## Tentamen [s]: EI1102/EI1100 Elkretsanalys, 2014-08-29 kl 08-13

Omtentan gäller en nerlagd kurs, vilken hade sista omgången HT13. Nästa omtentan planeras att ske samtidig med tentan i kursen EI1120, VT15.
Hjälpmedel: Ett A4-ark med studentens anteckningar (båda sidor). Dessutom, pennor!
Svar får anges på svenska eller engelska. En kort ordlista finns på sista sidan.
Tentan har 3 tal i del A (15p), och 3 tal i del B (15p).
Godkänt vid $\geq 25 \%$ på del A och del B individuellt, och $\geq 50 \%$ på delarna A och B tillsammans. Betyget räknas sedan från summan av A och B. Se KursPM:et angående betygsgränser, rättningsnormer och överklagande.

Läs varje tal noggrant innan du försöker svara.
Senare deltal kan vara betydligt svårare än de första deltalen.
Satsa inte för mycket tid på bara en uppgift om du fastnar: ta hänsyn till poängvärden på uppgifterna, och att man måste både delar av tentan.

Tänk på att använda återstående tid till att kolla igenom varje svar: man kan göra dimensionsanalys, rimlighetsbedömning (t.ex. "är det rätt att $y$ går ner medan $x$ går ner?"), och lösning genom en alternativ metod. Lösningar ska förenklas om inte annat är specificerat.

Examinator: Nathaniel Taylor

## Del A. Likström och Transienter.

1) $[5 \mathrm{p}]$

Kända: $R_{1}, R_{2}, R_{3}, R_{4}, U, I_{1}, I_{2}$.
a) [1p] Bestäm effekten levererat till $R_{2}$.
b) $[1 \mathrm{p}]$ Bestäm den markerade spänningen $u_{x}$.
c) $[1 \mathrm{p}]$ Bestäm effekten levererat till $R_{3}$.
d) $[2 \mathrm{p}]$ Bestäm effekten levererat från källan $U$.

2) $[5 p]$

Kända: $R_{1}, R_{2}, R_{3}, R_{4}, U, I, g, h$.
Använd nodanalys för att skriva ekvationer som går att lösa för de okända potentialerna $v_{1}, v_{2}, v_{3}, v_{4}$.

Du måste inte lösa ekvationerna, och måste inte skriva om dem i förenklad eller matris form. Det finns flera möjliga svar (alla med samma lösning).

3) $[5 \mathrm{p}]$

Kända: $U, C, L_{1}, L_{2}, R_{1}, R_{2}, R_{3}$.
Kretsen är i jämviktsläge vid tiden $t=0^{-}$, d.v.s. just innan $t=0$ när brytaren stängs.
a) $[2 \mathrm{p}]$ Bestäm $i\left(0^{-}\right)$.
b) $[2 \mathrm{p}]$ Bestäm $u\left(0^{+}\right)$.
c) $[1 \mathrm{p}]$ Bestäm funktionen $i(t)$, för $t>0$.


## Del B. Växelström

4) $[5 \mathrm{p}]$

Kända: $R_{\mathrm{i}}, R, L, R_{\mathrm{o}}$.
Spänningen $u_{\mathrm{i}}$ på polerna till vänster orsaker spänningen $u_{o}$ mellan polerna till höger.
a) $[2 \mathrm{p}]$ Bestäm kretsens nätverksfunktion, $\frac{u_{\mathrm{o}}(\omega)}{u_{\mathrm{i}}(\omega)}$.
b) [3p] Nätverksfunktionen från deltal 'a)' borde kunna skrivas i formen,

$$
H(\omega)=1+\frac{\mathrm{j} \omega / \omega_{\mathrm{z}}}{1+\mathrm{j} \omega / \omega_{\mathrm{p}}}
$$

(Du måste inte bevisa det, men du kanske vill dubbelkolla del 'a'!) Skissa ett Bode amplituddiagram av $H(\omega)$, med viktiga punkter och lutningar markerade. Anta att $\omega_{\mathrm{p}} \gg \omega_{\mathrm{z}}$. Är du osäker om hantering av ' $1+$ ' termen, så kan du lämna ut den utan stort avdrag.

