# ÉCOLE POLYTECHNIQUE FÉDÉRALE DE LAUSANNE 

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## Problem 1.

(a) (1) Additive noise: This model is one in which the impairment to the signal is an addition of a noise, a stochastic process independent of the signal.
(2) Distortion: A distortion is the alteration of the original shape (or other characteristic) of a waveform or other form of information or representation. In general a noise-free system can be characterized by a transfer function, such that the output $y(t)$ can be written as a function of the input $x(t)$ as

$$
y(t)=F(x(t))
$$

This form is memoryless distortion.
(3) Interference: Interference is anything which alters, modifies, or disrupts a signal as it travels along a channel between a source and a receiver. The term typically refers to the addition of unwanted signals from other transmitters to a useful signal.
(4) Fading: Fading is the random attenuation that a carrier-modulated signal experiences over certain propagation media. The fading may vary with time, geographical position and/or radio frequency, and is often modeled as a random process. A fading channel is a communication channel that experiences fading. In wireless systems, fading may either be due to multipath propagation, referred to as multipath induced fading, or due to shadowing from obstacles affecting the wave propagation, sometimes referred to as shadow fading.
(5) Doppler shift: It is the change that occurs in the frequency of a wave for an observer with a non-zero relative velocity to the source.
All channels are energy and bandwidth limited, but note that some channels are more restrictive on one area than the other. In general channels may be classified as bandwidth or energy limited according to whether they permit transmission at high spectral efficiency or not. There is no sharp line, but usually we could take spectral efficiency $2(\mathrm{~b} / \mathrm{s}) / \mathrm{Hz}$ as the boundary, corresponding to the highest spectral efficiency that can be achieved with binary transmission.
(i) In this region, the capacity (achievable spectral efficiency) increases linearly with $S N R$. In these kind of channel, we will be able to use binary coding and modulation.
(ii) In this region, the capacity (achievable spectral efficiency) increases logarithmic with $S N R$. In these kind of channels, we must use non-binary (multilevel) modulation.
(b) (1) Wireline channels (non-optical): Additive noise, Interference and Bandwidth limited
(2) Optical channels: Low attenuation, high bandwidth. No Interference
(3) Deep space satellite channels: Additive noise, Energy limited
(4) Microwave links: Line of sight links so no Interference, distortion and not much fading, High power
(5) Underwater acoustic channels:Additive background noise, distortion and doppler shift, fading and band limited
(6) Cellphone wireless communication channels: Additive noise, fading, doppler shift, interference, energy limited, bandwidth limited

Problem 2. $p_{V W}(v, w)$.
(a)

$$
\begin{align*}
E(V+W) & =\iint(v+w) p_{V W}(v, w) d v d w  \tag{1}\\
& =\iint\left(v p_{V W}(v, w)+w p_{V W}(v, w)\right) d v d w  \tag{2}\\
& =\iint v p_{V W}(v, w) d v d w+\iint w p_{V W}(v, w) d v d w  \tag{3}\\
& =\int v \int p_{V W}(v, w) d w d v+\int w \int p_{V W}(v, w) d v d w  \tag{4}\\
& =\int v p_{V}(v) d v+\int w p_{W}(w) d w  \tag{5}\\
& =E(V)+E(W) \tag{6}
\end{align*}
$$

(b)

$$
\begin{align*}
E(V \cdot W) & =\iint(v \cdot w) p_{V W}(v, w) d v d w  \tag{7}\\
& =\iint(v \cdot w) p_{V}(v) p_{W}(w) d v d w  \tag{8}\\
& =\int v p_{V}(v) d v \cdot \int w p_{W}(w) d w  \tag{9}\\
& =E(V) \cdot E(W) \tag{10}
\end{align*}
$$

