## SPRING 2023: MATH 791 DAILY UPDATE

Wednesday, May 3. The assigned teams worked on Exam 3.
Monday, May 1. The assigned teams worked on Exam 3.
Friday, April 28. The assigned teams worked on Exam 3.
Wednesday, April 26. We spent most of the class using the proof of the Inverse Galois Problem for cyclic groups to construct a field $K$ such that $K$ is a Galois extension of $\mathbb{Q}$ and $\operatorname{Gal}(K / \mathbb{Q}) \cong \mathbb{Z}_{8}$. For this, the first prime $p$ such that $p \equiv 1 \bmod 8$ is 17 . Thus, if $\epsilon$ is a primitive 17 th root of unity, $\operatorname{Gal}(\mathbb{Q}(\epsilon) / \mathbb{Q}) \cong \mathbb{Z}_{17}^{*}$, which is a cyclic group of order 16. To find a Galois extension $K$ of $\mathbb{Q}$ whose Galois group is isomorphic to $\mathbb{Z}_{8}$, we must take the fixed field of $H$, where $H \subseteq \operatorname{Gal}(\mathbb{Q}(\epsilon) / \mathbb{Q})$ is a subgroup of index 8 . Now, $\overline{5}$ is a cyclic generator for $\mathbb{Z}_{17}^{*}$ (check this), and thus, the proof of Theorem 38.1 tells us that $\sigma \in \operatorname{Gal}(\mathbb{Q}(\epsilon) / \mathbb{Q})$ given by $\sigma(\epsilon)=\epsilon^{5}$ is a cyclic generator of $\operatorname{Gal}(\mathbb{Q}(\epsilon) / \mathbb{Q})$. Therefore $\sigma^{8}$ is an element of order two, and so it generates a subgroup of index 8 . If we call this subgroup $H$, then $\gamma \in \mathbb{Q}(\epsilon)$ is in the fixed field of $H$ if and only if $\sigma^{8}(\gamma)=\gamma$
Now, the minimal polynomial for $\epsilon$ over $\mathbb{Q}$ is $\Phi_{17}(x)=x^{16}+x^{15}+\cdots+x+1$, so that $1, \epsilon, \ldots, \epsilon^{15}$ is a basis for $\mathbb{Q}(\epsilon)$ over $\mathbb{Q}$. Thus, if $\gamma \in \mathbb{Q}(\epsilon)$, we can write

$$
\gamma=a_{0}+a_{1} \epsilon+a_{2} \epsilon^{2}+\cdots+a_{15} \epsilon^{15}
$$

and therefore,

$$
\sigma^{8}(\gamma)=a_{0}+a_{1} \sigma^{8}(\epsilon)+a_{2} \sigma^{8}\left(\epsilon^{2}\right)+\cdots+a_{15} \sigma^{8}\left(\epsilon^{15}\right)
$$

A somewhat tedious (though not difficult) calculation shows that

$$
\begin{aligned}
\sigma^{8}(\epsilon) & =-1-\epsilon-\epsilon^{2}-\cdots-\epsilon^{15} \\
\sigma^{8}\left(\epsilon^{2}\right) & =\epsilon^{15}, \sigma^{8}\left(\epsilon^{3}\right)=e^{14}, \sigma^{8}\left(\epsilon^{4}\right)=\epsilon^{13} \\
\sigma^{8}\left(\epsilon^{5}\right) & =\epsilon^{12}, \sigma^{8}\left(\epsilon^{6}\right)=\epsilon^{11}, \sigma^{8}\left(e^{7}\right)=\epsilon^{10} \\
\sigma^{8}\left(\epsilon^{8}\right) & =\epsilon^{9}, \sigma^{8}\left(\epsilon^{9}\right)=\epsilon^{8}, \sigma^{8}\left(\epsilon^{10}\right)=\epsilon^{7} \\
\sigma^{8}\left(\epsilon^{11}\right) & =\epsilon^{6}, \sigma^{8}\left(\epsilon^{12}\right)=\epsilon^{5}, \sigma^{8}\left(\epsilon^{13}\right)=\epsilon^{4} \\
\sigma^{8}\left(\epsilon^{13}\right) & =\epsilon^{4}, \sigma^{8}\left(\epsilon^{14}\right)=\epsilon^{3}, \sigma^{8}\left(\epsilon^{15}\right)=\epsilon^{2} .
\end{aligned}
$$

It follows that

$$
\sigma^{8}(\gamma)=\left(a_{0}-a_{1}\right)+\left(-a_{1}\right) \epsilon+\left(a_{15}-a_{1}\right) \epsilon^{2}+\left(a_{14}-a_{1}\right) \epsilon^{3}+\cdots+\left(a_{2}-a_{1}\right) \epsilon^{15}
$$

Setting $\gamma=\sigma^{8}(\gamma)$, we obtain

$$
a_{0}=a_{0}, a_{1}=0, a_{2}=a_{15}, a_{3}=a_{14}, a_{3}=a_{13}, \ldots, a_{14}=a_{3}, a_{15}=a_{2}
$$

Therefore,

$$
\gamma=a_{0}+a_{2}\left(\epsilon^{2}+\epsilon^{15}\right)+a_{3}\left(\epsilon^{3}+\epsilon^{14}\right)+\cdots+a_{6}\left(\epsilon^{6}+\epsilon^{11}\right)+a_{7}\left(\epsilon^{7}+\epsilon^{10}\right)+a_{8}\left(\epsilon^{8}+\epsilon^{9}\right)
$$

It follows that $K:=\mathbb{Q}\left(e^{2}+\epsilon^{15}, \epsilon^{3}+\epsilon^{14}, \ldots, \epsilon^{6}+\epsilon^{11}, \epsilon^{7}+\epsilon^{10}, \epsilon^{8}+\epsilon^{9}\right)$ is the fixed field of $H$, and $\operatorname{gal}(K / \mathbb{Q}) \cong \mathbb{Z}_{8}$, as required. We can simplify $K$ by noting that

$$
\begin{aligned}
& \left(\epsilon^{2}+\epsilon^{15}\right)^{2}=\epsilon^{4}+2+\epsilon^{13} \\
& \left(\epsilon^{3}+\epsilon^{14}\right)^{2}=\epsilon^{6}+2+\epsilon^{11} \\
& \left(\epsilon^{5}+\epsilon^{12}\right)^{2}=\epsilon^{10}+2+\epsilon^{7} \\
& \left(\epsilon^{4}+\epsilon^{13}\right)^{2}=\epsilon^{8}+2+\epsilon^{9}
\end{aligned}
$$

It follows that $K=\mathbb{Q}\left(\epsilon^{2}+\epsilon^{15}, \epsilon^{3}+\epsilon^{14}, \epsilon^{5}+\epsilon^{12}\right)$. Can you find a simpler expression for $K$ ?

We ended class by sketching a proof of the following theorem:
Theorem. Let $G$ be a finite group. Then there exists a finite, Galois extension of fields $F \subseteq K$ such that $\operatorname{Gal}(K / F) \cong G$.
If $|G|=n$, the idea of the proof was to let $K$ be the rational function field in $n$ variables over $\mathbb{Q}$ and to let $S_{n}$ be the group of automorphisms of $K$ obtained by permuting the given variables. Letting $F_{0}$ denote the field of symmetric rational functions, then $\operatorname{Gal}\left(K / F_{0}\right)=S_{n}$. Identifying $G$ as a subgroup of $S_{n}$ and taking $F$ to be the fixed field $K^{G}$, it followed from the Galois Correspondence Theorem that $F \subseteq K$ is a Galois extension with Galois group $G$.
Monday, April 24. After reviewing the the statement of the Galois Correspondence Theorem, we illustrated the theorem by calculating the intermediate fields and corresponding subgroups of the Galois groups for the two extensions: $\mathbb{Q} \subseteq \sqrt{2}, \sqrt{3})$ and $\mathbb{Q} \subseteq \mathbb{Q}(\sqrt[3]{2}, \epsilon)$, where $\epsilon=e^{\frac{2 \pi i}{3}}$. This was made easy by the fact that we had previously calculated both Galois groups and the fixed fields of the first extension. The second example had the advantage of having intermediate fields not Galois over $\mathbb{Q}$ that correspond to subgroups of the Galois group that are not normal subgroups.

We then mentioned the well-known (and not fully determined!):
Inverse Galois Problem. Is every finite group the Galois group of a Galois extension of $\mathbb{Q}$ ?
We ended class by showing that the Inverse Galois Problem has a positive solution for $G=\mathbb{Z}_{n}$, assuming the following special case of Direchlet's Theorem: For $n$ fixed, there exist (infinitely many) prime numbers in the arithmetic progression $\{m n+1\}_{m \geq 1}$. Another key ingredient of the proof was that for $p$ prime, the multiplicative group $\left(\mathbb{Z}_{p}\right)^{*}$ is cyclic, a fact we proved when we showed that a finite extension of finite fields always has a primitive element.
Friday, April 24. We continued our discussion of Galois groups and Galois extensions by first offering the following theorem.

Theorem. Suppose that $F \subseteq K$ is a finite extension with a primitive element, so that $K=F(\alpha)$. Let $p(x)$ denote the minimal polynomial of $\alpha$ over $F$ and write $d=\operatorname{deg}(p(x))$. Then $K$ is Galois over $F$ if and only if $p(x)$ has $d$-distinct roots in $K$.
We noted that the theorem follows immediately from the Crucial Proposition from the lecture of April 14. We also noted, that it is not enough to just assume $K$ is the splitting field of $p(x)$ over $F$, since in positive characteristic, $p(x)$ need not have distinct roots. However, if $\mathbb{Q} \subseteq F$ or $F$ is finite, then the theorem shows that $K$ is Galois over $F$ if and only $F$ is the splitting field of $p(x)$ over $F$, since in these cases, the irreducible polynomial $p(x)$ is guaranteed to have distinct roots.

We then noted the following property of splitting fields, stated and proved here for the special case that a primitive element exists.
Theorem. Let $K=F(\alpha)$ be a finite extension of $F$ and assume that $K$ is the splitting field of the minimal polynomial of $\alpha$ over $F$. Then if $f(x) \in F[x]$ is a non-constant, irreducible polynomial with a root in $K$, then $f(x)$ splits over $K$.