5) $[5 \mathrm{p}]$

Kända: $I, C, R, L, \omega$.
$I$ är en komplex storhet som beskriver en växelströmskälla av vinkelfrekvensen $\omega$.
a) [2p] Bestäm Norton-ekvivalenten av kretsen, mellan polerna 'a' och 'b'. Visa ekvivalanten med ett diagram. Var inte förvånad om inte alla kända komponentvärdena behövs.

b) [3p] En impedans kopplas till kretsen, mellan polerna 'a' och 'b'. Rita ett diagram av hur impedansen kan skapas med två komponenter (välj mellan motstånd, spole, kondensator) på ett sätt som maximerar överföringen av aktiveffekt till impedansen. Det finns flera än ett sätt. Uttryck de två komponentvärden i de kända variablerna.
6) $[5 p]$

Kända: $\hat{U}, \hat{I}, \omega, \theta, R, L, C$.
Tidsvariabeln är $t$.
a) $[3 \mathrm{p}]$ Bestäm $i_{x}(t)$. Anta att lösningen kan göras med växelströmsanalys, d.v.s. att det inte finns transienter kvar. Observera att källorna har olika frekvenser.
b) $[2 \mathrm{p}]$ Bestäm den genomsnittliga effekten
 ("aktiveffekt" i växelströms terminologi) i $R$.

Ordlista över mindre självklara översättningar: resistor motstånd, capacitor kondensator, inductor (coil) spole, current ström, voltage spänning, power effekt, active power aktiveffekt source källa, switch brytare, terminal pol, angular(radian) frequency vinkelfrekvens,

## Solutions (EI1102, Omtenta 2014-08-29)

Q1
a) Resistor $R_{2}$ is in series with a current source $I_{2}$, which determines its current. Power in a resistor can be calculated from current and resistance, without having to consider the direction of current. The power into the resistor is $I_{2}^{2} R_{2}$.
b) The marked voltage $u_{x}$ across $R_{3}$ can be found by voltage division across the series-connected pair $R_{3}$ and $R_{4}$. This is convenient because this pair is directly connected to the voltage source $U$; it might not be very obvious, but you can follow the nodes round to see that it's true. Therefore, $u_{x}=U \frac{R_{3}}{R_{3}+R_{4}}$. One should be careful to check that the definition directions of $U$ and $u_{x}$ don't make a negative sign necessary!
c) The voltage across $R_{3}$ is $u_{x}$, already known from the above.

The power dissipated in a resistor can be found from the voltage across it and its resistance, i.e. $\frac{u_{x}^{2}}{R_{3}}$.
This power is $\left(U \frac{R_{3}}{R_{3}+R_{4}}\right)^{2} / R_{3}$, which simplifies to $\frac{U^{2} R_{3}}{\left(R_{3}+R_{4}\right)^{2}}$.
d) The voltage of a voltage source is known by definition. By finding the current coming out of the + terminal of source $U$, and multiplying by this voltage $U$, the power out of the source is found. So, what is the current? Consider the node at the + terminal of the voltage source:

* the current from this node into the current source $I_{1}$ is $-I_{1}$;
* the current into the resistor $R_{3}$ is $\frac{U}{R_{3}+R_{4}}$ as found in part 'b)';
and so, by KCL, the current out from the voltage source's + terminal is $\frac{U}{R_{3}+R_{4}}-I_{1}$.
The power out from the voltage source is then $\frac{U^{2}}{R_{3}+R_{4}}-U I_{1}$.