(c) Assume $V=W$ and $\operatorname{Pr}(V=1)=\operatorname{Pr}(V=-1)=\frac{1}{2}$. We compute $E(V)=E(W)=0$ and $E(V W)=1$, so $E(V W) \neq E(V) E(W)$
Now suppose $(V, W)$ takes values of $(1,1),(1,-1),(-1,1),(-1,-1),(0,0)$ with equal probability $\frac{1}{5}$. Because $\operatorname{Pr}(W=0 \mid V=1)=0 \neq \frac{1}{5}=\operatorname{Pr}(W=0), V$ and $W$ are not independent. We compute $E(V)=E(W)=0$ and $E(V W)=0$, so $E(V W)=$ $E(V) E(W)$
(d) Assume that $V$ and $W$ are independent and let $\sigma_{V}^{2}$ and $\sigma_{W}^{2}$ be the variances of $V$ and $W$, respectively. Show that the variance of $V+W$ is given by $\sigma_{V+W}^{2}=\sigma_{V}^{2}+\sigma_{W}^{2}$.

$$
\begin{align*}
\sigma_{V+W}^{2} & =E\left((V+W)^{2}\right)-E((V+W))^{2}  \tag{11}\\
& =E\left(V^{2}\right)+E\left(W^{2}\right)+2 E(V W)-(E(V)+E(W))^{2}  \tag{12}\\
& =E\left(V^{2}\right)+E\left(W^{2}\right)+2 E(V) E(W)-E(V)^{2}-E(W)^{2}-2 E(V) E(W \chi 13) \\
& =E\left(V^{2}\right)-E(V)^{2}+E\left(W^{2}\right)-E(W)^{2}  \tag{14}\\
& =\sigma_{V}^{2}+\sigma_{W}^{2} \tag{15}
\end{align*}
$$

## Problem 3.

(a)

$$
\begin{align*}
\sum_{n>0} \operatorname{Pr}(N \geq n) & =\sum_{n=1}^{\infty} \sum_{m=n}^{\infty} \operatorname{Pr}(N=m)  \tag{16}\\
& =\sum_{m=1}^{\infty} \sum_{n=1}^{m} \operatorname{Pr}(N=m)  \tag{17}\\
& =\sum_{m=1}^{\infty} m \operatorname{Pr}(N=m)  \tag{18}\\
& =E(N) \tag{19}
\end{align*}
$$

(b)

$$
\begin{align*}
\int_{0}^{\infty} \operatorname{Pr}(x \geq a) d a & =\int_{0}^{\infty} \int_{a}^{\infty} f_{x}(t) d t d a  \tag{20}\\
& =\int_{0}^{\infty} \int_{0}^{t} f_{x}(t) d a d t  \tag{21}\\
& =\int_{0}^{\infty} t f_{x}(t) d t  \tag{22}\\
& =E(X) \tag{23}
\end{align*}
$$

(c) The main point is to note that $G(t)=P(X \geq t)$ is a non-increasing function of $t$. So for any fixed value of $a>0$, the rectangle between point $(0,0)$ and ( $a, G(a)$ ) lies below the function $G(t)$. In conclusion, it follows from the discussion above that

$$
a G(a) \leq \int_{0}^{a} G(a) d t \leq \int_{0}^{a} G(t) d t \leq \int_{0}^{\infty} G(t) d t
$$

which means

$$
a \operatorname{Pr}(X \geq a) \leq E(X)
$$

(d) Assume

$$
X=(Y-E(Y))^{2} \quad X \geq 0
$$

Using part (c), we have

$$
a \operatorname{Pr}(X \geq a) \leq E(X)
$$

Therefore, one could conclude that

$$
a \operatorname{Pr}\left((Y-E(Y))^{2} \geq a\right) \leq E\left((Y-E(Y))^{2}\right)
$$

Setting $b=\sqrt{a}$, we have

$$
\operatorname{Pr}(|Y-E(Y)| \geq b)=\operatorname{Pr}\left((Y-E(Y))^{2} \geq b^{2}\right) \leq \frac{E\left((Y-E(Y))^{2}\right)}{b^{2}}=\frac{\sigma_{Y}^{2}}{b^{2}}
$$