We then defined the concept fixed field. Let $F \subseteq K$ be a finite extension with Galois group $\operatorname{Gal}(K / F)$. For $\sigma \in \operatorname{Gal}(K / F), K^{\sigma}:=\{\alpha \in K \mid \sigma(\alpha)=\alpha\}$ if the fixed field of $\sigma$. For $H$ a subgroup of $\operatorname{Gal}(K / F)$, $K^{H}:=\{\alpha \in K \mid \sigma(\alpha)=\alpha$, for all $\sigma \in H\}$ is the fixed field of $H$. We showed that the fixed fields are intermediate fields between $F$ and $K$ and also calculated the fixed fields of the subgroups of $\operatorname{Gal}(K / F)$ when $F=\mathbb{Q}$ and $K=\mathbb{Q}(\sqrt{2}, \sqrt{3})$. We ended class by stating the:
Galois Correspondence Theorem. Let $F \subseteq K$ be a finite Galois extension and set $G:=\operatorname{Gal}(K / F)$. Then:
(i) There is a 1-1, onto, order reversing correspondence between the subgroups $H \subseteq G$ and the intermediate fields $F \subseteq E \subseteq K$, given by $H \rightarrow K^{H}$ and $E \rightarrow \operatorname{Gal}(K / E)$. In particular, for all $H$ and $E$, $H=\operatorname{Gal}\left(K / K^{H}\right)$ and $E=K^{\operatorname{Gal}(K / E)}$.
(ii) If $H$ and $E$ correspond, then $[G: H]=[E: F]$.
(iii) For any intermediate field $E, K$ is Galois over $E$.
(iv) An intermediate field $E$ is Galois over $F$ if and only if $\operatorname{Gal}(K / E)$ is a normal subgroup of $G$. In which case, $\operatorname{Gal}(E / F) \cong G / \operatorname{Gal}(K / E)$.
In particular, there are only finitely many intermediate fields between $F$ and $K$.
Wednesday, April 19. We began class by proving the crucial proposition stated in the lecture of April 14. A key point was that if $\sigma: F_{1} \rightarrow F_{2}$ is an isomorphism of fields, then $\hat{\sigma}: F_{1}[x] \rightarrow F_{2}[x]$ given by $\hat{\sigma}(f(x)):=f(x)^{\sigma}$ is an isomorphism of rings. We then discussed and gave proofs of the following applications of the crucial proposition. The proofs below clarify the discussion in class.

First Application. Let $F_{1} \subseteq K_{1}, F_{2} \subseteq K_{2}$ be fields, $\left.p_{1}(x) \in F_{1}\right], p_{2}(x) \in F_{2}[x]$ non-constant polynomials, and $K_{1}$ the splitting of $p_{1}(x)$ over $F_{1}$ and $K_{2}$ the splitting field of $p_{2}(x)$ over $F_{2}$. Suppose $\sigma: F_{1} \rightarrow F_{2}$ is an isomorphism such that $p_{2}(x)=p_{1}(x)^{\sigma}$. Then there exists an isomorphism $\tilde{\sigma}: K_{1} \rightarrow K_{2}$ extending $\sigma$. In particular, if $K_{1}$ and $K_{2}$ are splitting fields of $f(x) \in F[x]$ over $F$, then there is an isomorphism $\tau: K_{1} \rightarrow K_{2}$ fixing $F$.

Proof. We induct on $n=\operatorname{deg}\left(p_{1}(x)\right)$. If $n=1$, there is nothing to prove. Suppose $n>1$. Let $\alpha_{1}, \ldots, \alpha_{n}$ be the (not necessarily distinct roots of $p_{1}(x)$ and $\beta_{1}, \ldots, \beta_{n}$ be the (not necessarily distinct roots of $p_{2}(x)$. Now, by the crucial proposition, there exists an isomorphism $\bar{\sigma}: F_{1}\left(\alpha_{1}\right) \rightarrow F_{2}\left(\beta_{1}\right)$ extending $\sigma$ (where $\bar{\sigma}$ takes $\alpha_{1}$ to a root of $q(x)^{\sigma}$, say $\beta_{1}$, where $q(x)$ is the minimal polynomial of $\alpha_{1}$ over $\left.F_{1}\right)$. If we write $p_{1}(x)=\left(x-\alpha_{1}\right) f_{1}(x) \in F_{1}\left(\alpha_{1}\right)[x]$ and $p_{2}(x)=\left(x-\beta_{1}\right) f_{2}(x) \in F_{2}\left(\beta_{1}\right)[x]$. Then, $K_{1}$ is the splitting field of $f_{1}(x)$ over $F_{1}\left(\alpha_{1}\right)$ and $K_{2}$ is the splitting field of $f_{2}(x)$ over $F_{2}\left(\beta_{1}\right)$. Now applying the induction hypothesis to $\bar{\sigma}, f_{1}(x)$ and $f_{2}(x)$ finishes the proof.

Second Application. Let $F \subseteq K$ be a finite extension of fields. Then the number of isomorphisms from $K$ to $\bar{F}$ fixing $F$ is less than or equal to $[K: F]$. In particular, $|\operatorname{Gal}(K / F)| \leq[K: F]$. In particular, for any finite extension $F \subseteq K,|\operatorname{Gal}(K / F)| \leq[K: F]$.

Proof. We can write $K=F\left(\alpha_{1}, \ldots, \alpha_{n}\right)$ and proceed by induction on $n$. When $n=1, K=F(\alpha)$, say. Any isomorphism $\sigma: K \rightarrow \bar{F}$ fixing $F$ must take $\alpha$ to another root of $p(x)$, the minimal polynomial of $p(x)$ over $F$. Since $p(x)$ has at most $\operatorname{deg}(p(x))=[K: F]$ roots in $\bar{F}$, there are at most $[K: F]$ such isomorphisms.

Suppose $n>1$ and set $E:=F\left(\alpha_{1}, \ldots, \alpha_{n-1}\right)$. By induction, there are at most $[E: F]$ isomorphisms $\tau: E \rightarrow \bar{F}$ fixing $F$. Now since $K=E\left(\alpha_{n}\right)$, any $\tau: K \rightarrow \bar{F}$ extending $\tau$ must take $\alpha_{n}$ to a root of $q(x)^{\tau}$, where $q(x)$ is the minimal polynomial of $\alpha_{n}$ over $E$. Since there are at most $[K: E]$ such roots, there are at most $[K: E]$ extensions of $\tau$. Thus, there are at most $[E: F] \cdot[K: E]$ isomorphisms from $K$ to $\bar{F}$ that are obtained by extending isomorphisms from $E$ to $\bar{F}$ to an isomorphism from $K$ to $\bar{F}$. However, given any isomorphism $\sigma: K \rightarrow \bar{F}$ fixing $F, \sigma$ is the extension of the isomorphism $\sigma_{\left.\right|_{E}}: E \rightarrow \bar{F}$, and hence we have accounted for all such isomorphisms. Therefore, the number of the number of isomorphisms from $K$ to $\bar{F}$ fixing $F$ is less than or equal to $[K: E] \cdot[E: F]=[K: F]$, as required.

The second statement is an immediate consequence. of the first.
We also noted (but did not prove in detail) that the crucial proposition can be used in conjunction with Zorn's Lemma to show that any two algebraic closures of $F$ are isomorphic via an isomorphism fixing $F$.

Monday, April 17. We began class by restating the proposition presented in the previous lecture and noting two consequences thereof:
(i) If $p(x) \in F[x]$ is irreducible over $F$ and $\alpha_{1}, \alpha_{2} \in \bar{F}$ are two roots of $p(x)$, then there is an isomorphism from $F\left(\alpha_{1}\right)$ to $F\left(\alpha_{2}\right)$ fixing $F$ and taking $\alpha_{1}$ to $\alpha_{2}$.
(ii) If $K=F(\alpha)$, for $\alpha$ algebraic over $F$ (e.g., $\mathbb{Q} \subseteq F$ or $|F|<\infty$ ), then $|\operatorname{Gal}(K / F)|$ equals the number of distinct roots of $p(x)$ in $K$, where $p(x)$ is the minimal polynomial of $\alpha$ over $F$.

We also noted, but did not prove, that $|\operatorname{Gal}(K / F)| \leq[K: F]$, for any finite extension $F \subseteq K$, and defined in general, the extension to be a Galois extension if $[K: F]=|\operatorname{Gal}(K / F)|$. We then spent the remainder of the class using the proposition referred to above to construct six automorphisms of the field $K:=\mathbb{Q}(\sqrt[3]{2}, \epsilon)$, where $\epsilon$ is a primitive 3rd root of unity, so that $K$ is the splitting field of $x^{3}-2$ over $\mathbb{Q}$. We further observed that in this case $\operatorname{Gal}(K / F) \cong S_{3}$ and that $K$ is a Galois extension of $\mathbb{Q}$.

Friday, Apri 14. After reviewing the definition of $\operatorname{Gal}(K / F)$ for the field extension $F \subseteq K$, we stated and discussed (but did not prove) the following proposition as a means to construct elements of $\mathrm{Gal}(K / F)$. .
Crucial Proposition. Let $F_{1} \subseteq K_{1}, F_{2} \subseteq K_{2}$ be fields, $\left.p_{1}(x) \in F_{1}\right], p_{2}(x) \in F_{2}[x]$ be monic irreducible polynomials of degree $d$, and $\alpha_{1} \in K_{1}, \alpha_{2} \in K_{2}$ roots of $p_{1}(x)$ and $p_{2}(x)$, respectively. Suppose $\sigma: F_{1} \rightarrow F_{2}$ is an isomorphism such that $p_{2}(x)=p_{1}(x)^{\sigma}$. Then there exists an isomorphism $\tilde{\sigma}: F_{1}\left(\alpha_{1}\right) \rightarrow F_{2}\left(\alpha_{2}\right)$ extending $\sigma$ such that $\tilde{\sigma}\left(\alpha_{1}\right)=\alpha_{2}$.

