal Shifts

The weighted shift operators are interesting for solving a lot of problems in operator theory, they can be used for examples and counterexamples to illustrate many properties of operators. So, it is natural to ask which weighted shifts are normal, which are hyponormal, and which are subnormal? The first question was already answered that there is no positive weight sequence that makes a unilateral weighted shift normal and only the constant weight sequence that makes a bilateral weighted shift normal. The answer to the second question is also easy that is the hyponormal weighted shifts are characterized by monotonically increasing weight sequences. The answer to the third question is not easy. One formulation was offered by stampfli [?]; another formulation due to C. Berger is very different. It is elegant and easy to state [?], [?].
We begin by stating the **Spectral Theorem for Normal Operators** (see [?]? and [?]) which will be used throughout this section. A vector \( e \in \mathcal{H} \) is said to be star-cyclic for an operator \( T \in \mathcal{L}(\mathcal{H}) \) if \( \mathcal{H} \) is the smallest closed invariant subspace for \( T \) and for its adjoint \( T^* \) containing \( e \). The operator \( T \) is said to be star-cyclic if it has a star-cyclic vector. Note that a vector \( e \in \mathcal{H} \) is star-cyclic for a normal operator \( T \) if and only if the closed linear vector subspace generated by \( \{T^nT^{*m}e : n, m \in \mathbb{N}\} \) coincide with \( \mathcal{H} \). Every injective bilateral weighted shift is star-cyclic.

For a positive measure \( \mu \) with compact support subset \( K \) of \( \mathbb{C} \), the linear map \( N^\mu \) defined on \( L^2\) space \( L^2(K, \mu) \) by \( N^\mu(f)(z) = zf(z) \) is bounded operator on the Hilbert space \( L^2(K, \mu) \); it is star-cyclic with star-cyclic vector the constant function 1 since the set \( C(K) \) of all complex continuous functions on \( K \) is dense in \( L^2(K, \mu) \).

**Theorem 3.2.1 (Spectral Theorem for Normal Star-Cyclic Operators.)** Let \( e_0 \) be a star-cyclic vector for a normal operator \( T \in \mathcal{L}(\mathcal{H}) \). Then there is a probability measure \( \mu \) in the closed disc \( \overline{D} \) of the complex plane \( \mathbb{C} \) of radius \( \|T\| \) such that there is a unique isomorphism \( V : \mathcal{H} \rightarrow L^2(\overline{D}, \mu) \) with \( Ve_0 = 1 \) and \( VTV^{-1} = N^\mu \).

Let \( T \) be a weighted shift on a Hilbert space \( \mathcal{H} \) with a positive weight sequence \( (\omega_n)_n \), that is

\[
Tc_n = \omega_ne_{n+1} = \frac{\beta_{n+1}}{\beta_n}c_{n+1}
\]

where \( (e_n)_n \) is a orthonormal basis of \( \mathcal{H} \) and \( \beta \) is the following sequence given by:

\[
\beta_n = \begin{cases} 
\omega_0...\omega_{n-1} & \text{if } n > 0 \\
1 & \text{if } n = 0 \\
\frac{1}{\omega_{-n}...\omega_{-1}} & \text{if } n < 0
\end{cases}
\]

**Theorem 3.2.2 (Berger)** If \( T \) is unilateral weighted shift then \( T \) is subnormal if and only if there is a probability measure \( \mu \) on the closed interval \([0, \|T\|]\) such that for \( n \geq 1 \)

\[
\beta_n^2 = \int \! t^{2n}d\mu(t).
\]

**Proof.** Without loss of generality, we assume that \( \|T\| = 1 \).

Suppose first that \( T \) is a subnormal operator and let \( S \) be its minimal normal extension; then by the corollary 3.1.11 we have \( \|S\| = 1 \). Let \( M \) be a closed invariant subspace for \( S \) and its adjoint \( S^* \) containing \( e_0 \), then the restriction \( S|M \) of \( S \) on \( M \) is a normal operator; since \( S^n e_0 = T^n e_0 = \beta_n e_{n} \) then \( e_n \in M \) for every \( n \geq 0 \). Hence \( \mathcal{H} \subset M \); therefore by the minimality of \( S \) it follows that \( M = \mathcal{K} \) and so, \( e_0 \) is star-cyclic vector for the normal operator \( S \). The Theorem 3.2.1 shows that there is a probability measure \( \mu \) on the closed unit disc \( D \) such that there is a unique isomorphism \( V : \mathcal{H} \rightarrow L^2(\overline{D}, \mu) \) with \( Ve_0 = 1 \) and \( VSV^{-1} = N^\mu \). Let \( \nu \) be the probability measure in the closed unit interval \( I \) defined for each Borel subset \( E \) of \( I \) by \( \nu(E) = \mu(S^{-1}(E)) \) where \( S : \overline{D} \rightarrow I \) is the mapping given by \( S(z) = |z| \). It follows that \( \int \! f(t)d\nu(t) = \int \! f \circ S(z)d\mu(z) \) for every \( L^1(I, \nu) \)-function \( f \). In particular, for every \( n \geq 0 \), we have

\[
\int \! t^{2n}d\nu(t) = \int \! |z|^{2n}d\mu(z) = \int \! |N^\mu_n(1)|^2d\mu(z)
\]

\[
= \|N^\mu_n(1)\|^2 = \|N^\mu_nV e_0\|^2 = \|VS^n e_0\|^2
\]

\[
= \|S^n e_0\|^2 = \|T^n e_0\|^2
\]

\[
= \beta_n^2
\]
and so, the proof of the necessity of the Theorem is complete.