## Q2

Some previous solutions of homeworks or exams have shown two methods of systematically writing an equation system from a circuit. Here, we will use just the 'simple' direct method, which results in more equations and variables, but has advantages of less work spent on making simplifications, and more direct correspondence between the equations and the circuit diagram. Other methods, e.g. by considering 'supernodes', produce a smaller equation system, which can be desirable for hand-calculation.
Define currents (unknown variables) in the voltage sources: let's call them $I_{\alpha}$ into the + pole of source $U$, and $I_{\beta}$ into the + pole of source $h i_{x}$.
Now do KCL on the outgoing current at all the nodes other than the ground node:
$\begin{array}{ll}\mathrm{KCL}(1): & I_{\beta}-I_{\alpha}+\frac{v_{1}-v_{4}}{R_{2}}=0 \\ \mathrm{KCL}(2): & -I_{\beta}+\frac{v_{2}}{R_{1}}=0 \\ \mathrm{KCL}(3): & g u_{y}+\frac{v_{3}-v_{4}}{R_{3}}+\frac{v_{3}-v_{4}}{R_{4}}-I=0 \\ \mathrm{KCL}(4): & \frac{v_{4}-v_{1}}{R_{2}}+\frac{v_{4}-v_{3}}{R_{3}}+\frac{v_{4}-v_{3}}{R_{4}}+I=0\end{array}$
Then we add equations describing the relations imposed between node potentials by the voltage sources. These two equations compensate for the two unknown currents $I_{\alpha}$ and $I_{\beta}$ in the voltage sources.
$v_{1}=-U$
$v_{1}-v_{2}=h i_{x}$
Finally, the two controlling variables of the dependent sources need to be defined in terms of already-introduced variables:
$u_{y}=v_{3}-v_{4}$
$i_{x}=\frac{v_{3}-v_{4}}{R_{4}}-I$

Counting up the unknowns and equations, we see $v_{1 \cdots 4}, I_{\alpha}, I_{\beta}, i_{x}$ and $u_{y}$ ( 8 unknowns), and 8 equations also. Having used a suitable systematic method, we can be confident that a solvable circuit will have given a solvable equation system.

## Q3

a) The source is dc (a constant value) and we are told to assume that there is equilibrium at this time $t=0^{-}$. We can therefore find voltages and currents by replacing inductors with short circuits, and capacitors with open circuits. This leaves a circuit with just the three resistors and voltage source.
Solving for current $i$, which is the current through resistor $R_{1}$, we can first find the current from the source, which is $\frac{U}{\frac{R_{1} R_{2}}{R_{1}+R_{2}}+R_{3}}$, then use current division between $R_{1}$ and $R_{2}$, giving $i\left(0^{-}\right)=\frac{U}{\frac{R_{1} R_{2}}{R_{1}+R_{2}}+R_{3}} \cdot \frac{R_{2}}{R_{1}+R_{2}}$.

Simplifying this, $i\left(0^{-}\right)=\frac{U R_{2}}{R_{1} R_{2}+R_{3}\left(R_{1}+R_{2}\right)}$.
b) The voltage $u(t)$ is defined across a capacitor: it is therefore a continuous variable, as it relates to the stored energy, so $u\left(0^{-}\right)=u\left(0^{+}\right)$. The task of finding $u\left(0^{+}\right)$is thus acheived by finding the equilibrium level, $u\left(0^{-}\right)$before the switch closed.

By the same process as part 'a)', the circuit can be reduced to three resistors and the voltage source, and $u$ is seen to be the same as the voltage across the parallel resistors $R_{1}$ or $R_{2}$.
By voltage division, this is $u\left(0^{-}\right)=U \frac{\frac{R_{1} R_{2}}{R_{1}+R_{2}}}{\frac{R_{1}+R_{2}}{R_{1}+R_{3}}}$.
Simplifying, and using $u\left(0^{-}\right)=u\left(0^{+}\right)$, the answer is $u\left(0^{+}\right)=U \frac{1}{1+R_{3} \frac{R_{1}+R_{2}}{R_{1} R_{2}}}$ or $u\left(0^{+}\right)=\frac{U R_{1} R_{2}}{R_{1} R_{2}+R_{2} R_{3}+R_{3} R_{1}}$.
c) The current $i$ is the current in an inductor, so it is a continuous variable. Its initial state, $i\left(0^{+}\right)$, is therefore the same as $i\left(0^{-}\right)$which was found in part 'a)'. The switch short-circuits $R_{2}$ and the right part of the circuit, so the current $i(t>0)$ runs in a loop of just $L_{1}$ and $R_{1}$. This loop has a time-constant of $\tau=L_{1} / R_{1}$. The final state of $i(t)$ as $t \rightarrow \infty$, must be zero, as there is no source but there is a resistor that dissipates energy.
The requested time-function is therefore $i(t)=i\left(0^{-}\right) \exp (-t / \tau)$, which has the desired property of going from the initial value to the final value, in an exponential shape with time-constant $\tau$, as expected for a first-order circuit.
This answer should ideally be expressed in the given quantities: $i(t)=\frac{U R_{2}}{R_{1} R_{2}+R_{3}\left(R_{1}+R_{2}\right)} \exp -t \frac{R_{1}}{L_{1}} \quad(t>0)$.