## Problem 4.

(a) $\operatorname{Pr}\left(X_{1} \leq X_{2}\right)=\frac{1}{2}$. We know because of independence we have, $f_{X_{1}, X_{2}}\left(x_{1}, x_{2}\right)=$ $f_{X_{1}}\left(x_{1}\right) f_{X_{2}}\left(x_{2}\right)$, and we want to find the probability of $x_{1}$ being minimum of two. This event partitions the probability space into two equal sub-sets, the other one is $x_{2}$ being the minimum of the two. The only problem is the boundary line $x_{1}=x_{2}$, which we assume is a part of first sub-set, but because $f_{x}$ is a continuous random variable the line $x_{1}=x_{2}$ has zero probability mass and because $f_{X_{1}}\left(x_{1}\right) f_{X_{2}}\left(x_{2}\right)$ is symmetric with respect to the line $x_{1}=x_{2}$, we conclude that the event $\min \left(x_{1}, x_{2}\right)=x_{1}$ partitions the whole probability space into two equally probable regions.
(b) $\operatorname{Pr}\left(X_{1} \leq X_{2} ; X_{1} \leq X_{3}\right)=\frac{1}{3}$; We follow the exact same argument as the part (a), this time the probability space is partitioned into three equally probable sub-sets, in each of sub-sets one of the three random variable is minimum.
(c) Similar to last parts, we can show that

$$
\operatorname{Pr}\left(X_{1} \leq X_{2} ; X_{1} \leq X_{3} ; \ldots ; X_{1} \leq X_{n-1} ; X_{1} \leq X_{n}\right)=\frac{1}{n}
$$

and

$$
\operatorname{Pr}\left(X_{1} \leq X_{2} ; X_{1} \leq X_{3} ; \ldots ; X_{1} \leq X_{n-1}\right)=\frac{1}{n-1}
$$

We know

$$
\begin{align*}
\operatorname{Pr}(N=n)= & \operatorname{Pr}\left(X_{1} \leq X_{2} ; X_{1} \leq X_{3} ; \ldots ; X_{1} \leq X_{n-1} ; X_{1}>X_{n}\right)  \tag{24}\\
= & \operatorname{Pr}\left(X_{1} \leq X_{2} ; X_{1} \leq X_{3} ; \ldots ; X_{1} \leq X_{n-1}\right) \\
& \quad-\operatorname{Pr}\left(X_{1} \leq X_{2} ; X_{1} \leq X_{3} ; \ldots ; X_{1} \leq X_{n-1} ; X_{1} \leq X_{n}\right)  \tag{25}\\
= & \frac{1}{n-1}-\frac{1}{n}=\frac{1}{n^{2}-n} . \quad n>1 \tag{26}
\end{align*}
$$

Using properties of telescopic series, we conclude

$$
\begin{align*}
\operatorname{Pr}(N \geq n) & =\sum_{m=n}^{\infty} \operatorname{Pr}(N=m)  \tag{27}\\
& =\frac{1}{n-1}-\frac{1}{n}+\frac{1}{n}-\frac{1}{n+1}+\ldots  \tag{28}\\
& =\frac{1}{n-1} . \quad n \geq 2 \tag{29}
\end{align*}
$$

(d) We use part (a) of Problem 3.

$$
E(N)=\sum_{n>0} \operatorname{Pr}(N \geq n)=\sum_{n>1} \frac{1}{n-1} \rightarrow \infty
$$