We noted that in the proposition, $p_{1}(x)^{\sigma}$ denotes the polynomial in $F_{2}[x]$ obtained by applying $\sigma$ to the coefficients of $p_{1}(x)$ and that the map $\phi: F_{1}[x] \rightarrow F_{2}[x]$ given by $\phi(f(x))=f(x)^{\sigma}$ is an isomorphism of rings. We also saw that if the degree of $p_{1}(x)$ equals $d$, then the isomorphism $\tilde{\sigma}$ is given by

$$
\tilde{\sigma}\left(\lambda_{0}+\lambda_{1} \alpha_{1}+\cdots+\lambda_{d-1} \alpha_{1}^{d-1}\right)=\sigma\left(\lambda_{0}\right)+\sigma\left(\lambda_{1}\right) \alpha_{2}+\cdots+\sigma\left(\lambda_{d-1}\right) \alpha_{2}^{d-1}
$$

We then applied the proposition above to the extension $\mathbb{Q} \subseteq \mathbb{Q}(\sqrt{2}, \sqrt{3})=\mathbb{Q}(\sqrt{2})+(\sqrt{3})$. Taking $F_{1}=\mathbb{Q}(\sqrt{2})=F_{2}$ and $\sigma: F_{1} \rightarrow F_{2}$ defined by $\sigma(a+b \sqrt{2})=a-b \sqrt{2}$, we noted that the proposition enabled us to extend $\sigma$ to $K=F_{1}(\sqrt{3})$ in two ways, $\tilde{\sigma_{1}}: F_{1}(\sqrt{3}) \rightarrow F_{2}(\sqrt{3})$ such that $\tilde{\sigma}_{1}(\sqrt{3})=\sqrt{3}$ and $\tilde{\sigma_{2}}: F_{1}(\sqrt{3}) \rightarrow F_{2}(\sqrt{3})$ such that $\tilde{\sigma_{2}}(\sqrt{3})=-\sqrt{3}$. Similarly, if we take $\tau: F_{1} \rightarrow F_{1}$ to be the identity, then the proposition enabled us to extend $\tau$ to $K=F_{1}(\sqrt{3})$ in two ways, first to $\tilde{\tau_{1}}$ which takes $\sqrt{3}$ to $\sqrt{3}$ and second, to $\tilde{\tau_{2}}$ which takes $\sqrt{3}$ to $-\sqrt{3}$. The maps described above are all automorphisms of $K$. We then noted that these automorphisms form a group isomorphic to $\mathbb{Z}_{2} \times \mathbb{Z}_{2}$. This also enabled us to see that for $p(x)=x^{4}-10 x^{2}+1$, the minimal polynomial of $\sqrt{2}+\sqrt{3}$ over $\mathbb{Q}, p(x)$ splits over $K$, and its roots are: $\sqrt{2}+\sqrt{3}, \sqrt{2}-\sqrt{3},-\sqrt{2}+\sqrt{3},-\sqrt{2}-\sqrt{3}$.
Wednesday, April 12. We spent most of the class proving the following important theorem:
Theorem. Let $F$ be a field. There is an algebraic extension $F \subseteq \bar{F}$ such that if $p(x) \in F[x]$ and the degree of $p(x)$ is $d>0$, then there exist $\alpha_{1}, \ldots, \alpha_{d} \in F$ (not necessarily distinct) such that $p(x)=\left(x-\alpha_{1}\right) \cdots\left(x-\alpha_{d}\right)$.
The field $\bar{F}$ above is called an algebraic closure of $F$. The proof of the theorem followed by iterating the construction in the final proposition from the previous lecture and taking the union of the resulting fields. We noted that $\mathbb{C}$ is the algebraic closure of $\mathbb{R}$, and sketched a proof using Louiville's theorem from complex analysis.

We ended class by defining the Galois group of a field extension $F \subseteq K$ : Let $\operatorname{Gal}(K / F)$ denote the set of automorphisms of $K$ fixing $F$. We call $\operatorname{Gal}(K / F)$ the Galois group of $K$ over $F$. We noted the following properties, as listed from Homework 27:
(i) $\operatorname{Gal}(K / F)$ is a group.
(ii) If $f(x) \in F[x], \alpha \in K$ satisfies $f(\alpha)=0$, then $f(\sigma(\alpha))=0$, for all $\sigma \in \operatorname{Gal}(K / F)$.
(iii) If $K=F(\alpha)$, for $\alpha \in K$ is a primitive element, then $\operatorname{Gal}(K / F)$ is finite. In particular, if $F \subseteq K$ is a finite extension, with $\mathbb{Q} \subseteq F$, then $\operatorname{Gal}(K / F)$ is a finite group.

Monday, April 10. We began class by establishing the following facts for a field extension $F \subseteq K$ :
(i) If $\alpha, \beta \in K$ are algebraic over $F$, then $\alpha \pm \beta, \alpha \beta, \alpha \beta^{-1}$ are algebraic over $F$.
(ii) The set $E$ of elements in $K$ that are algebraic over $F$ form a subfield of $K$ containing $F$, called the algebraic closure of $F$ in $K$.

We then noted that if $E$ is the algebraic closure of $F$ in $K$, then $E$ is algebraically closed in $K$, i.e., $E$ equals the algebraic closure of $E$ in $K$. This followed from the proposition showing that if $F \subseteq E \subseteq K$ is an extension of fields, with $E$ algebraic over $F$ and $K$ is algebraic over $E$, then $K$ is algebraic over $F$.

We then proved the following proposition and its corollary (which followed from the proposition and the Primitive Element Theorem), which are relevant to a study of Galois theory.
Proposition. Let $F \subseteq K$ be a finite extension, with $F$ an infinite field. Then there is a primitive element for the extension if and only if there are (only) finitely many intermediate fields $F \subseteq E \subseteq K$.

Corollary. Let $F \subseteq K$ be a finite extension of fields, with $\mathbb{Q} \subseteq F$. Then there are only finitely many intermediate fields $F \subseteq E \subseteq K$.

We ended class with the first key step in showing the existence of algebraic closures.
Proposition. Let $F$ be a field. Then there exists a field extension $F \subseteq \tilde{F}$ with the following property: For all $0 \neq f(x) \in F[x]$, there exists $\alpha \in \tilde{F}$ such tht $f(\alpha)=0$.
Friday, April 7. We began class by giving the following definition.
Definition. Let $F \subseteq K$ be an extension of fields.
(i) $\alpha \in K$ is algebraic over $F$ if there is a non-constant polynomial $p(x) \in F[x]$ such that $p(\alpha)=0$.
(ii) $K$ is an algebraic extension of $F$ is every element of $K$ is algebraic over $F$.

This was followed by establishing the following important proposition.
Proposition. For $F \subseteq K$ fields and $\alpha \in K, \alpha$ is algebraic over $K$ if and only if $[F(\alpha): F]<\infty$.
The if direction followed using the same determinantal trick as the one used in the example from the previous lecture, showing $\sqrt{2}+\sqrt{3}$ is a root of $x^{4}-10 x^{2}+1$.

We then stated and proved the following version of the primitive element theorem.
Theorem. Suppose $F \subseteq K$ is an extension of fields satisfying $[K: F]<\infty$. If $\mathbb{Q} \subseteq F$ or $F$ is finite, then there exists a primitive element $\alpha \in K$ such that $K=F(\alpha)$.
When $\mathbb{Q} \subseteq F$, the proof quickly reduced to the case that $K=F(u, v)$ and ultimately showed that, in this case, for all but finitely many $\lambda \in F, \alpha:=u+\lambda v$ is a primitive element. When $F$ is finite, we appealed to the standard fact that a finite abelian group is a direct product of cyclic groups. This lead to the fact that the multiplicative group $(K \backslash\{0\}, \cdot)$ is cyclic, and if $\alpha \in K$ is a cyclic generator, then $K=F(\alpha)$.

We ended class by showing that if $F:=\mathbb{Z}_{2}\left(x^{2}, y^{2}\right)$ is the rational function field in $x^{2}, y^{2}$ over $\mathbb{Z}_{2}$ and $K=\mathbb{Z}_{2}(x, y)$, then $K$ is a finite extension of $F$ (of degree four) and there is no primitive element for this extension.
Wednesday, April 5 . We began class by showing that $\mathbb{Q}(\sqrt{2}, \sqrt{3})=\mathbb{Q}(\sqrt{2}+\sqrt{3})$. The key point here was that $p(x)=x^{4}-10 x^{2}+1$ is irreducible and has $\sqrt{2}+\sqrt{3}$ as a root. We arrived at $p(x)$ by starting with the basis $1, \sqrt{2}, \sqrt{3}, \sqrt{6}$, multiplying each basis element by $\sqrt{2}+\sqrt{3}$, and writing the result in terms of the basis with coefficients in $\mathbb{Q}$. This gave rise to a $4 \times 4$ homogeneous system of equations with non-trivial solution. The determinant of the coefficient matrix for this system was $p(\sqrt{2}+\sqrt{3})$, which is therefore zero.

We then defined the concept of primitive element: For a finite extension of fields $F \subseteq K, \alpha \in K$ is a primitive element, if $K=F(\alpha)$. The purpose of the example above was to motivate the following version of the Primitive Element Theorem.
Theorem. Suppose $F \subseteq K$ is an extension of fields satisfying $[K: F]<\infty$. If $\mathbb{Q} \subseteq F$, then there exists a primitive element $\alpha \in K$ such that $K=F(\alpha)$.

As a tool for proving the Primitive Element Theorem, we proved the following proposition:
Proposition. If $F$ is a field containing $\mathbb{Q}$ and $p(x) \in F[x]$ is irreducible, then $p(x)$ has distinct roots in $K$, the splitting of $p(x)$ over $F$.
This was followed by noting that the proposition above fails, if $F$ does not contain $\mathbb{Q}$. Taking $F:=\mathbb{Z}_{2}\left(t^{2}\right)$ and $K:=\mathbb{Z}_{t}(t)$, the rational function fields over $\mathbb{Z}_{2}$ in the variables $t^{2}$ and $t$, then $p(x)=x^{2}-t^{2}$ is irreducible over $F$, but has the repeated root $t$ in its splitting field $K$, since $p(x)=(x-t)^{2}$ over $K$. Thus, $p(x)$ is an irreducible polynomial with a repeated root in its splitting field.
Monday, April 3. We began by observing that the splitting fields for $x^{2}-2$ and $x^{2}+1$ over $\mathbb{Q}$ are $\mathbb{Q}(\sqrt{2})$ and $\mathbb{Q}(i)$, respectively, and that $[\mathbb{Q}(\sqrt{2}: \mathbb{Q}]=2=[\mathbb{Q}(i): \mathbb{Q}]$. We then noted that the splitting field for $x^{3}-2$ over $\mathbb{Q}$ is $\mathbb{Q}\left(\sqrt[3]{2}, \sqrt[3]{2} \epsilon, \sqrt[3]{2} \epsilon^{2}\right)=\mathbb{Q}(\sqrt[3]{2}, \epsilon)$, where $\epsilon:=e^{\frac{2 \pi i}{3}}$ is a primitive cube root of 1 . Since $\epsilon$ and $\epsilon^{2}$ are roots of $x^{2}+x+1$, we saw that $\epsilon, \epsilon^{2}($ in order $)=\frac{-1 \pm \sqrt{3} i}{2}$. To calculate $[\mathbb{Q}(\sqrt[3]{2}, \epsilon): \mathbb{Q}]$, we needed the following proposition, which yields: $[\mathbb{Q}(\sqrt[3]{2}, \epsilon): \mathbb{Q}]=[\mathbb{Q}(\sqrt[3]{2}, \epsilon): \mathbb{Q}(\sqrt[3]{2})] \cdot[\mathbb{Q}(\sqrt[3]{2}): \mathbb{Q}]=2 \cdot 3=6$.
Proposition. Let $F \subseteq K \subseteq L$ be fields. Then $[L: F]$ is finite if and only if $[L: K]$ and $[K: F]$ are finite, in which case, $[L: F]=[L: K] \cdot[K: F]$.