## Q4

Yes, I could have mentioned that the opamp is ideal. This time I thought I wouldn't, as this is an assumption we make for all the components. But previous exams always said this, so I understand it was a bit worrying if something was different in this case.
a) The function $u_{\mathrm{o}} / u_{\mathrm{i}}$ can be found by assuming that the opamp's feedback is forcing the inverting input to have potential $u_{\mathrm{i}}$ (to follow the non-inverting input). Then voltage division of $u_{\mathrm{o}}$ between the two impedances provides the necessary relation between $u_{\mathrm{o}}$ and $u_{\mathrm{i}}: u_{\mathrm{i}}=u_{\mathrm{o}} \frac{R_{i}}{R_{i}+\frac{R . j \omega L}{R+j \omega L}}$.
The desired network function requires this to be inverted to give the ratio $u_{\mathrm{o}} / u_{\mathrm{i}}: \frac{u_{\mathrm{o}}}{u_{\mathrm{i}}}=1+\frac{\mathrm{j} \omega L / R_{i}}{1+\mathrm{j} \omega L / R}$.
This could be written in different forms, e.g. by changing terms like $R / L$ for a frequency $\omega_{x}$, or by keeping terms $(R+\mathrm{j} \omega L)$ instead of making them dimensionless; the question did not state that the canonical form is required.
b)

Ignoring the ' $1+$ ' part of $H(\omega)$, there is a $+20 \mathrm{~dB} /$ decade slope for $\omega<\omega_{\mathrm{p}}$, and a $0 \mathrm{~dB} /$ decade slope for $\omega>\omega_{\mathrm{p}}$. In the region $\omega=\omega_{z}$, the amplitude passes through 0 dB .

This shape can be understood by combining the separate effects of the numerator, $\mathrm{j} \omega / \omega_{\mathrm{z}}$, which has $+20 \mathrm{~dB} /$ decade slope everywhere, and the denominator, $1+\mathrm{j} \omega / \omega_{\mathrm{p}}$, which has $-20 \mathrm{~dB} /$ decade slope at $\omega>\omega_{\mathrm{p}}$ and flat 0 dB at $\omega<\omega_{\mathrm{p}}$ ).

To include the ' $1+$ ' term, notice that this ensures that $|H(\omega)|$ does not fall below $1(0 \mathrm{~dB})$. (This claim depends on the main term not having a negative real part; in our case this is true, as the main term's phase varies only from $90^{\circ}$ to $0^{\circ}$.

In the given situation of $\omega_{\mathrm{z}} \ll \omega_{\mathrm{p}}$, all the region $\omega>\omega_{\mathrm{z}}$ would have $|H(\omega)|>1$, so the ' $1+$ ' is negligible. The Bode plot can be modified to include the ' $1+$ ' by forcing the plot up to a constant 0 dB for $\omega<\omega_{\mathrm{z}}$.

One could also mention that the gain at the high frequencies, $\omega>\omega_{\mathrm{p}}$, is approximately $20 \log _{10} \frac{\omega_{p}}{\omega_{z}}$. Pedantically, we could say the points on the horizontal axis should be marked as e.g. $\log \omega_{\mathrm{p}}$ instead of $\omega_{\mathrm{p}}$.