Now, let \( \mu \) be a probability measure in the closed unit interval \( I \) such that \( \beta_n^2 = \int t^{2n} d\mu(t) \) for every \( n \geq 0 \). We consider the probability measure \( \nu = d\mu(r) \frac{dr}{2\pi} \) on the closed unit disc \( \overline{D} \). The operator \( N_\nu \) of “multiplication by \( z \)” on the Hilbert space \( L^2(\overline{D}, \nu) \) is a bounded normal operator. Let \( H^2(\nu) \) be the closed subspace spanned by \( \{z^n : n \geq 0\} \). Then \( H^2(\nu) \) is invariant closed subspace for \( N_\nu \) and \( N_\nu \) restricted to \( H^2(\nu) \) is a subnormal operator \( R \in \mathcal{L}(H^2(\nu)) \). On the other hand, for every \( n, m \in \mathbb{N} \) we have

\[
(z^n, z^m) = \int_{\overline{D}} r^{n+m} e^{i(n-m)\theta} d\nu(r e^{i\theta}) = \frac{1}{2\pi} \int_0^{2\pi} r^{n+m} dr \int_0^{2\pi} e^{i(n-m)\theta} d\theta.
\]

Hence \( (v_n)_{n \in \mathbb{N}} \) is orthonormal basis for \( H^2(\nu) \) where \( v_n = \frac{z^n}{\|v_n\|} = \frac{z^n}{\beta_n} \) for every \( n \in \mathbb{N} \). And so, for every \( n \in \mathbb{N} \) we have

\[
Rv_n = N_\nu\left(\frac{z^n}{\beta_n}\right) = \frac{1}{\beta_n} z^{n+1} = \omega_n v_{n+1}.
\]

It follows that \( T \) is a subnormal operator. \( \square \)

**Remark 3.2.3** If \( \mu \) is a probability measure on the interval \([0, \|T\|]\) such that for \( n \geq 1 \)

\[
\beta_n^2 = \int t^{2n} d\mu(t),
\]

then using the Cauchy-Schwartz Inequality we get

\[
\beta_n^2 \leq \beta_{n-1}\beta_{n+1} \quad \text{for every } n \geq 1.
\]

This condition is equivalent that the weight sequence \( (\omega_n)_n \) is increasing that is the weighted shift \( T \) is hyponormal.

We close this section with the characterization of subnormal bilateral shifts.

**Theorem 3.2.4** If \( T \) is bilateral weighted shift then \( T \) is subnormal if and only if there is a probability measure \( \mu \) on the closed interval \([0, \|S\|]\) such that for \( n \geq 1 \) the functions \( t^n, t^{-n} \in L^1(\mu) \) and

\[
\frac{\beta_n^2}{\beta_{-n}^2} = \int t^{2n} d\mu(t) \quad \text{and} \quad \frac{1}{\beta_{-n}^2} = \int t^{-2n} d\mu(t).
\]
Chapter 4

Bounded Point Evaluations for Cyclic Operators

4.1 Bounded Point Evaluations

Throughout this section, \( T \) will be a cyclic bounded operator on a Hilbert space \( \mathcal{H} \) with cyclic vector \( x \). A complex number \( \lambda \in \mathbb{C} \) is said to be a bounded point evaluation of \( T \) if there is a constant \( M > 0 \) such that

\[
\left| p(\lambda) \right| \leq M \| p(T)x \|
\]

for every complex polynomial \( p \). The set of all bounded point evaluations of \( T \) will be denoted by \( B(T) \). Note that it follows from the Riesz Representation Theorem (see [?] and [?]) that \( \lambda \in B(T) \) if and only if there is a unique vector denoted \( k_\lambda \in \mathcal{H} \) such that \( p(\lambda) = \langle p(T)x , k_\lambda \rangle \) for every complex polynomial \( p \).

**Proposition 4.1.1** \( B(T) = \Gamma(T) \) the compression spectrum of \( T \).

**Proof.** Let \( \lambda \in B(T) \), then \( x \) can not be in the closure of the range, \( \overline{\text{Im}(T - \lambda I)} \), of \( (T - \lambda I) \) otherwise \( T^nx \) is in \( \text{Im}(T - \lambda I) \) for every \( n \geq 0 \) since \( T \) and \( (T - \lambda I) \) commutes; hence \( \text{Im}(T - \lambda I) = \mathcal{H} \) and \( k_\lambda = 0 \) which is impossible since \( 1 = \langle x , k_\lambda \rangle \). Conversely, let \( \lambda \in \Gamma(T) \) then \( x \) is not in the closure of the range of \( T - \lambda I \) since \( x \) is also cyclic vector for \( T - \lambda I \). Let \( y \in \mathcal{H} \) be the orthogonal projection of \( x \) onto the orthogonal complement of the range of \( T - \lambda I \) and set \( k_\lambda = \frac{1}{\langle x , y \rangle} y \). Since for every complex polynomial \( p \) there is a complex polynomial \( q \) such that \( p(z) = (z - \lambda)q(z) + p(\lambda) \), then

\[
\langle p(T)x , k_\lambda \rangle = \langle [(T - \lambda I)q(T) + p(\lambda)]x , k_\lambda \rangle
\]

\[
= \langle (T - \lambda I)q(T)x , k_\lambda \rangle + \langle p(\lambda)x , k_\lambda \rangle
\]

\[
= 0 + p(\lambda)\langle x , k_\lambda \rangle = p(\lambda).
\]

**Proposition 4.1.2** Let \( \lambda \in \mathbb{C} \), the following statements are equivalent:

(i) \( \lambda \in B(T) \).