The proof of the proposition followed by showing that if $w_{1}, \ldots, w_{t} \in L$ form a basis for $L$ over $K$ and $v_{1}, \ldots, v_{r} \in K$ form a basis for $K$ over $F$, then $\left\{v_{i} w_{j}\right\}_{1 \leq i \leq r, 1 \leq j \leq t}$ forms a basis for $L$ over $F$. We ended class by noting that: (i) If $F \subseteq K$ is a field extension with $[K: F]=p$, prime, then there are no intermediate fields properly between $F$ and $K$ and (ii) If $\alpha$ is a root of the irreducible polynomial $x^{2}+x+2 \in \mathbb{Z}_{3}[x]$, then $\alpha, 2-\alpha$ are the roots of $x^{2}+x+2, \mathbb{Z}_{3}(\alpha)$ is its splitting field and $\mathbb{Z}(\alpha)$ is a field with nine elements.
Friday, March 31. We gave a hands on proof of the following theorem: Let $F$ be a field, $p(x) \in F[x]$ a non-constant polynomial. Then there exists a field $K$ containing $F$ and $\alpha \in K$ such that $p(\alpha)=0$. By hands on, we mean we did not just write $K:=F[x] /\langle p(x)\rangle$ in the case $p(x)$ is irreducible and argue that $K$ is a field and $\alpha:=\bar{x}$ is a root of $p(x)$, which is the standard proof. Though we did note this at the end of class by observing $p(\bar{x})=\overline{p(x)}=\overline{0}$ in $F[x] /\langle p(x)\rangle$. Instead, we took a new variable which we called $z$ and considered the set $K:=\left\{a_{0}+a_{1} z+\cdots+a_{d-1} z^{d-1} \mid a_{j} \in F\right\} \subseteq F[z]$, and defined addition and multiplication on $K$ as follows: For $A=a_{0}+a_{1} z+\cdots+a_{d-1} z^{d-1}$ and $B=b_{0}+b_{1} z+\cdots+b_{d-1} z^{d-1}$,

$$
A+B:=\left(a_{0}+b_{0}\right)+\left(a_{1}+b_{1}\right) \alpha+\cdots+\left(a_{d-1}+b_{d-1}\right) \alpha^{d-1} \quad \text { and } \quad A * B:=r(z)
$$

where $A(x) B(x)=p(x) q(x)+r(x)$, according to the division algorithm, and $A(x)=a_{0}+a_{1} x+\cdots+a_{d-1} x^{d-1}$ and $B(x)=b_{0}+b_{1} x+\cdots+b_{d-1} x^{d-1}$. We then showed that, with these operations, $K$ is a field containing $F$ and $p(z)=0$. The point of this construction is to observe that, when we already have $K$ and $\alpha$, a root of $p(x)$, as in the previous lecture, addition and multiplication in $F(\alpha)$ tells us what we should expect when we do not, a priori, have $\alpha$ and $K$. Of course this construction just makes apparent what is really going on in the standard abstract proof noted above. We finished class by observing that if $p(x)$ has degree $d$, then the theorem can be iterated to show that there exists a field $\tilde{K}$ containing $F$ and $\alpha_{1}, \ldots, \alpha_{d} \in \tilde{K}$, not necessarily distinct, such that $p(x)=\left(x-\alpha_{1}\right) \cdots\left(x-\alpha_{d}\right)$. We say that $p(x)$ splits over $\tilde{K}$. The field $F\left(\alpha_{1}, \ldots, \alpha_{d}\right)$ is called a splitting field for $p(x)$ over $F$.
Wednesday, March 29. We continued our discussion of fields by noting that we want to generalize the constructions from the previous lecture, where we began with the field $\mathbb{Q}$, particular irreducible polynomials over $\mathbb{Q}$ with a designated root, and used this information to construct new fields. We therefore, want to start with fields $F \subseteq K, \alpha \in K$ a root of an irreducible polynomial over $F$ and construct a new field $F(\alpha)$. We first began by proving properties of the monic polynomial $p(x) \in F[x]$ of least degree with $p(\alpha)=0$. We noted that $p(x)$ has the properties: (i) $p(x)$ is irreducible over $F$; (ii) $p(x)$ divides any $g(x) \in F[x]$ having $\alpha$ as a root and (iii) $p(x)$ is unique. We noted $p(x)$ is called the minimal polynomial of $\alpha$ over $F$. We then showed that if $F \subseteq K$ are fields, and $\alpha \in K$ has minimal polynomial $p(x)$ with degree $d>0$, then $F(\alpha):=\left\{a_{0}+a_{1} \alpha+\cdots+a_{d-1} \alpha^{d-1}\right\}$ is a field, and that $F(\alpha)$ is the smallest subfield of $K$ containing $F$ and $\alpha$.

This was followed by giving the following important definition.
Definition. Let $F \subseteq K$ be field, and $\alpha \in K$. Then $\alpha$ is algebraic over $F$ if $\alpha$ is a root of a polynomial with coefficients in $F$. It follows then that $\alpha$ also has a minimal polynomial over $F$.
We also noted that the minimal polynomial of $\alpha$ over $F$ depends on $F$ : If we set $F:=\mathbb{Q}(\sqrt{2}), \alpha:=\sqrt[4]{2}$, and $K:=\mathbb{R}$, then the minimal polynomial of $\alpha$ over $\mathbb{Q}$ is $x^{4}-2$, while the minimal polynomial of $\alpha$ over $F$ is $x^{2}-\sqrt{2}$.

We ended class by defining $F(\alpha)$ in the case $\alpha$ is not algebraic over $F$.
Definition. Suppose $F \subseteq K$ are fields, and $\alpha \in K$ is not algebraic over $F$. We set $F(\alpha)$ equal to the intersections of all intermediate fields $F \subseteq E \subseteq K$ such that $\alpha \in E$.

We then noted that for $\alpha$ in the definition above,

$$
F(\alpha)=\left\{f(\alpha) g(\alpha)^{-1} \in K \mid f(x), g(x) \in F[x]\right\}=\left\{\left.\frac{f(\alpha)}{g(\alpha)} \in K \right\rvert\, f(x), g(x) \in F[x]\right\}
$$

Monday, March 27. We began by defining a field $F$ to be a commutative ring in which every non-zero element has a multiplicative inverse. We noted the following familiar examples of fields: $\mathbb{Q}, \mathbb{R}, \mathbb{C}$, and $\mathbb{Z}_{p}, p$ a prime. We also noted rings like $\mathbb{Z}$ and $\mathbb{Q}[x]$ are not fields. We then constructed several other examples of fields, including the following:
(i) $\mathbb{Q}(\sqrt{2}):=\{a+b \sqrt{2} \mid a, b \in \mathbb{Q}\}$, a subfield of $\mathbb{R}$.
(ii) $\mathbb{Q}(i):=\{a+b i \mid a, b \in \mathbb{Q}\}$, a subfield of $\mathbb{C}$.
(iii) $\mathbb{Q}(\sqrt[3]{2}):=\{a+b \sqrt[3]{2}+c \sqrt[3]{4} \mid a, b, c \in \mathbb{Q}\}$, a subfield of $\mathbb{R}$.

For (iii), we noted that $E:=\{a+b \sqrt[3]{2} \mid a, b \in \mathbb{Q}\}$ is not a field, since $\sqrt[3]{2} \cdot \sqrt[3]{2}$ does not belong to $E$. The proof of this required the observation that $x^{3}-2$ is irreducible over $\mathbb{Q}$. Using Bezout's Principle in $\mathbb{Q}[x]$, we were able to show that non-zero elements in $\mathbb{Q}(\sqrt[3]{2})$ as defined in (iii) have multiplicative inverses in $\mathbb{Q}(\sqrt[3]{2})$, which is crucial for this set to be a field.

We ended class by noting that if $F \subseteq K$ are fields, then $K$ can be regarded as a vector space over $F$. As such, we refer to the dimension of $K$ as an $F$ vector space as the degree of $K$ over $F$, which we denote by $[K: F]$.
Friday, March 24. We began by proving the following proposition.
Proposition. Let $R$ be a commutative ring.
(i) An ideal $P \subseteq R$ is a prime ideal if and only if $R / P$ is an integral domain.
(ii) An ideal $M \subseteq R$ is a maximal ideal if and only if $R / M$ is a field.

We then defined a commutative ring to be Noetherian if it satisfies any one, and hence all, of the conditions in the next proposition.
Proposition. The following conditions are equivalent for the commutative ring $R$ :
(i) $R$ satisfies the ascending chain condition, i.e., given a chain of ideals $I_{1} \subseteq I_{2} \subseteq \cdots$, there exists $n_{0}$ such that for all $n \geq n_{0}, I_{n_{0}}=I_{n}$.
(ii) $R$ satisfies the maximal condition, i.e., any non-empty collection of ideals of $R$ has a maximal element.
(iii) Every ideal of $R$ is finitely generated.

We also noted that the Noetherian condition can be stated for non-commutative rings, but one must define the notions of left Noetherian and right Noetherian by restricting to left or right ideals in the previous proposition.