Q5
a) Component $C$ is irrelevant to the rest of the circuit, as it is in series with a current source. As we are only interested in what happens outside the branch of $I$ and $C$, we can replace these two by just $I$.

The circuit is then almost exactly a Norton equivalent already.
The Norton current source is $I_{\mathrm{N}}=I$, still in the same direction with regard to terminals 'a' and 'b'.
The conductance $\frac{1}{R}$, and susceptance $\frac{1}{\omega L}$, are combined into a single admittance, $Y_{\mathrm{N}}=\frac{1}{R}+\frac{1}{\mathrm{j} \omega L}$.
As an alternative (less conventional for Norton sources) we can write this as an impedance, $Z_{\mathrm{N}}=\frac{\mathrm{j} \omega L R}{R+\mathrm{j} \omega L}$.
The same result could be found by the more general way, of calculating short-circuit current and open-circuit voltage.
b) The ac maximum power-transfer condition is that the load's impedance (or admittance) equals the complex conjugate of the impedance (or admittance, respectively) of the source.
So we want to make a load admittance $Y_{x}$ such that $Y_{x}=Y_{\mathrm{N}}^{*}=\frac{1}{R}-\frac{1}{\mathrm{j} \omega L}$, or equivalently, $Y_{x}=\frac{1}{R}+\mathrm{j} \frac{1}{\omega L}$.
Notice that this admittance has a positive real part and positive imaginary part. A resistor has purely positive real impedance or admittance. An inductor or capacitor has a purely imaginary impedance or admittance, which can be negative or positive depending on the component and on whether we use impedance or admittance.
When connecting admittances in parallel, they sum directly to the total admittance (impedances sum when series). The positive real and imaginary parts of the required admittance can therefore be independently implemented with a resistor $R_{x}$ and capacitor $C_{x}$. These give an admittance of $\frac{1}{R_{x}}+\mathrm{j} \omega C_{x}$. We have to choose $R_{x}$ and $C_{x}$ to give the required value of $Y_{x}$.
Considering the real and imaginary parts separately, we see $R_{x}=R$, and $\mathrm{j} \omega C_{x}=\mathrm{j} \frac{1}{\omega L}$, i.e. $C_{x}=\frac{1}{\omega^{2} L}$.
The other way to implement the right load for maximum power would be a series connection of a resistor and capacitor. In this case, the necessary values would be different from the ones needed for a parallel connection: let's call them $R_{y}$ and $C_{y}$. The criterion would be $R_{y}-\mathrm{j} \frac{1}{\omega C_{y}}=Z_{y}=Z_{\mathrm{N}}^{*}=\frac{1}{Y_{\mathrm{N}}^{*}}$. It is useful to get this last term into a form where the real and imaginary parts are separate so that they can be equated with $R_{y}$ and $\frac{1}{\omega C_{y}}$.
Hence $\frac{1}{Y_{\mathrm{N}}^{*}}=\frac{1}{\frac{1}{R}-\frac{1}{\mathrm{j} \omega L}}=\frac{\mathrm{j} \omega L R}{\mathrm{j} \omega L-R}=\frac{\omega^{2} L^{2} R-\mathrm{j} \omega L R^{2}}{R^{2}+\omega^{2} L^{2}}=R_{y}-\mathrm{j} \frac{1}{\omega C_{y}}$.
From this, $R_{y}=\frac{\omega^{2} L^{2} R}{R^{2}+\omega^{2} L^{2}}$, and $-\mathrm{j} \frac{1}{\omega C_{y}}=-\mathrm{j} \frac{\omega L R^{2}}{R^{2}+\omega^{2} L^{2}}$, i.e. $C_{y}=\frac{R^{2}+\omega^{2} L^{2}}{\omega^{2} L R^{2}}$.
Notice that the choice of parallel components was simpler in this case, as it better matches the known components $R$ and $L$ in the source circuit.