(ii) \( \ker ((T - \lambda)^*) \) is one dimensional.

(iii) \( \overline{x} \in \sigma_p(T^*) \).

**Proof.** It is clear that (ii) implies (iii) and also (iii) and (i) are equivalent since \( \Gamma(T) = \sigma_p(T^*) \) (see [?]).
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(i) implies (ii). It is clear that $k_{\lambda} \in \ker ((T - \lambda)^*)$. Let $u$ be a non zero element of $\ker ((T - \lambda)^*)$, so it suffices to prove that there is $\alpha \neq 0 \in \mathbb{C}$ such that $u = \alpha k_{\lambda}$. Let $\alpha = \frac{1}{\langle u, k_{\lambda} \rangle}$. Since for every polynomial $p$ there is a polynomial $q$ such that $p(z) = (z - \lambda)q(z) + p(\lambda)$, then

$$
\langle p(T)u, \alpha k_{\lambda} \rangle = \langle [(T - \lambda I)q(T) + p(\lambda)]u, \alpha k_{\lambda} \rangle \\
= \langle (T - \lambda I)q(T)u, \alpha k_{\lambda} \rangle + \langle p(\lambda)u, \alpha k_{\lambda} \rangle \\
= \langle q(T)u, (T - \lambda I)^*\alpha k_{\lambda} \rangle + p(\lambda)(u, \alpha k_{\lambda}) \\
= p(\lambda).
$$

Hence, $u = \alpha k_{\lambda}$. Thus proves that $\ker ((T - \lambda)^*)$ is one dimensional. This finishes the proof. \(\blacksquare\)

**Definition 4.1.3** A subset $O$ of $B(T)$ which is open in $\mathbb{C}$ is said to be an analytic set for $T$ if for every $z \in \mathcal{H}$ the function $z \rightarrow \|k\|$ is analytic on $O$. The largest analytic set for $T$ is will denoted by $B_a(T)$ and every point of its will be called analytic bounded point evaluation for $T$.

**Lemma 4.1.4** A subset $O$ of $B(T)$ which is open in $\mathbb{C}$ is an analytic set for $T$ if and only if the function $\lambda \rightarrow \|k\|$ is bounded on compact subsets of $O$.

**Proof.** First suppose that $O$ is an analytic set for $T$ and $K$ is a compact subset of $O$. For every $z \in \mathcal{H}$ the function $z \rightarrow \|k\|$ is analytic on $O$, in particular, $\sup_{\lambda \in K} \|z, k\| < +\infty$. So, it follows from the Uniform Boundedness Principle that

$$
\sup_{\lambda \in K} \|k\| < +\infty.
$$

Conversely, suppose that the function $\lambda \rightarrow \|k\|$ is bounded on compact subsets of $O$. Let $z \in \mathcal{H}$, then there is a sequence of polynomial $(p_n)_n$ such that $\lim_{n \to +\infty} \|p_n(T)x - z\| = 0$. And so, for every compact subset $K$ of $O$ it follows by using Cauchy-Schwartz inequality that

$$
\sup_{\lambda \in K} |p_n(\lambda) - \tilde{z}(\lambda)| \leq \sup_{\lambda \in K} \|k\||p_n(T)x - z\|.
$$

Hence, $\tilde{z}$ is analytic function on $O$ (see $\|\|$). \(\blacksquare\)

**Proposition 4.1.5 (Williams)** $\Gamma(T) \setminus \sigma_{ap}(T) \subset B_a(T)$.

**Proof.** Since $\sigma(T) = \Gamma(T) \cup \sigma_{ap}(T)$ then $\sigma(T) \setminus \sigma_{ap}(T) = \Gamma(T) \setminus \sigma_{ap}(T)$. On the other hand, $\sigma(T) \setminus \sigma_{ap}(T) = \text{int}(\sigma(T)) \setminus \sigma_{ap}(T)$ since the boundary of the spectrum of $T$ is contained in $\sigma_{ap}(T)$. So, the set $O = \Gamma(T) \setminus \sigma_{ap}(T)$ is a subset of $B(T)$ which is open in $\mathbb{C}$. Let now $\lambda \in O$ then there is a positive constant $C$ such that $\|z\| \leq C\|(T - \lambda I)z\|$ for every $z \in \mathcal{H}$. So, for $\mu \in \mathbb{C}$ such that $|\mu - \lambda| \leq \frac{1}{2C}$ we have

$$
\|z\| \leq C\|(T - \lambda I)z\| = C\|(T - \mu I)z + (\mu - \lambda)z\| \\
\leq C\|(T - \mu I)z\| + |\mu - \lambda|\|z\| \\
\leq C\|(T - \mu I)z\| + \frac{1}{2}\|z\|.
$$