We then stated and proved the celebrated:
Hilbert's Basis Theorem. Let $R$ be a Noetherian commutative ring. Then $R[x]$, the polynomial ring in one variable over $R$, is Noetherian.
We ended clss by noting that two immediate consequences of the Hilbert Basis Theorem are that the polynomial rings $F\left[x_{1}, \ldots, x_{n}\right]$, with $F$ a field, and $\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right]$ are Noetherian. Since a homomorphic image of a Noetherian ring is clearly Noetherian, it followed that homomorphic images of these polynomial rings are Noetherian, which implies that most of the rings encountered in algebraic geometry and algebraic number theory are Noetherian. In particular, since any Noetherian ring automatically satisfies the ascending chain condition on principal ideals, it follows that every non-zero, non-unit in a Noetherian domain can be factored as a product of irreducible elements, a fact that applies to the rings in algebraic geometry and algebraic number theory. However, such rings need not be UFDs - which shows that it is the uniqueness of factorization that eludes the non-UFDs encountered in algebraic geometry and algebraic number theory.
Wednesday, March 22. We worked through the details showing that $R:=\mathbb{R}[x, y] /\left\langle x^{2}+y^{2}-1\right\rangle$, is not a UFD. This fact follows because the image of $x$ in $R$ is an irreducible element that is not a prime element. Crucial steps in the proof required showing:
(i) $R$ is an integral domain. Here we used Eisenstein's criterion in the exact same way as in the previous lecture.
(ii) $\bar{x}$, the image of $x$ in $R$, is not a prime element. This followed since $R /\langle\bar{x}\rangle \cong \mathbb{R}[y] /\left\langle y^{2}-1\right\rangle$, which is not an integral domain.
(iii) Every polynomial in $f \in \mathbb{R}[x, y]$ can be written uniquely as $f=f_{0}+f_{1}+\cdots+f_{n}$, where each $f_{j} \in \mathbb{R}[x, y]$ is homogeneous of degree $j$.
(iv) $x^{2}+y^{2}$ is an irreducible polynomial in $\mathbb{R}[x, y]$.
(v) $\bar{x} \in R$ is an irreducible element.

For part (iv), we needed the following lemma:
Lemma. Suppose $f \in \mathbb{R}[x, y$ [ is homogeneous of degree two and $f$ is not irreducible. Then $f=g h$, where $h \in \mathbb{R}[x, y]$ are homogeneous of degree one. To see this, one writes $g=g_{n}+\cdots+g_{0}$ and $h=h_{m}+\cdots+h_{0}$, where each $g_{i}, h_{j}$ are homogeneous of degrees $i$ and $j$. Since $f$ has no terms of degree zero or one, we must have $g_{0}=g_{1}=h_{1}=h_{0}=0$. On the other hand, $f$ has no terms of degree greater than two, so we must have $n=m=1$, which gives what we want.
We then addressed the question: Why doesn't this same approach show that $S:=\mathbb{C}[x, y] /\left\langle x^{2}+y^{2}-1\right\rangle$ is not a UFD? The answer is that while the image of $x$ in $S$ is still not prime, the image of $x$ in $S$ is no longer irreducible. In fact, in $S$ we have

$$
x \equiv \frac{1}{2}(x-i y) \cdot\{(x+i y+i)(x+i y-i)\}
$$

where $x-i y$ is a unit in $S$, but neither $x+i y+i$ nor $x+i y-i$ are units in $S$.
Monday, March 20. We worked through the details showing that the ring $R:=\mathbb{C}[x, y] /\left\langle x^{2}+y^{2}-1\right\rangle$ is a UFD. Crucial steps in the proof required showing:
(i) $R$ is an integral domain. For this, we first proved Eisenstein's criterion and applied it to the polynomial $x^{2}+y^{2}-1$.
(ii) If $p$ is a prime element in the integral domain $S$, then $S /\langle p\rangle$ is an integral domains (and conversely).
(iii) For an indeterminate $U$, the Laurent polynomial ring $\mathbb{C}\left[U, U^{-1}\right]$ is a UFD. This followed from a property of localization, which is part of the Bonus problem on Exam 2.
(iv) $R \cong \mathbb{C}[U, V] /\langle U V-1\rangle$.
(vi) $\mathbb{C}[U, V] /\langle U V-1\rangle \cong \mathbb{C}\left[U, U^{-1}\right]$. This followed by showing that the ring homomorphism $\phi: \mathbb{C}[U, V] \rightarrow \mathbb{C}\left[U, U^{-1}\right]$ given by $\phi(f(U, V))=f\left(U, U^{-1}\right)$ is surjective with kernel equal to $\langle U V-1\rangle$.

Friday, March 10. Today we did group work on homework problems.
Wednesday, March 8. We continued along our path leading to the one of the main theorems in the course: If $R$ is a UFD, then $R[x]$ is a UFD. We began by recalling the quotient field of an integral domain as well as Proposition A and Gauss's Lemma from the previous lecture. We also observed that if $R$ is a UFD with quotient field $K$, then any $h(x) \in K[x]$ can be written as $\frac{a}{b} \cdot h_{0}(x)$, where $a, b \in R$ have o common factor and $h_{0}(x) \in R[x]$ is a primitive polynomial. We then proceeded to prove the following two propositions.
Proposition B. Suppose $R$ is a UFD with quotient field $K$ and $f(x) \in R[x]$ is primitive. Then $f(x)$ is irreducible in $R[x]$ if and only if it is irreducible in $K[x]$.
Proposition C. Suppose $R$ is a UFD and $f(x) \in R[x]$ is primitive and irreducible. Then $f(x)$ is a prime element.

After presenting these propositions we were able to give a proof of the following theorem.
Theorem. If $R$ is a UFD, then $R[x]$, the polynomial ring in one variable over $R$, is also a UFD.
The idea behind the proof of the theorem was the following: Given $f(x) \in R[x]$ we may write $f(x)=a f_{0}(x)$, where $a \in R$ and $f_{0}(x) \in R[x]$ is primitive. Since $a$ can be factored as a product of primes, which remain prime in $R[x]$ (by Proposition A), it suffices to show that $f_{0}(x)$ is a product of primes in $R[x]$ and for this (by Proposition C), it suffices to show that $f_{0}(x)$ is a product of irreducible primitive polynomials in $R[x]$. This followed by factoring $f_{0}(x)$ as a product of irreducible polynomials in $K[x]$ and then using the observation stated in the first paragraph above together with Proposition B. We ended class by noting that, by induction, the following rings are UFDs: $F\left[x_{1}, \ldots, x_{n}\right], F$ any field; $\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right] ; R\left[x_{1}, \ldots, x_{n}\right]$, for $R$ a UFD.
Monday, March 6. We began with the construction the quotient field $K$ of an arbitrary integral domain $R$, as described in Homework 16. We noted that (up to isomorphism) $R$ can be identified with a subring of $K$ and that $K$ is the smallest field containing $R$, in the sense that $K$ is contained in any field containing $R$.

We then began a discussion of the basic strategy for proving one of the main results of this part of the course, namely if $R$ is a UFD, then $R[x]$ is a UFD. The idea is to use in tandem the facts that $R$ is a UFD and $K[x]$ is a UFD, for $K$ the quotient field of $R$. Illustrating the idea with the ring $\mathbb{Z}[x]$, we noted that $f(x) \in \mathbb{Z}[x]$ can be written as $a f_{0}(x)$, where $a \in \mathbb{Z}$ and $f_{0}(x) \in \mathbb{Z}[x]$ has the property that there is no common
divisor among the coefficients of $f_{0}(x)$. We then pointed out that on the one hand, we will show that the prime factorization of $a$ in $R$ remains a prime factorization of $a$ as an element of $R[x]$, while on the other hand, the factorization of $f_{0}(x)$ in $K[x]$ as a product of irreducible polynomials is actually a factorization of $f_{0}(x)$ in $R[x]$ as a product of irreducible elements.

Before setting out on our path towards the proof of the theorem, we began with the following observation. For $0 \neq a \in$ and $f(x) \in R[x], a \mid f(x)$ in $R[x]$ if and only if, in $R, a$ divides every coefficient of $f(x)$. We then proved the following fact (one of many versions of Gauss's lemma):

Proposition A. For a UFD $R$, if $p \in R$ is a prime element, then $p$ is also a prime element in $R[x]$.
We then defined $f(x) \in R[x]$ to be a primitive polynomial, if for all prime elements $p \in R, p \nmid f(x)$ in $R[x]$. This immediately gave rise to the following more common version of Gauss's lemma:

Gauss's lemma. Let $R$ be a UFD. Then the product of primitive polynomials is primitive.
Friday, March 3. We began class by noting the following properties relating division of elements to containments of principal ideals. For an integral domain $R$ :
(i) $a \mid b$ if and only if $\alpha b\rangle \subseteq \alpha a\rangle$.
(ii) $\alpha a\rangle=\alpha b\rangle$ if and only if $b=a u$, for some unit $u \in R$.
(iii) $q \in R$ is irreducible if and only $\alpha q\rangle$ is maximal among principal ideals.
(iv) $p \in R$ is prime if and only if whenever $a b \in \alpha p\rangle, a \in \alpha p\rangle$ or $b \in \alpha p\rangle$.

We then proved the following sequences of propositions.
Proposition A. For an integral domain $R$, consider the following statements:
(i) $R$ satisfies the ascending chain condition on principal ideals.
(ii) Every non-empty collection of principal ideals has a maximal element.
(iii) Every non-zero, non-unit in $R$ is a product of finitely many irreducible elements.

Then statements (i) and (ii) are equivalent, and imply statement (iii).
Proposition B. Let $R$ be a PID. Then $R$ satisfies the ascending chain condition on principal ideals.
Proposition C. Let $R$ be a PID. Then every irreducible element is a prime element.
We ended class by noting that it follows immediately from Propositions B and C, that every PID is a UFD.
Wednesday, March 1. We continued our discussion of integral domains, by first reviewing the crucial definitions of prime element and irreducible element. We then gave a proof of the following
Proposition. Let $R$ be an integral domain. Then the following are equivalent:
(a) Every non-zero, non-unit of $R$ can be written as a product of prime elements.
(b) Every non-zero, non-unit in $R$ can be written uniquely (up to order and unit multiples) as a product of irreducible elements.
We noted that a key point in the proof of the proposition was that irreducible elements are prime in the presence of either condition (a) or (b). We then defined a Unique Factorization Domain, or UFD, to be any ring satisfying the conditions of the previous proposition. We noted again that any ring with a division algorithm is a UFD, this includes such rings as $\mathbb{Z}, F[x]$ and the Gaussian integers $\mathbb{Z}[i]:=\{a+b i \mid a, b \in \mathbb{Z}\}$. We also indicated that the UFD properly is rather subtle by noting the following examples:
(i) The ring $\mathbb{R}[x, y] /\left\langle x^{2}+y^{2}-1\right\rangle$ is not a UFD, while the ring $\mathbb{C}[x, y] /\left\langle x^{2}+y^{2}-1\right\rangle$ is a UFD.
(ii) The ring $\mathbb{R}[x, y, z] /\left\langle x^{2}+y^{2}+z^{2}-1\right\rangle$ is a UFD, while the ring $\mathbb{C}[x, y, z] /\left\langle x^{2}+y^{2}+z^{2}-1\right\rangle$ is not a UFD.