(We know that series $\frac{1}{n}$ is divergent.)
(e) The symmetry of the $f_{X_{1}}\left(x_{1}\right) f_{X_{2}}\left(x_{2}\right)$ still holds because of independence but in the discrete case it is possible to put some probability mass on line $x_{1}=x_{2}$. Therefore in the discrete case the event $x_{1} \leq x_{2}$ does not partition the whole probability space into two equally probable sub-spaces. The same as before we can conclude that $\operatorname{Pr}\left(X_{1}<\right.$ $\left.X_{2}\right)=\operatorname{Pr}\left(X_{2}<X_{1}\right)$. We know $\operatorname{Pr}\left(X_{1}<X_{2}\right)+\operatorname{Pr}\left(X_{1}=X_{2}\right)+\operatorname{Pr}\left(X_{2}<X_{1}\right)=1$. From these two we conclude that $\operatorname{Pr}\left(X_{1} \leq X_{2}\right) \geq \frac{1}{2}$. Similarly we conclude that

$$
\operatorname{Pr}\left(X_{1} \leq X_{2} ; X_{1} \leq X_{3} ; \ldots ; X_{1} \leq X_{n-1} ; X_{1} \leq X_{n}\right) \geq \frac{1}{n}
$$

Following the steps in part (d), we can show that

$$
E(N) \geq \sum_{n>1} \frac{1}{n-1} \rightarrow \infty
$$

Problem 5. Let's consider the case where $n=2$ first, we have

$$
P(Z=0)=P\left(X_{1} \oplus X_{2}=0\right)=P\left(X_{1}=0, X_{2}=0\right)+P\left(X_{1}=1, X_{2}=1\right)=\frac{1}{2}
$$

in which we used independence of $X_{1}$ and $X_{2}$.
By induction, one could easily show that for arbitrary $n$, we have

$$
P(Z=0)=\frac{1}{2} .
$$

(a)

$$
\begin{align*}
P\left(Z=z \mid X_{1}=x_{1}\right) & =P\left(X_{1} \oplus X_{2} \oplus \cdots \oplus X_{n}=z \mid X_{1}=x_{1}\right)  \tag{30}\\
& =P\left(X_{2} \oplus \cdots \oplus X_{n}=z \oplus x_{1} \mid X_{1}=x_{1}\right)  \tag{31}\\
& =P\left(X_{2} \oplus \cdots \oplus X_{n}=z \oplus x_{1}\right)  \tag{32}\\
& =\frac{1}{2}=P(Z=z) \tag{33}
\end{align*}
$$

in (32) we used that $X_{i}$ 's are independent. We conclude that $Z$ is independent of $X_{1}$
(b)

$$
\begin{align*}
P\left(Z=z \mid X_{1}, \ldots, X_{n-1}=x_{1}, \ldots, x_{n-1}\right) & =  \tag{34}\\
P\left(X_{1} \oplus X_{2} \oplus \cdots \oplus X_{n}=z \mid X_{1}, \ldots, X_{n-1}=x_{1}, \ldots, x_{n-1}\right) & =  \tag{35}\\
P\left(X_{n}=z \oplus x_{1} \oplus \cdots \oplus x_{n-1} \mid X_{1}, \ldots, X_{n-1}=x_{1}, \ldots, x_{n-1}\right) & =  \tag{36}\\
P\left(X_{n}=z \oplus x_{1} \oplus \cdots \oplus x_{n-1}\right) & =  \tag{37}\\
& =\frac{1}{2}  \tag{38}\\
& =P(Z=z) \tag{39}
\end{align*}
$$

in (37) we used that $X_{i}$ 's are independent. We conclude that $Z$ is independent of $X_{1}, \ldots, X_{n-1}$.
(c) No, $Z$ is a deterministic function of $X_{1}, \ldots, X_{n}$, which means

$$
P\left(Z=z \mid X_{1}, \ldots, X_{n}=x_{1}, \ldots, x_{n}\right)
$$

is either 0 or 1 depending on the values of $x_{1}, \ldots, x_{n}$ and $z$.
(d) Suppose $\operatorname{Pr}\left(X_{i}=1\right)=\frac{3}{4}$, we have

$$
P(Z=0)=P\left(X_{1} \oplus X_{2}=0\right)=P\left(X_{1}=0, X_{2}=0\right)+P\left(X_{1}=1, X_{2}=1\right)=\frac{9+1}{16}=\frac{5}{8} .
$$

but

$$
\begin{align*}
P\left(Z=0 \mid X_{1}=0\right) & =P\left(X_{1} \oplus X_{2}=0 \mid X_{1}=0\right)  \tag{40}\\
& =P\left(X_{2}=0 \mid X_{1}=0\right)  \tag{41}\\
& =\frac{1}{4} \neq \frac{5}{8}=P(Z=0), \tag{42}
\end{align*}
$$

in which we used that $X_{1}$ and $X_{2}$ are independent. We conclude that $Z$ is not independent of $X_{1}$.