Therefore $\|z\| \leq 2C\|(T - \mu I)z\|$ for each $z \in \mathcal{H}$ and for each complex number $\mu$ satisfying $|\mu - \lambda| \leq \frac{1}{2C}$. In particular, for every polynomial $p$ we have

$$(*) \quad |p(\lambda)| \leq \|p(T)x\||k\| \leq 2C\|(T - \mu I)p(T)x\||k\| \quad \text{for} \quad \mu \in \mathbb{C}, \ |\mu - \lambda| \leq \frac{1}{2C}.$$
Now, let $\mu \in \mathbb{C}$ then for every polynomial $p$ there is a polynomial $q$ such that $p(t) = (t - \mu)q(t) + p(\mu)$. Hence for $\mu \in \mathbb{C}$, $|\mu - \lambda| \leq \frac{1}{2C}$

$$|p(\mu)| \leq |p(\lambda)| + |\lambda - \mu||q(\lambda)|$$

$$\leq \|p(T)x\||k_{\lambda}| + 2C|\lambda - \mu||(T - \mu I)q(T)x||k_{\lambda}||$$

by (1)

$$\leq \|p(T)x\||k_{\lambda}| + 2C|\lambda - \mu||p(T)x - p(\mu)x||k_{\lambda}||$$

$$\leq \|p(T)x\||k_{\lambda}| + 2C|\lambda - \mu|[\|p(T)x\| + |p(\mu)||x||]||k_{\lambda}||.$$

So, if in addition, $\mu \in \mathbb{C}$, $|\mu - \lambda| \leq M = \min \left(\frac{1}{2C}, \frac{1}{4C\|k_{\lambda}\||x||}\right)$ then

$$|p(\mu)| \leq \|p(T)x\||k_{\lambda}| + \frac{1}{2\|x\||p(T)x\| + \frac{1}{2}|p(\mu)|.}

It follows that $|p(\mu)| \leq \|p(T)x\||k_{\lambda}| + \frac{1}{2\|x\||p(T)x\| + \frac{1}{2}|p(\mu)|$ for each $\mu \in \mathbb{C}$, $|\mu - \lambda| \leq M$ and for every polynomial $p$. And so, $\|k_{\mu}\| \leq 2\|k_{\lambda}\| + \frac{1}{\|x\||}$ for every $\mu \in \mathbb{C}$, $|\lambda - \mu| \leq M$ since the set $\{p(T)x : p \text{ is a polynomial}\}$ is dense in $\mathcal{H}$ and every vector $k_{\mu}$ define a bounded linear functional on $\mathcal{H}$ of norm $\|k_{\mu}\|$. Therefore, the function $\mu \mapsto \|k_{\mu}\|$ is bounded on compact subsets of $O$. Hence $O$ is an analytic set for $T$, therefore $\Gamma(T) \setminus \sigma_{ap}(T) \subset B_{a}(T)$.

Tavan Trent \cite{Tavan Trent} proved that the converse of this proposition holds for the operator $S_{\mu}$ of multiplication by $z$ on $H^{2}(\mu)$, the closure of the polynomials in the $L^{2}(\mu)$ space where $\mu$ is a positive finite compactly supported Borel measure.

**Proposition 4.1.6 (Tavan Trent)** Let $\mu$ be a positive finite compactly supported Borel measure. Then

$$\Gamma(S_{\mu}) \setminus \sigma_{ap}(S_{\mu}) = B_{a}(S_{\mu}).$$

**Proof.** We first observe that the constant function 1 is a cyclic vector for $S_{\mu}$ and $p(S_{\mu})1 = p$ for every polynomial $p$. Since $B_{a}(S_{\mu}) \subset B(S_{\mu}) = \Gamma(S_{\mu})$ (see Proposition 4.1.1) then it suffices to prove that $B_{a}(S_{\mu}) \cap \sigma_{ap}(S_{\mu}) = \emptyset$. Now, let $\lambda \in B_{a}(S_{\mu})$ then there is $r > 0$ such that $B_{a}(S_{\mu})$ contained a closed disc $D$ centered at $\lambda$ and of radius $r$ since $B_{a}(S_{\mu})$ is open set. So,

$$\sup \{\|k_{\xi}\| : \xi \in D\} = C < \infty.$$

Let $p$ be a polynomial then for $\gamma$ in the boundary of $D$ i.e: $|\gamma - \lambda| = r$, we have $|(\gamma - \lambda)p(\gamma)| \leq \|(S_{\mu} - \lambda)p\||k_{\gamma}||$ since $\gamma$ is a bounded point evaluation for $S_{\mu}$. Hence,

$$|p(\gamma)| \leq \frac{C}{r} \|(S_{\mu} - \lambda)p\| \quad \text{for every } \gamma \in \mathbb{C}, |\gamma - \lambda| = r.$$