We ended class by defining what it means for $R$ to be a Principal ideal Domain (PID) and demonstrated how a ring with a division algorithm is a PID by proving that $F[x]$ is a PID for any field $F$. We also noted (but did not verify) the classical example of a ring which is a PID, but does not admit a division algorithm, namely $S:=\{a+b \omega \mid a, b \in \mathbb{Z}\}$, where $\omega:=\frac{1}{2} \cdot(1+\sqrt{-19})$.
Monday, February 27. We began class by introducing the type of rings we will be studying for the next few weeks: A commutative ring $R$ is said to be an integral domain (ID) if the product of non-zero elements is
always non-zero. For example $\mathbb{Z}$ and $F[x]$, with $F$ a field are examples of integral domains. We also noted that if $R$ is an integral domain, then $R[x]$, the polynomial ring over $R$, is also an integral domain - since in this case, if $f(x), g(x) \in R[x]$, then the degree of $f(x) g(x)$ is the sum of the degrees of $f(x)$ and $g(x)$. We also noted that $\mathbb{Z}_{4}$ and the ring of continuous functions from $\mathbb{R} \rightarrow \mathbb{R}$ are not integral domains. This was followed by noting

Cancellation holds in an integral domains. Suppose $R$ is an integral domain, and $a, b, c \in R$, with $a \neq 0$. If $a b=a c$, then $b=c$.

This was followed by giving a number of definitions for elements in an integral domain $R$, including the definitions of: $a \mid b$; a unit; associates; prime elements; irreducible elements. In particular, a non-zero, non-unit $p \in R$ is prime if whenever $p|a b, p| a$ or $p \mid b$ and $q \in R$ is an irreducible element if whenever $q=a b$, then $a$ or $b$ is a unit. We noted that primes are always irreducible, but irreducible elements need not be prime, an example being $3 \in \mathbb{Z}+\sqrt{-5} \mathbb{Z}$.

We ended class by proving the following proposition (whose conclusion was expected, based upon our proof of the uniqueness part of the Fundamental Theorem of Arithmetic) and discussing its relevance to a general theory of unique factorization
Proposition. Let $R$ be an integral domain and $p_{1}, \ldots, p_{r} q_{1}, \ldots, q_{t}$ be primes in $R$. If $p_{1} \cdots p_{r}=q_{1} \cdots q_{s}$, then $r=s$ and after re-indexing, $q_{i}=p_{i} u_{i}$, for units $u_{i}$ and $1 \leq i \leq r$.
Friday, February 24. Today's lecture was devoted to a discussion concerning, and a proof of:
Fundamental Theorem of Arithmetic. Every positive integer $n$ can be written uniquely as as a product $n=p_{1}^{e_{1}} \cdots p_{r}^{e_{r}}$, where each $p_{i}$ is prime and $e_{i} \geq 1$. Here uniqueness mens, if $n=q_{1}^{f_{1}} \cdots q_{s}^{f_{s}}$, with each $q_{j}$ prime and $f_{i} \geq 1$, then $r=s$, and after re-indexing, we have $1 \leq i \leq r, q_{i}=p_{i}$ and $e_{i}=f_{i}$, for all $1 \leq i \leq r$.

The proof of the theorem required a few preliminary results:
(i) If $a, b \in \mathbb{Z}$, with $b>0$, then there exist unique $q, r \in \mathbb{Z}$ such that $0 \leq r<b$. (Division Algorithm).
(ii) Given $a, b \in \mathbb{Z}$, the greatest common divisor of $a$ and $b$ exists.
(iii) If $a, b \in \mathbb{Z}$ and $d=\operatorname{GCD}(a, b)$, then there exist $n, m \in \mathbb{Z}$ such that $d=m a+n b$. (Bezout's Principle)

The division algorithm was proven by finding a least positive integer $r$ among all expressions of the form $a-b k$, with $k \in \mathbb{Z}$ and $a-b k \geq 0$. We then used the division alrgoithm to show (ii) and (iii), the key observation being if $a, b \in \mathbb{Z}$ with $b>0$ and $a=b q-r$, with $0 \leq r<b$, then the set of common divisors of $a$ and $b$ is the same as the set foc common divisors of $b$ and $r$. We then used Bezout's principle to prove the following important properties of primes:

$$
\text { If } p, a, b \in \mathbb{Z} \text {, with } p \text { prime, and } p \mid a b \text {, then } p \mid a \text { or } p \mid b \text {. }
$$

We then gave a proof of the Fundamental Theorem of Arithmetic, by first showing existence of a factorization. This followed by noting that if factorization failed, the set $X$ of positive integers not admitting a factorization would have a least element, say $n$. Since $n$ would not be prime $n=a b$, with $a, b$ positive integers less than $n$, and thus $a, b \notin X$. Therefore each of $a, b$ have a prime factorization, and hence $a b$ has a factorization, a contradiction forcing $X$ to be empty. Uniqueness of factorization followed by induction and the important property of primes stated above.

We ended class by discussing how exactly the same steps used in proving the Fundamental Theorem of Arithmetic can be used to prove the following fact: For $F$ a field, every monic polynomial $f(x) \in F[x]$ can be factored uniquely as a product $f(x)=p_{1}(x)^{e_{1}} \cdots p_{r}(x)^{e_{r}}$, with each $p_{i}(x)$ monic and irreducible over $F$. This works because $F[x]$ also has a (familiar) division algorithm.

Wednesday, February 22. We continued our discussion of rings in the abstract, beginning with a discussion concerning ideals in the ring $R$ generated by the set $X \subseteq R$. We noted that if one defines the left ideal of $R$ generated by $X$, denoted $\langle X\rangle_{L}$, to be the intersection of all left ideals of $R$ containing $X$, then $\langle X\rangle_{L}$ equals the set of all finite expressions of the form $r_{1} x_{1}+\cdots+r_{n} x_{n}$., with each $r_{i} \in R$ and $x_{i} \in X$. We then defined the right ideal of $R$ generated by $X$ and the two-sided ideal of $R$ generated by $X$ in similar ways and gave corresponding intrinsic descriptions.

We then showed that if $I \subseteq R$ is a two-sided ideal, the abelian group $(R / I,+)$ has a natural ring structure, where coset multiplication is defined as $(a+I) \cdot(b+I):=a b+I$. After defining what it means for a map $f: R \rightarrow S$ between rings to be a ring homomorphism, we noted that a subset $I \subseteq R$ is a two-sided ideal if and only if $I$ is the kernel of a ring homomorphism. We then noted that one has essentially the same isomorphism theorems for rings as for groups, in particular
(i) If $f: R \rightarrow S$ is a surjective ring homomorphism, then $R / \operatorname{ker}(f) \cong S$.
(ii) If $f: R \rightarrow S$ is a surjective ring homomorphism, then there is a one-to-one correspondence between the two-sided (resp., left or right) ideals of $S$ and the two-sided (resp., left or right) ideals of $R$ containing $\operatorname{ker}(f)$.

We ended class by discussing the Fundamental Theorem of Arithmetic for $\mathbb{Z}$ and some of its more subtle points, especially what uniqueness of factorization means when we factor a positive integer uniquely into a product of primes.

Monday, February 20. We began our discussion of ring theory by defining a (not necessarily commutative) ring as a set with two binary operations,+ such that:
(i) $(R,+)$ is an abelian group.
(ii) Multiplication is associative.
(iii) $a \cdot(b+c)=a \cdot b+a \cdot c$ and $(b+c) \cdot a=b \cdot a+c \cdot a$, for all $a, b, c \in R$.
(iv) $R$ has a multiplicative identity, denoted as 1 , satisfying $1 \cdot a=a=a \cdot 1$, for all $a \in R$.

We then gave several examples of rings, including $\mathbb{Z}, \mathbb{R}, \mathbb{C}, \mathbb{Q}$, rings of functions on a set taking values in a commutative ring, and matrix rings. This was followed by a discussion of various types of ideals in a ring: left ideals, right ideals, and two-sided ideals. Examples of each were given. We noted that for $R$, the ring of $2 \times 2$ matrices over $\mathbb{R}$ :
(i) The set of matrices of the form $\left(\begin{array}{ll}a & 0 \\ b & 0\end{array}\right)$ is a left ideal.
(ii) The set of matrices of the form $\left(\begin{array}{ll}a & b \\ 0 & 0\end{array}\right)$ is a right ideal
(iii) The only non-zero two-sided ideal has to be the whole ring.

We also noted that if $R=\mathrm{M}_{n}(\mathbb{Z})$, then the matrices whose entries are even integers gives an example of a proper two-sided ideal in $R$. Similarly, $\mathrm{M}_{2}(n \mathbb{Z})$ is a two=sided ideal of $R$. This ultimately led to us recording the fact that if $S$ is a commutative ring, then the two-sided ideals of $R:=\mathrm{M}_{n}(S)$ are all of the form $\mathrm{M}_{n}(J)$, where $J \subseteq S$ is an ideal.

Friday, February 17. Today's class was devoted to a proof of the following theorem.
Theorem. Let $G$ be a simple group of order 60. Then $G$ is isomorphic to $A_{5}$.
Proof. Suppose we could find a group homomorphism from $G$ to $S_{5}$. Since $G$ is simple, the kernel of this homomorphism is $\{e\}$. Thus, $G$ is isomorphic to a subgroup $G^{\prime} \subseteq S_{5}$. Since $\left|G^{\prime}\right|=60,\left[S_{5}: G^{\prime}\right]=2$, so that $G^{\prime}$ is a proper normal subgroup of $S_{5}$. By Homework 8 or the Bonus Problem on Exam $1, G^{\prime}=A_{5}$, so that $G$ is isomorphic to $A_{5}$, as required.