Problem 6. (1)

$$
\begin{align*}
V\left(\pi_{0}\right) & =\pi_{0} \operatorname{Pr}\left(\text { error } \mid H_{0}\right)+\left(1-\pi_{0}\right) \operatorname{Pr}\left(\text { error } \mid H_{1}\right)  \tag{43}\\
& =\pi_{0} \int_{R_{H_{0}}} \operatorname{Pr}\left(y \mid H_{0}\right) d y+\left(1-\pi_{0}\right) \int_{R \backslash R_{H_{0}}} \operatorname{Pr}\left(y \mid H_{1}\right) d y, \tag{44}
\end{align*}
$$

where we define

$$
R_{H_{0}}=\left\{y \in R \mid \pi_{0} \operatorname{Pr}\left(y \mid H_{0}\right) \geq\left(1-\pi_{0}\right) \operatorname{Pr}\left(y \mid H_{1}\right)\right\} .
$$

(2) We define

$$
R_{\pi_{0}}=\left\{y \in R \mid \pi_{0} \operatorname{Pr}\left(y \mid H_{0}\right) \geq\left(1-\pi_{0}\right) \operatorname{Pr}\left(y \mid H_{1}\right)\right\} .
$$

and

$$
R_{\pi_{0}^{\prime}}=\left\{y \in R \mid \pi_{0}^{\prime} \operatorname{Pr}\left(y \mid H_{0}\right) \geq\left(1-\pi_{0}^{\prime}\right) \operatorname{Pr}\left(y \mid H_{1}\right)\right\} .
$$

We know that these regions are optimal decision regions for $V\left(\pi_{0}\right)$ and $V\left(\pi_{0}^{\prime}\right)$, respectively, which means they minimize $V$ for the given prior among all other decision regions. Now, we define
$R_{\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}}=\left\{y \in R \mid\left(\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}\right) \operatorname{Pr}\left(y \mid H_{0}\right) \geq\left(1-\left(\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}\right)\right) \operatorname{Pr}\left(y \mid H_{1}\right)\right\}$, and we conclude that

$$
\begin{align*}
V\left(\pi_{0}\right) & =\pi_{0} \int_{R_{\pi_{0}}} \operatorname{Pr}\left(y \mid H_{0}\right) d y+\left(1-\pi_{0}\right) \int_{R \backslash R_{\pi_{0}}} \operatorname{Pr}\left(y \mid H_{1}\right) d y  \tag{45}\\
& \leq \pi_{0} \int_{R_{\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}}} \operatorname{Pr}\left(y \mid H_{0}\right) d y+\left(1-\pi_{0}\right) \int_{R \backslash R_{\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}}} \operatorname{Pr}\left(y \mid H_{1}\right) d y,(4 \tag{46}
\end{align*}
$$

and

$$
\begin{align*}
V\left(\pi_{0}^{\prime}\right) & =\pi_{0}^{\prime} \int_{R_{\pi_{0}^{\prime}}} \operatorname{Pr}\left(y \mid H_{0}\right) d y+\left(1-\pi_{0}^{\prime}\right) \int_{R \backslash R_{\pi_{0}^{\prime}}} \operatorname{Pr}\left(y \mid H_{1}\right) d y  \tag{47}\\
& \leq \pi_{0}^{\prime} \int_{R_{\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}}} \operatorname{Pr}\left(y \mid H_{0}\right) d y+\left(1-\pi_{0}^{\prime}\right) \int_{R \backslash R_{\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}}} \operatorname{Pr}\left(y \mid H_{1}\right) d y,( \tag{48}
\end{align*}
$$