Since the polynomial $p$ is analytic function on $D$ then by the Maximum Modulus Principle,

$$|p(\gamma)| \leq \frac{C}{r} \|(S_{\mu} - \lambda)p\| \quad \text{for every } \gamma \in D.$$
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Hence,

$$\|p\|^2 = \int_{\mathbb{D}} |p|^2 d\mu + \int_{\mathbb{C}\setminus\mathbb{D}} |p|^2 d\mu$$

$$\leq \left( \frac{C}{r} \right)^2 \|(S_\mu - \lambda)p\|^2 \mu(\mathbb{D}) + \int_{\mathbb{C}\setminus\mathbb{D}} \frac{1}{|z - \lambda|^2} |z - \lambda p(z)|^2 d\mu(z)$$

$$\leq \left( \frac{C}{r} \right)^2 \|(S_\mu - \lambda)p\|^2 \mu(\mathbb{D}) + \frac{1}{r^2} \int_{\mathbb{C}\setminus\mathbb{D}} \left| (S_\mu - \lambda)p(z) \right|^2 d\mu(z)$$

$$\leq \left( \frac{C}{r} \right)^2 \|(S_\mu - \lambda)p\|^2 \mu(\mathbb{D}) + \frac{1}{r^2} \|(S_\mu - \lambda)p\|^2.$$
that two quasisimilar cyclic operators have the same bounded point evaluations and the same analytic bounded point evaluations.

**Theorem 4.1.8** Let $H_1$ and $H_2$ be two Hilbert spaces. Two quasisimilar cyclic operators $T \in \mathcal{L}(H_1)$ and $S \in \mathcal{L}(H_2)$ have the same bounded point evaluations and the same analytic bounded point evaluations, i.e: $B(T) = B(S)$ and $B_a(T) = B_a(S)$.

**Proof.** Suppose that there exist two bounded transformations $X : H_1 \rightarrow H_2$ and $Y : H_2 \rightarrow H_1$ having trivial kernels and dense ranges such that $XR = SX$ and $RY = YS$. Since $X$ has dense range then $0 \notin \Gamma(X)$ so, $X^*$ is injective since $\Gamma(X) = \sigma_p(X^*)$. If $\lambda \in \mathbb{C}$ is eigenvalue of $S^*$ with the corresponding eigenvector $x \in H_2$ then $R^*(X^*x) = \lambda X^*x$, and so, $\lambda \in \sigma_p(R^*)$ since $X^*x \neq 0$. Hence, $\sigma_p(S^*) \subseteq \sigma_p(R^*)$. By symmetry, $\sigma_p(S^*) = \sigma_p(R^*)$, thus $B(R) = B(S)$. Suppose that $v \in H_2$ is a cyclic vector for $S$ then $u = Yv$ is a cyclic vector for $R$. Let $\lambda$ be in $B(R) = B(S)$ there are vectors $k_{\lambda} \in H_1$ and $h_{\lambda} \in H_2$ such that $p(\lambda) = \langle p(R)u , k_{\lambda} \rangle = \langle p(S)v , h_{\lambda} \rangle$ for each polynomial. On the other hand $p(R)Y = Yp(S)$ since $RY = YS$ for every polynomial $p$. Ans so,

$$\langle p(S)v , Y^*k_{\lambda} \rangle = \langle Yp(S)v , k_{\lambda} \rangle = \langle p(R)Yv , k_{\lambda} \rangle = \langle p(R)u , k_{\lambda} \rangle = p(\lambda) = \langle p(S)v , h_{\lambda} \rangle$$

for every polynomial $p$ and $\lambda \in B(R) = B(S)$. Hence, $h_{\lambda} = Y^*k_{\lambda}$ for every $\lambda \in B(R) = B(S)$. Since $\lambda \mapsto \|k_{\lambda}\|$ is bounded on compact subsets of $B_a(R)$ then $\lambda \mapsto \|h_{\lambda}\|$ is bounded on compact subsets of $B_a(S)$. It follows from Lemma 4.1.4 that $B_a(R) \subseteq B_a(S)$. By symmetry, the desired result holds i.e: $B_a(R) = B_a(S)$.

### 4.2 Analytic Bounded Point Evaluations for Unilateral Weighted Shift

In [3], A.L. Shields represented a weighted shift operator as ordinary shift operator (that is, as "multiplication by $z$") on a Hilbert space of formal power series (in the unilateral case) or formal Laurent series (in the bilateral case). Thus he defined the concept of bounded point evaluations of a weighted shift to examine which power series and Laurent series represent analytic functions. In fact, this concept of bounded point evaluations for injective unilateral weighted shift coincide with the one defined by L.R. Williams in [3].