Thus, we seek a group homomorphism $\psi: G \rightarrow S_{5}$. For this recall that we have seen in class that such a homomorphism exists, if $G$ acts on a set with 5 elements. We first note that $60=2^{2} \cdot 3 \cdot 5$, so that any Sylow 2 -subgroup has four elements, any Sylow 3 -subgroup has three elements and any Sylow 5 -subgroup has five elements. By the Third Sylow Theorem, for $p=2,3,5$, the number of $p$-Sylow subgroups divides 60 and is congruent to 1 modulo $p$. Based upon this we have:

Possible number of Sylow 2-subgroups $=1,3,5,15$
Possible number of Sylow 3 -subgroups $=1,4,10$
Possible number of Sylow 5 -subgroups $=1,6$
In each case, we can eliminate the possibility of one Sylow subgroup, for then that subgroup would be normal in $G$, contradicting the simplicity of $G$. Suppose we had four Sylow 3-subgroups. Then if $G$ acts on the set of Sylow 3 -subgroups via conjugation, and thus we would have a group homomorphism $\psi: G \rightarrow S_{4}$. Since
$\left|S_{4}\right|=24, \operatorname{ker}(\psi) \neq\{e\}$. But then $G$ would have a non-trivial normal subgroup, which is a contradiction. Thus, $G$ must have ten Sylow 3-subgroups. Now, any two of the Sylow 3-subgroup have to intersect in e, since the intersection would is a subgroup of each. Thus, there are $10 \cdot 2=20$ elements in $G$ having order 3.
Since we cannot have one Sylow 5 -subgroup, there must be six of them, and since they each have prime order, the intersection of any two of them must be $e$. Thus, there are $6 \cdot 4=24$ elements of order 5 in $G$. We have now accounted for $20+24=44$ elements of $G$.

Regarding the number of Sylow 2-subgroups, we cannot have just three of them, otherwise, there would exist a group homomorphism $\psi: G \rightarrow S_{3}$, which would necessarily have a non-trivial kernel. If the number of Sylow 2-subgroups equals five, then we let $G$ act on the set of Sylow 2-subgroups via conjugation. This then gives a group homomorphism $\psi: G \rightarrow S_{5}$, which is what we want.
Suppose the number of Sylow 2-subgroups equals fifteen. First assume the following: Given any two such subgroups, say, $P_{1}$ and $P_{2}, P_{1} \cap P_{2}=e$. Then $G$ would contain $15 \cdot 3=45$ elements of order two or four, which is a contradiction, since we already have 44 elements of order three or five. Thus, there must be at least two $P_{1}, P_{2}$ with $P_{1} \cap P_{2} \neq e$. Then $\left|P_{1} \cap P_{2}\right|=2$. If $e \neq x \in P_{1} \cap P_{2}$, then since $P_{1}$ and $P_{2}$ are abelian, $x$ commutes with every element in both $P_{1}$ and $P_{2}$. Thus, both $P_{1}$ and $P_{2}$ are contained in $C_{G}(x)$, the centralizer of $x$. Since $\left|P_{1} P_{2}\right|=8$, and $P_{1} P_{2} \subseteq C_{G}(x)$ (as sets), $\left|C_{G}(x)\right| \geq 8$ and $\left|C_{G}(x)\right|$ divides 60 . Moreover, since $P_{1} \subseteq C_{G}(x)$, 4 divides $\left|C_{G}(x)\right|$. Thus, $\left|C_{G}(x)\right|=12$ or 20.

Suppose $\left|C_{G}(x)\right|=12$. Then $\left[G: C_{G}(x)\right]=5$. Thus, if we let $G$ act on the set of left cosets of $C_{G}(x)$ in $G$ via translation, there exists a group homomorphism $\psi: G \rightarrow S_{5}$, which is what we want. If $\left|C_{G}(x)\right|=20$, then, in a similar way, there would exist $\psi: G \rightarrow S_{3}$, which would have non-trivial kernel, so this case does not exist. Thus, if the number of Sylow 2-subgroups equals fifteen, then the required map $\psi: G \rightarrow S_{5}$ exists and this completes the proof of the theorem.

We ended class by very informally discussing the program resulting in the classification finite simple groups, noting that centralizers of involutions play a key role in many of the resulting theorems, as they did in the proof above.

Wednesday, February 15. We began by stating the three Sylow theorems. Here are the second and third Sylow theorem (the first Sylow theorem is stated in the previous update).
Second Sylow Theorem. Let $G$ be a finite group such that $|G|=p^{n} m$, where $p$ is prime and $p$ does not divide $m$. Suppose $H \subseteq G$ is a subgroup of order $p^{i}$, with $1 \leq i \leq n$ and $P$ is a Sylow p-subgroup. Then there exists $a \in G$ such that $H \subseteq a P a^{-1}$. In particular, any two Sylow p-subgroups are conjugate.

Third Sylow Theorem. Let $G$ be a finite group such that $|G|=p^{n} m$, where $p$ is prime and $p$ does not divide $m$ and write $n_{p}$ for the number of Sylow $p$-subgroups. Then $n_{p}$ divides $|G|$ and is congruent to 1 mod $p$.

We then gave applications of these theorems by showing any group of order 36 or 105 has a non-trivial normal subgroup, We then provided proofs of the second and third Sylow theorems. The proofs of these theorems relied on the following lemma:
Lemma. Let $G$ be a group of order $p^{t}$, with $p$ prime, and assume $G$ acts on the finite set $X$. If $r$ denotes the number of orbits with just one element, then $|X| \equiv r(\bmod p)$.

The proof of the second theorem followed by letting $H$ act on $X$, the left cosets of some Sylow $p$-subgroup $P$, by left translation. We then used the Lemma to show that that $r \neq 0$. Thus, there exists $g P \in X$ such that $h g P=g P$, for all $h \in h$, so $h g \in g P$, for all $h \in H$, which gives $H \subseteq g P g^{-1}$. The third Sylow theorem followed by first letting $G$ act on the set $X$ of Sylow $p$-subgroups, so $n_{p}=|X|$, in this case. By the second Sylow theorem, orb $(P)=X$, for any $P \in X$, so $n_{p}=|\operatorname{orb}(P)|$ divides $G$. Then letting $P$ act on $X$, we noted that $P$ is the only element of $X$ having one element, so by the Lemma, $n_{p}=|X| \equiv 1(\bmod p)$.
Monday, February 13. We stated and proved Cauchy's Theorem for abelian groups: If $G$ is a finite abelian group, and $p$ is a prime dividing $|G|$, then $G$ has an element of order $p$. Equivalently, $G$ has a subgroup of order $p$, namely the cyclic subgroup generated by an element of order $p$. We then defined the notion of a Sylow $p$-subgroup for a finite group whose order is divisible by $p$. We then stated and proved:

First Sylow Theorem. Let $G$ be a finite group such that $|G|=p^{n} m$, where $p$ is prime and $p$ does not divide $m$. Then $G$ has a Sylow p-subgroup. That is, there exists a subgroup $P \subseteq G$ such that $|P|=p^{n}$.
The idea of the proof was the following: If $|G|=p$, the result is clear. Proceed by induction on $|G|$. Using the class equation, if $p$ does not divide $|Z(G)|$, then $p$ does not divide the index of some $C_{G}\left(x_{i}\right)$, where $x_{i}$ is an element whose conjugacy class contains more than one element. It follows that $p^{n}$ divides $\left|C_{G}\left(x_{i}\right)\right|$, so by induction, $C_{G}\left(x_{i}\right)$, and hence $G$, has a subgroup of order $p^{n}$. Otherwise, if $p$ divides $|Z(G)|$, by Cauchy's theorem, there exists $x \in Z(G)$, with $o(x)=p$. One finishes the proof by applying induction to $G /\langle x\rangle$ and using the group correspondence theorem. We then recorded two consequences of this theorem:
(i) If $|G|=p^{n} m$, as in Sylow's theorem, then for each $1 \leq i \leq n$, there exist subgroups $H_{1} \subseteq \cdots \subseteq H_{n}$ such that $\left|H_{i}\right|=p^{i}$.
(ii) If $|G|=p q^{n}$, with $p<q$ primes, then $G$ has a normal Sylow $q$-subgroup (which is the unique Sylow $q$-subgroup).

The first statement follows from the Sylow theorem and the property proved in the previous lecture about groups of order $p^{n}$ and the second statement follows from the Sylow theorem and the fact that a subgroup whose index is the smallest prime dividing the order of the group is normal.
Friday, February 10. Today we did group work on homework problems.
Wednesday, February 8. We began class by recalling the basic ideas behind a group $G$ acting on a set $X$, in particular, the important fact that if $G$ acts on a set with $n$ elements, then there is a group homomorphism fro $G$ to $S_{n}$. This led to a discussion involving the orbit of $x \in X$, i.e., $\operatorname{orb}(x):=\{g \cdot x \mid g \in G\}$ and the stabilizer of $x$, i.e., $G_{x}:=\{g \in G \mid g \cdot x=x\}$. We also noted that the distinct orbits under the action partition $X$. We proved the:
Proposition. Assume the group $G$ acts on the set $X$. Fix $x \in X$. The there is a 1-1, onto set map between $\operatorname{orb}(x)$ and the set of distinct left cosets of $G_{x}$ given by $g \cdot x \longrightarrow g G_{x}$. In particular, if $|\operatorname{orb}(x)|$ or $\left[G: G_{x}\right]$ is finite, then $|\operatorname{orb}(x)|=\left[G: G_{x}\right]$, and it follows that $|\operatorname{orb}(x)|$ divides $|G|$, if $G$ is finite.
We then noted that in the special case $G$ acts on itself via conjugation, orb $(x)=\left\{g x g^{-1} \mid g \in G\right\}$, the conjugacy class of $G$ and $G_{x}:=\{g \in G \mid g x=x g\}$, the centralizer if $X$. We denoted the conjugacy class of $x$ by $c(x)$ and the centralizer of $x$ by $C_{G}(x)$. Thus, $\left.\mid c(x)\right) \mid=\left[G: C_{G}(x)\right]$, when either one of these is finite. We were thus led to a proof of the very important:

Class Equation. Let $G$ be a finite group. Then:

$$
\begin{aligned}
|G| & =|Z(G)|+\Sigma_{i=1}^{r}\left|c\left(x_{i}\right)\right| \\
& =|Z(G)|+\Sigma_{i=1}^{r}\left[G: C_{G}\left(x_{i}\right)\right]
\end{aligned}
$$

where the sum is taken over the distinct conjugacy classes with more than one element. Here $Z(G)$ denotes the center of $G$, where $Z(G):=\{g \in G \mid g x=x g$, for all $x \in G\}$.
Using the class equation, we were able to prove the following theorem.
Theorem. Let $G$ be a finite group with $|G|=p^{n}$, with $p$ prime and $n \geq 1$. Then:
(i) $Z(G) \neq\{e\}$
(ii) For each $1 \leq i<n, G$ has a subgroup of order $p^{i}$.