Q6
a) The independent sources have different frequencies. (It might have been clearer just to have written them as $\omega_{1}$ and $\omega_{2}$, instead of $\omega$ and $2 \omega$.) Superposition is therefore needed if we are going to take advantage of the convenient method of steady-state sinusoidal analysis, which solves for a specific frequency.

Consider the contribution of each source separately, to the marked current $i_{x}(t)$.
Case 1. With just the current source active, the voltage source is a short-circuit. The circuit can be redrawn with the irrelevant inductor omitted and the voltage source being a short-circuit; then we see that the current
at frequency $2 \omega$ divides through $R$ and $C$. Let's define the current source in the frequency-domain as being a current phasor $I=\hat{I} \angle \underline{0}$, which means we've used a sine convention and peak-value convention.
Then by current-division, $i_{x(1)}(2 \omega)=I \frac{\frac{1}{\mathrm{j} 2 \omega C}}{R+\frac{1}{\mathrm{j} 2 \omega C}}=\frac{1}{1+\mathrm{j} 2 \omega C R} I=\frac{\hat{I}}{\sqrt{1+4 \omega^{2} C^{2} R^{2}}} /-\tan ^{-1} \frac{2 \omega C R}{1}$.
Now returning to the time-domain, remembering to use the same sine- and peak conventions as we used at the start, $i_{x(1)}(t)=\frac{\hat{I}}{\sqrt{1+4 \omega^{2} C^{2} R^{2}}} \sin \left(2 \omega t-\tan ^{-1}(2 \omega C R)\right)$.
Case 2. Now with just the voltage source active, the current source can be ignored (open circuit), and the circuit is simply a loop of $U, R, C$. Let us define the voltage source as $U=\hat{U} \not \theta$, which means we've used a cosine reference and peak values. Notice that the reference can be independently chosen for each superposition-case; what is important is to choose the same reference for the time-frequency conversion as for the frequency-time conversion, for any particular case, in order to get the correct time-functions.
The marked current in this second case is $i_{x(2)}(\omega)=\frac{U}{R-\frac{j}{\omega C}}=\frac{\hat{U}}{\sqrt{R^{2}+\frac{1}{\omega^{2} C^{2}}}}\left\langle\theta+\tan ^{-1} \frac{1}{\omega C R}\right.$.
In the time-domain, using the cosine reference that we used when converting to the frequency-domain, this current is $i_{x(2)}(t)=\frac{\hat{U}}{\sqrt{R^{2}+\frac{1}{\omega^{2} C^{2}}}} \cos \left(\omega t+\theta+\tan ^{-1} \frac{1}{\omega C R}\right)$.
Finally: By superposition, the actual current $i_{x}(t)$ is the sum of both the above results,
$i_{x}(t)=\frac{\hat{I}}{\sqrt{1+4 \omega^{2} C^{2} R^{2}}} \sin \left(2 \omega t-\tan ^{-1} 2 \omega C R\right)+\frac{\hat{U}}{\sqrt{R^{2}+\frac{1}{\omega^{2} C^{2}}}} \cos \left(\omega t+\theta+\tan ^{-1} \frac{1}{\omega C R}\right)$.
b) The mean power (over multiple periods of the source frequencies) could be found by integrating instantaneous power over some time $T$ during which a whole number of cycles of both frequencies have occurred, i.e. $\bar{P}=\frac{1}{T} \int_{t}^{t+T} R i^{2}(t) \mathrm{d} t$.
But there is a more convenient method, power superposition, which is valid only for this situation where the sources being handled by superposition have different frequencies.
Looking back to the frequency-domain expressions for $i_{x}(\cdot)(\omega)$ in the working of part 'a)', the active power in each superposition case is easily found by squaring the magnitude of the current, multiplying by the resistance, and dividing by 2 (because of peak values being used).
Summing the results for the two frequencies, this is $\frac{R}{2}\left(\frac{\hat{I}^{2}}{1+4 \omega^{2} C^{2} R^{2}}+\frac{\hat{U}^{2}}{R^{2}+\frac{1}{\omega^{2} C^{2}}}\right)$.