We now describe the set of bounded point evaluations and the set of analytic bounded point evaluations for an arbitrary injective unilateral weighted shift. Let $S$ be a unilateral weighted shift on a Hilbert space $H$ with a positive weight sequence $(\omega_n)_{n \geq 0}$, that is

$$Sc_n = \omega_ne_{n+1} = \frac{\beta_{n+1}}{\beta_n} e_{n+1}$$

where $(e_n)_{n \geq 0}$ is a orthonormal basis of $H$ and $\beta$ is the following sequence given by:

$$\beta_n = \begin{cases} 
\omega_0 \cdots \omega_{n-1} & \text{if } n > 0 \\
1 & \text{if } n = 0
\end{cases}$$
The unilateral weighted shift \( S \) is cyclic of cyclic vector \( e_0 \). It follows from Corollary 1.3.6 and Theorem 4.1.8 that the set of the bounded point evaluations and the set of analytic bounded point evaluations for the unilateral weighted shift \( S \) have a circular symmetry about the origin. Also, it follows from Theorem 1.3.12 and Proposition 4.1.1 that \( B(S) = \{0\} \) if \( r_2(S) = 0 \) otherwise
\[
\{ \lambda \in \mathbb{C} : |\lambda| < r_2(S) \} \subset B(S) \subset \{ \lambda \in \mathbb{C} : |\lambda| \leq r_2(S) \},
\]
where \( r_2(S) = \liminf \frac{1}{n} \beta_n^\frac{1}{n} \).

**Theorem 4.2.1** If \( r_2(S) > 0 \) then \( B_a(S) = \{ \lambda \in \mathbb{C} : |\lambda| < r_2(S) \} \).

**Proof.** Let \( \lambda \in B(S) \) then there is \( k_\lambda = \sum_{n \geq 0} \hat{k}_\lambda(n)e_n \in \mathcal{H} \) such that \( p(\lambda) = \langle p(S)e_0, k_\lambda \rangle \) for every polynomial \( p \). In particular, for every \( n \geq 0 \) we have
\[
\hat{k}_\lambda(n) = \langle e_n, k_\lambda \rangle = \langle \frac{1}{\beta_n} S^n e_0, k_\lambda \rangle = \frac{\lambda^n}{\beta_n}.
\]
Now, let \( h = \sum_{n \geq 0} a_n e_n \) then \( \tilde{h}(\lambda) = \langle h, k_\lambda \rangle = \sum_{n \geq 0} \frac{a_n}{\beta_n} \lambda^n \). In fact this series is absolutely convergent since \( \sum_{n \geq 0} |a_n| e_n \) is also in \( \mathcal{H} \) and \( |\lambda| \) is also in \( B(S) \). Hence, \( \tilde{h} \) is analytic in the interior of \( B(S) \) which is exactly the disc \( \{ \lambda \in \mathbb{C} : |\lambda| < r_2(S) \} \). This proves the theorem. ■

For the weighted shift \( S \), the spectrum \( \sigma(S) \) is known to be the disk \( \{ \lambda \in \mathbb{C} : |\lambda| \leq r(S) \} \), (see Theorem 1.3.8) and the approximate point spectrum \( \sigma_{ap}(S) \) is known to be the annulus \( \{ \lambda \in \mathbb{C} : r_1(S) \leq |\lambda| \leq r(S) \} \) (see Theorem 1.3.10). Therefore, \( \Gamma(S) \setminus \sigma_{ap}(S) = \sigma(S) \setminus \sigma_{ap}(S) = \{ \lambda \in \mathbb{C} : |\lambda| < r_1(S) \} \). And so, \( \Gamma(S) \setminus \sigma_{ap}(S) \subsetneq B_a(S) \) if and only if \( r_1(S) < r_2(S) \). Hence, a negative answer to the question 1 (see [?]) can be given by a unilateral weighted shifts \( S \) for which \( r_1(S) < r_2(S) \). Let us consider an example of a such weighted shift. For \( s \in \mathbb{N} \) there are unique \( n, k \in \mathbb{N} \) such that \( s = n! + k \) with \( 0 < k < (n+1)! - n! - 1 \). We set
\[
\beta_s = \beta_{n! + k} = e^k.
\]
And so, for every \( k \in \mathbb{N} \), we have
\[
\frac{\beta_{k+1}}{\beta_k} = \begin{cases} e & \text{if } n! \leq k < (n+1)! - 1 \\ \frac{1}{e^{(n+1)! - n! - 1}} & \text{if } k = (n+1)! - 1 \end{cases}
\]
Hence, the unilateral weighted shift \( S \) corresponding to the weight \( (\frac{\beta_{k+1}}{\beta_k})_{n \geq 0} \) is bounded (see Proposition 1.3.1). For every \( n \geq 2 \), set \( k_n = n! - n \). Clearly, we have \( (n-1)! \leq k_n < n! \) and \( \beta_{k_n} = e^{(n!-(n-1)!-n)} \). And so,
\[
\inf_{k \geq 0} \frac{\beta_{n+k}}{\beta_k} \leq \frac{\beta_{n+k_n}}{\beta_{k_n}} = \frac{1}{e^{(n!-(n-1)!-n)}}.
\]
Hence,
\[
r_1(S) = \lim_{n \to \infty} \left[ \inf_{k \geq 0} \frac{\beta_{n+k}}{\beta_k} \right]^\frac{1}{n} = 0.
\]
On the other hand, it is clear that
\[ r_2(S) = \liminf_{n \to \infty} \left[ \beta_n \right]^\frac{1}{n} = 1. \]
Therefore, \( B_a(S) = \{ \lambda \in \mathbb{C} : |\lambda| < 1 \} \) and \( \Gamma(S) \setminus \sigma_{ap}(S) = \emptyset. \)
For the unilateral weighted shift \( S \) the set of its analytic bounded point evaluations is exactly the interior of the set of its bounded point evaluations. This suggests the following question, which was first posed by J.B. Conway in [?] page 65:

**Question 2**: Is always the interior of \( B(T) \) coincide with \( B_a(T) \) for an arbitrary cyclic operator \( T \in \mathcal{L}(\mathcal{H}) \)?