In (46) and (48), we used the fact that $R_{\pi_{0}}$ and $R_{\pi_{0}^{\prime}}$ minimize $V$ for the given prior among all other decision regions. Multiplying (46) by $\lambda$ and (48) by $1-\lambda$ and adding them up, we will have

$$
\begin{align*}
\lambda V\left(\pi_{0}\right)+(1-\lambda) V\left(\pi_{0}^{\prime}\right) \leq & \left(\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}\right) \int_{R_{\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}}} \operatorname{Pr}\left(y \mid H_{0}\right) d y \\
& +\left(1-\left(\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}\right)\right) \int_{R \backslash R_{\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}}} \operatorname{Pr}\left(y \mid H_{1}\right) d y \\
= & V\left(\lambda \pi_{0}+(1-\lambda) \pi_{0}^{\prime}\right) \tag{49}
\end{align*}
$$

(3) Concavity implies that $V$ is continuous on $[0,1]$ and the minimum is on the boundary and for a given prior $\pi_{\text {max }}$, most likely in the interior of $] 0,1[$ (unless $V$ is a linear function), $V$ is maximized over all $\pi \in[0,1]$.

Problem 7. We define

$$
C\left(x_{i}\right)=2 \sigma^{2} \log \operatorname{Pr}\left(x_{i}\right)
$$

It is easy to show that for the optimal decision maker (MAP) in Gaussian noise, the detector finds $x_{i}$ so that

$$
\left\langle x_{i}, x_{i}\right\rangle-2\left\langle y, x_{i}\right\rangle-C\left(x_{i}\right)
$$

is minimized.
We know the following for any $j \neq i$

$$
\begin{align*}
& \left\langle x_{i}, x_{i}\right\rangle-2\left\langle y_{1}, x_{i}\right\rangle-C\left(x_{i}\right) \leq\left\langle x_{j}, x_{j}\right\rangle-2\left\langle y_{1}, x_{j}\right\rangle-C\left(x_{j}\right)  \tag{50}\\
& \left\langle x_{i}, x_{i}\right\rangle-2\left\langle y_{2}, x_{i}\right\rangle-C\left(x_{i}\right) \leq\left\langle x_{j}, x_{j}\right\rangle-2\left\langle y_{2}, x_{j}\right\rangle-C\left(x_{j}\right) . \tag{51}
\end{align*}
$$

Now let us consider the following,

$$
\begin{aligned}
\left\langle x_{i}, x_{i}\right\rangle-2\left\langle\alpha y_{1}+(1-\alpha) y_{2}, x_{i}\right\rangle-C\left(x_{i}\right)= & \left\langle x_{i}, x_{i}\right\rangle-2 \alpha\left\langle y_{1}, x_{i}\right\rangle \\
& -2(1-\alpha)\left\langle y_{2}, x_{i}\right\rangle-C\left(x_{i}\right) \\
= & \alpha\left[\left\langle x_{i}, x_{i}\right\rangle-2\left\langle y_{1}, x_{i}\right\rangle-C\left(x_{i}\right)\right]+ \\
& (1-\alpha)\left[\left\langle x_{i}, x_{i}\right\rangle-2\left\langle y_{2}, x_{i}\right\rangle-C\left(x_{i}\right)\right] \\
\leq & \alpha\left[\left\langle x_{j}, x_{j}\right\rangle-2\left\langle y_{1}, x_{j}\right\rangle-C\left(x_{j}\right)\right]+ \\
& (1-\alpha)\left[\left\langle x_{j}, x_{j}\right\rangle-2\left\langle y_{2}, x_{j}\right\rangle-C\left(x_{j}\right)\right] .
\end{aligned}
$$