### 4.3 The Local Spectra Through Bounded Point Evaluations

In what follows, \( \mathcal{H} \) denote a Hilbert space. In this section we will show that for a cyclic operator \( T \in \mathcal{L}(\mathcal{H}) \) with Dunford’s Condition C and without point spectrum, the local spectra of \( T \) at \( x \) is equal to the spectrum of \( T \) for each \( x \) in a dense subset of \( \mathcal{H} \).

**Lemma 4.3.1** Suppose that \( \mathcal{H} = \mathcal{H}_1 \oplus \mathcal{H}_2 \) where \( \mathcal{H}_1 \) and \( \mathcal{H}_2 \) are two Hilbert spaces such that \( \mathcal{H}_2 \) is finite dimensional. If \( T \in \mathcal{L}(\mathcal{H}) \) has the single valued extension property and \( \mathcal{H}_1 \) is invariant subspace for \( T \) then \( A = T|_{\mathcal{H}_1} \) has the single valued extension property and \( \sigma_A(x) = \sigma_T(x) \) for every \( x \in \mathcal{H}_1 \).

**Proof.** It is clear that \( A \) has the single valued extension property and \( \sigma_T(x) \subset \sigma_A(x) \) for every \( x \in \mathcal{H}_1 \). Conversely, let \( x \in \mathcal{H}_1 \) then \( \widetilde{x} = F_1 + F_2 \) on \( \rho_T(x) \) where \( F_1 = P_1 \widetilde{x} \) and \( P_i : \mathcal{H} \rightarrow \mathcal{H}_i \) are the orthogonal projections \( i = 1, 2 \). And so, for every \( \lambda \in \rho_T(x) \)

\[
x = (T - \lambda I)\widetilde{x}(\lambda) \\
= (T - \lambda I)F_1(\lambda) + (T - \lambda I)F_2(\lambda) \\
= (A - \lambda I)F_1(\lambda) + P_1(T - \lambda I)F_2(\lambda) + P_2(T - \lambda I)F_2(\lambda) \\
= (A - \lambda I)F_1(\lambda) + P_1TF_2(\lambda) + (P_2T - \lambda I)F_2(\lambda) \\
= \left[ (A - \lambda I)F_1(\lambda) + P_1TF_2(\lambda) \right] + (P_2T - \lambda I)F_2(\lambda).
\]

It follows that \( (P_2T - \lambda I)F_2(\lambda) = 0 \) for every \( \lambda \in \rho_T(x) \). Since \( \sigma(P_2T) \) is a finite set then \( F_2 \) is identically zero function. Hence,

\[
x = (T - \lambda I)\widetilde{x}(\lambda) = (A - \lambda I)F_1(\lambda) \quad \text{for all } \lambda \in \rho_T(x).
\]

Therefore, \( \sigma_A(x) \subset \sigma_T(x) \). The proof is complete.

**Lemma 4.3.2** Suppose that \( \mathcal{H} = \mathcal{H}_1 \oplus \mathcal{H}_2 \) where \( \mathcal{H}_1 \) and \( \mathcal{H}_2 \) are two Hilbert spaces such that \( \mathcal{H}_2 \) is finite dimensional. If \( \mathcal{H}_1 \) is an invariant subspace for an operator \( T \in \mathcal{L}(\mathcal{H}) \) which satisfies DCC then the operator \( A = T|_{\mathcal{H}_1} \) satisfies the DCC.

**Proof.** It follows from the last Lemma that for every closed subset \( F \) of \( \mathcal{C} \), \( \mathcal{H}_{1A}(F) = \mathcal{H}_2(F) \cap \mathcal{H}_1 \). So, the desired result holds.
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Lemma 4.3.3 Suppose that \( \mathcal{H} = \mathcal{H}_1 \oplus \mathcal{H}_2 \) where \( \mathcal{H}_1 \) and \( \mathcal{H}_2 \) are two Hilbert spaces such that \( \mathcal{H}_2 \) is finite dimensional. If \( T \in \mathcal{L}(\mathcal{H}) \) is an operator without point spectrum such that \( \mathcal{H}_1 \) is invariant subspace for \( T \) then \( \sigma(T) = \sigma(A) \) where \( A = T|_{\mathcal{H}_1} \).

Proof. We first observe that there are two operators \( B : \mathcal{H}_2 \rightarrow \mathcal{H}_1 \) and \( C : \mathcal{H}_2 \rightarrow \mathcal{H}_2 \) such that for every \( x = x_1 \oplus x_2 \in \mathcal{H} \), \( Tx = [Ax_1 + Bx_2] + Cx_2 \). Now, suppose that \( A \) is invertible in \( \mathcal{L}(\mathcal{H}_1) \) then \( TH_1 = H_1 \) and \( H_1 \cap TH_2 = \{0\} \). Let \( x_2 \in \mathcal{H}_2 \) such that \( Cx_2 = 0 \) then \( Tx_2 = Bx_2 \in H_1 \cap TH_2 \), hence \( C \) is injective and so, by the finite-dimensionality, \( C \) is invertible. Therefore the operator \( T \) is invertible with inverse given by:

\[
T^{-1}x = [A^{-1} x_1 - A^{-1} B C^{-1} x_2] + C^{-1} x_2 \quad \text{for every } x = x_1 \oplus x_2 \in \mathcal{H}.
\]

Thus \( \sigma(T) \subset \sigma(A) \). The converse follows from Theorem 2.1.2 and Lemma 4.3.1.

Theorem 4.3.4 Let \( T \in \mathcal{L}(\mathcal{H}) \) be a cyclic operator with cyclic vector \( x \in \mathcal{H} \) and let \( S \in \mathcal{L}(\mathcal{H}) \) be an operator commutes with \( T \) such that \( \ker(S^*) \) is finite dimensional. If \( T \) satisfies DCC and \( \sigma_p(T) = \emptyset \) then \( \sigma_T(Sx) = \sigma(A) \) where \( A = T \).

Proof. Let \( \mathcal{H}_1 \) be the closed linear subspace generated by \( \{T^n Sx : n \geq 0\} \) then it is clear that \( \mathcal{H}_1 \) is invariant subspace for \( T \). Since \( TS = ST \) and \( x \) is a cyclic vector for \( T \), then \( \mathcal{H}_1 = \overline{\text{Im}(S)} \).

Thus \( \mathcal{H} = \mathcal{H}_1 \oplus \mathcal{H}_2 \) where \( \mathcal{H}_2 = \ker(S^*) \). It follows from the Lemma 4.3.1 that \( \sigma_T(Sx) = \sigma_A(Sx) \) where \( A = T|_{\mathcal{H}_1} \). Since \( A \) is a cyclic operator with cyclic vector \( Sx \) and satisfies DCC (Lemma 4.3.2) then it follows from Proposition 2.2.2 that \( \sigma_A(Sx) = \sigma(A) \). Since \( \sigma_p(T) = \emptyset \), it follows from Lemma 4.3.3 that \( \sigma(T) = \sigma(A) \). And so, \( \sigma(T) = \sigma(A) = \sigma_A(Sx) = \sigma_T(Sx) \). The proof is complete.

Remark 4.3.5 Let \( T \in \mathcal{L}(\mathcal{H}) \) be a cyclic operator with cyclic vector \( x \in \mathcal{H} \) and satisfies DCC such that \( \sigma_p(T) = \emptyset \). Since for every non zero polynomial \( p \), there exists \( a, a_1, \ldots, a_n \in \mathbb{C} \) such that \( p(T)^* = a(T^* - a_1 I) \cdots (T^* - a_n) \). Then it follows from Proposition 4.1.2 that \( \ker(p(T)^*) \) is finite dimensional. Hence, \( \sigma_T(p(T)x) = \sigma(T) \) for every non zero polynomial \( p \) (see Theorem 4.3.4). Therefore, \( \sigma_T(y) = \sigma(T) \) holds for all non zero \( y \) in a dense subset of \( \mathcal{H} \). L.R. Williams proved in Theorem 2.5 of [?] that if \( T \) is a non normal hyponormal (unilateral or bilateral) weighted shift operator then \( \sigma_T(x) = \sigma(T) \) for every non zero element \( x \in \mathcal{H} \).

We state and we give a simple proof of Theorem 2.5 of [?] using the fact that a non zero analytic function has isolate zeroes.

Theorem 4.3.6 Let \( T \) be a non normal hyponormal weighted shift on \( \mathcal{H} \). Then for every a non zero element \( x \in \mathcal{H} \), \( \sigma_T(x) = \sigma(T) \).

Proof. First suppose that \( T \) is a non normal hyponormal unilateral weighted shift, then \( r(T) = r_1(T) = r_2(T) = \|T\| > 0 \). Let now \( x \in \mathcal{H} \) such that there is \( \lambda \in \sigma(T) \setminus \sigma_T(x) \). So, there is vector valued analytic function \( F \) on an open neighbourhood \( V \) of \( \lambda \) such that

\[
(T - \mu I)F(\mu) = x \quad \text{for every } \mu \in V.
\]

Since \( \emptyset \neq V \cap \text{int}(\sigma(T)) \subset B_\alpha(T) \), then for every \( \mu \in V \cap \text{int}(\sigma(T)) \) we have,

\[
\tilde{x}(\mu) = \langle x, k_\mu \rangle = \langle (T - \mu I)F(\mu), k_\mu \rangle = \langle F(\mu), (T - \mu I)^* k_\mu \rangle = 0.
\]

Hence, \( \tilde{x} \equiv 0 \). And so, \( x = 0 \).

The case of a non normal hyponormal bilateral weighted shift is similar using the bounded point evaluations of the sense of A.L. Shields [?].
Remark 4.3.7 Using the same proof of this Theorem, one can see that for every injective unilateral weighted shift $T$, $B_0(T) \subset \sigma_T(x)$ for every non zero element $x \in \mathcal{H}$. The same for every bilateral weighted shift.

Question 3. For which hyponormal operators $T$ we have $\sigma_T(h) = \sigma(T)$ $\forall h \neq 0$?
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