In the last step we used 50 and 51 . We conclude

$$
\left\langle x_{i}, x_{i}\right\rangle-2\left\langle\alpha y_{1}+(1-\alpha) y_{2}, x_{i}\right\rangle-C\left(x_{i}\right) \leq\left\langle x_{j}, x_{j}\right\rangle-2\left\langle\alpha y_{1}+(1-\alpha) y_{2}, x_{j}\right\rangle-C\left(x_{j}\right)
$$

for all $j \neq i$. Therefore, the decoder decodes $\alpha y_{1}+(1-\alpha) y_{2}$ as $x_{i}$.
Problem 8. Based on what we learned in the class, we find the optimal decision region as follows,

$$
R_{-3}=\left\{y \in \mathbb{R} \mid \pi_{-3} \operatorname{Pr}(y \mid-3)>\pi_{3} \operatorname{Pr}(y \mid 3)\right\},
$$

which simplifies to

$$
\begin{align*}
\frac{2}{3} \operatorname{Pr}(z=y+3) & >\frac{1}{3} \operatorname{Pr}(z=y-3)  \tag{52}\\
\frac{2}{3} \frac{1}{\pi\left(1+(y+3)^{2}\right)} & >\frac{1}{3} \frac{1}{\pi\left(1+(y-3)^{2}\right)}  \tag{53}\\
2 y^{2}-12 y+20 & >y^{2}+6 y+10  \tag{54}\\
y^{2}-18 y+10 & >0 . \tag{55}
\end{align*}
$$

Therefore, we have

$$
R_{-3}=\{y \in \mathbb{R} \mid y \geq 17.4 \text { or } y \leq 0.57\}
$$

and

$$
R_{3}=\{y \in \mathbb{R} \mid 0.57 \leq y \leq 17.4\} .
$$

Note that $R_{-3}$ is not convex!
The error probability is

$$
\begin{align*}
P_{e} & =\frac{1}{3} \int_{-\infty}^{-2.43} \frac{1}{\pi\left(1+z^{2}\right)} d z+\frac{1}{3} \int_{14.4}^{\infty} \frac{1}{\pi\left(1+z^{2}\right)} d z+\frac{2}{3} \int_{3.57}^{20.4} \frac{1}{\pi\left(1+z^{2}\right)} d z  \tag{56}\\
& =\frac{1}{3}\left(1-\int_{-2.43}^{14.4} \frac{1}{\pi\left(1+z^{2}\right)} d z\right)+\frac{2}{3} \int_{3.57}^{20.4} \frac{1}{\pi\left(1+z^{2}\right)} d z . \tag{57}
\end{align*}
$$

For Gaussian noise we have

$$
R_{-3}=\left\{y \in \mathbb{R} \mid \pi_{-3} \operatorname{Pr}(y \mid-3)>\pi_{3} \operatorname{Pr}(y \mid 3)\right\},
$$

which simplifies to

$$
\begin{align*}
\frac{2}{3} \operatorname{Pr}(z=y+3) & >\frac{1}{3} \operatorname{Pr}(z=y-3)  \tag{58}\\
\frac{2}{3} \frac{1}{\sqrt{2 \pi}} e^{-\frac{(y+3)^{2}}{2}} & >\frac{1}{3} \frac{1}{\sqrt{2 \pi}} e^{-\frac{(y-3)^{2}}{2}}  \tag{59}\\
\ln (2) & >6 y  \tag{60}\\
\frac{\ln (2)}{6} & >y . \tag{61}
\end{align*}
$$

Therefore, we have

$$
R_{-3}=\left\{y \in \mathbb{R} \left\lvert\, y<\frac{\ln (2)}{6}=0.1155\right.\right\}
$$

and

$$
R_{3}=\{y \in \mathbb{R} \mid y>0.1155\} .
$$

The error probability is

$$
P_{e}=\frac{1}{3} Q(3-0.1155)+\frac{2}{3} Q(3+0.1155)=0.0013 .
$$