Part (i) followed immediately from the class equation, since each $\left[G: C_{G}\left(x_{i}\right)\right]$ divides $p^{n}$, while part (ii) follows by induction on $n$, applied to $G /\langle y\rangle$, where $y \in G$ is an element of order $p$. We ended class by noting that the proof of the theorem we gave can be slightly modified to show that a finite group of order $p^{n}$ is solvable, i.e., there exists a sequence of subgroups $\{e\} \subseteq H_{1} \subseteq \cdots \subseteq H_{n-1} \subseteq H_{n}=G$, such that each $H_{i}$ is a normal subgroup of $H_{i+1}$ and each quotient $H_{i+1} / H_{i}$ is cyclic of prime order.
Monday, February 6. Given a set $X$ and a group $G$, we defined what it means for $G$ to act on X: There is a binary map $G \times X \rightarrow X$ satisfying: (i) $e \cdot x=x$ and $(a b) \cdot x=a \cdot(b \cdot x)$, for all $a, b \in G$ and $x \in X$. We then discussed the following examples of group actions:

Examples Let $G$ be a group and $X$ a set:
(i) Taking $X=G$, then $G$ acts on itself via left translation: $g \cdot x:=g x$, for all $g \in G$ and $x \in X$.
(ii) If $H \subseteq G$ is a subgroup, and $X$ denotes the set of distinct left cosets of $H$, then $G$ acts on $X$ via left translation: $g \cdot(a H):=g a H$, for all $g \in G$ and $a H \in X$.
(iii) $G=S_{n}$ acts on $X=\{1,2 \ldots, n\}$, via $\sigma \cdot i:=\sigma(i)$, for all $\sigma \in S_{n}$ and $i \in X$.
(iv) Taking $X=G$, then $G$ acts on itself via conjugation: $g \cdot x:=g x g^{-1}$, for all $g \in G, x \in X$.
(v) Suppose $G$ has a subgroup of order $n$. If we let $X$ denote the set of all subgroups of order $n$, then $G$ acts on $X$ via conjugation: $g \cdot H=g H g^{-1}$, for all $g \in G$ and $H \in X$.
(vi) If $G:=\mathrm{Gl}_{n}(\mathbb{R})$ and $X:=\mathbb{R}^{n}$, thought of as column vectors, then matrix multiplication defines a group action: $A \cdot v:=A v$, for all $A \mid i n G$ and $v \in X$.

We then observed that if $G$ acts on $X$, then for fixed $g \in G$, the map $X \xrightarrow{g} X$ is one-to-one and onto. Thus, multiplication by $g$ permutes the elements of $X$. This gave rise to the very important:

Proposition. If the group $G$ acts on a set with $n$ elements, then there is a group homomorphism $\phi: G \rightarrow S_{n}$.
The important point of the proposition is that the association between $G$ and the elements of $S_{n}$ is given by a group homomorphism. We also noted that the converse is left as a homework exercise, namely, if there is a group homomorphism $\phi: G \rightarrow S_{n}$, then $\phi$ induces a group action on any set with $n$ elements. Hence: To give an action of the group $G$ on a set with $n$ elements is equivalent to giving a group homomorphism from $G$ to $S_{n}$.

As an application of the previous proposition and a number of other results established, we then proved the following theorem.

Theorem. Let $G$ be a finite group and $H \subseteq G$ a subgroup. Suppose $[G: H]=p$, where $p$ is the smallest prime dividing the order of $G$. Then $H$ is normal in $G$.

The ideal behind of the proof is there exists a group homomorphism $\phi: G \rightarrow S_{p}$ with kernel $K \subseteq H$. Thus, $G / K$ is isomorphic to a subgroup of $S_{p}$ and hence any prime $q$ dividing $|G / K|$ divides $p$ !. This forces $|G / K|=p$, which in turns gives $H=K$.

We finished class by noting that if $G$ acts on the set $X$, the relation $x_{1} \sim x_{2}$ if and only if $x_{2}=g \cdot x_{1}$, for some $g \in G$, with $x_{1}, x_{2} \in X$, is an equivalence relation on $X$. The equivalence class of $x \in X$ is called the orbit of $x$.

Friday, February 3. Today's lecture was devoted to the proof of the following very important theorem.
Theorem. $A_{n}$ is a simple group for $n \geq 5$. In other words, there are no proper normal subgroups of $A_{n}$, for $n \geq 5$.

Proof. The proof proceeds in four steps.
Step 1. We first note that $A_{n}$ is the subgroup of $S_{n}$ generated by the set of all 3-cycles. To see this, if ( $\left.u, v, w\right)$ is a 3 -cycle, then $(u, v, w)=(u, w)(u, v)$, so that every 3 -cycle belongs to $A_{n}$. Conversely, every element of $A_{n}$ is a product of permutations of the form $(u, v)(s, t)$ or $(u, v)(u, s)$, for distinct elements $u, v, w, s \in X_{n}$. But, $(u, v)(s, t)=(u, s, v)(u, s, t)$ and $(u, v)(u, s)=(u, s, v)$, which shows that $A_{n}$ is the subgroup of $S_{n}$ generated by the set of 3 -cycles.

Step 2. Fix $1 \leq a \neq b \leq n$. Then $A_{n}$ is generated by the 3-cycles of the form $(a, b, c)$ with $c \in X_{n} \backslash\{a, b\}$. To see this, let us note that any 3 -cycle is of the form: $(a, b, u),(a, u, b),(a, u, v),(b, u, v), o r(u, v, w)$, for $a, b, u, v, w$ distinct elements of $X_{n}$. However, direct calculation shows that

$$
\begin{aligned}
(a, u, b) & =(a, b, u)^{2} \\
(a, u, v) & =(a, b, v)(a, b, u)^{2} \\
(b, u, v) & =(a, b, v)^{2}(a, b, u) \\
(u, v, w) & =(a, b, u)^{2}(a, b, w)(a, b, v)^{2}(a, b, u)
\end{aligned}
$$

which gives what we want.

Step 3. If $N$ is a normal subgroup of $A_{n}$ and $N$ contains a 3 -cycle, then $N=A_{n}$. To see this, let $(a, b, u) \in N$ be a 3-cycle. Then, by Step 2, it suffices to show that $(a, b, c) \in N$, for all $c \in X_{n} \backslash\{a, b\}$. However, using that $(a, b)(u, c)$ is its own inverse, we have

$$
\{(a, b)(u, c)\}(a, b, u)^{2}\{(a, b)(u, c)\}^{-1}=(a, b)(u, c)(a, b, u)^{2}(a, b)(u, c)=(a, b, c)
$$

which belongs to $N$ since $N$ is normal in $A_{n}$.
Step 4. If $N$ is a normal subgroup of $A_{n}$, then $N$ contains a 3-cycle. This is the hardest step, and requires consideration of several cases, where we analyze the decomposition of $\sigma \in N$ into a product of disjoint cycles.

Case (a). $N$ contains an element $\sigma$ of the form $\sigma:=\left(i_{1}, \ldots, i_{k}\right) \tau$, where $k \geq 4$ and $\tau$ is a product of disjoint cycles, each of which is disjoint from $\left(i_{1}, \ldots, i_{k}\right)$. Set $\gamma:=\left(i_{1}, i_{2}, i_{3}\right)$. Then $\sigma^{-1}\left(\gamma \sigma \gamma^{-1}\right) \in N$. However,

$$
\sigma^{-1}\left(\gamma \sigma \gamma^{-1}\right)=\tau^{-1}\left(i_{1}, i_{k}, i_{k-1}, \ldots, i_{2}\right)\left(i_{1}, i_{2}, i_{3}\right)\left(i_{1}, \ldots, i_{k}\right) \tau\left(i_{1}, i_{3}, i_{2}\right)=\left(i_{1}, i_{3}, i_{k}\right)
$$

which shows that $N$ contains a 3 -cycle.
Case (b). $N$ contains $\sigma=\left(i_{1}, i_{2}, i_{3}\right)\left(i_{4}, i_{5}, i_{6}\right) \tau$, a product of disjoint cycles. Set $\gamma:=\left(i_{1}, i_{2}, i_{4}\right)$. Then $\sigma^{-1}\left(\gamma \sigma \gamma^{-1}\right) \in N$. However,

$$
\sigma^{-1}\left(\gamma \sigma \gamma^{-1}\right)=\tau^{-1}\left(i_{4}, i_{6}, i_{5}\right)\left(i_{1}, i_{3}, i_{2}\right)\left(i_{1}, i_{2}, i_{4}\right)\left(i_{1}, i_{2}, i_{3}\right)\left(i_{4}, i_{5}, i_{6}\right) \tau\left(i_{1}, i_{4}, i_{2}\right)=\left(i_{1}, i_{4}, i_{2}, i_{6}, i_{3}\right)
$$

and thus, $N$ contains a 5 -cycle. By Case (a), $N$ contains a 3 -cycle.
Case (c). $N$ contains $\sigma=\left(i_{1}, i_{2}, i_{3}\right) \tau$, where $\tau$ is a product of disjoint 2-cycles, disjoint from $\left(i_{1}, i_{2}, i_{3}\right)$. Then $\sigma^{2} \in N$, and since disjoint cycles commute,

$$
\sigma^{2}=\left(i_{1}, i_{2}, i_{3}\right)^{2} \tau^{2}=\left(i_{1}, i_{2}, i_{3}\right)^{2}=\left(i_{1}, i_{3}, i_{2}\right)
$$

so $N$ contains a 3 -cycle.
Case (d). One of the previous three cases must hold. If not, then every element of $N$ is a product of an even number of disjoint 2 -cycles. Let $\sigma \in N$, and write $\sigma=\left(i_{1}, i_{2}\right)\left(i_{3}, i_{4}\right) \tau$ be the cycle decomposition of $\sigma$. Set $\gamma:=\left(i_{1}, i_{2}, i_{3}\right)$, so that $\sigma^{-1}\left(\gamma \sigma \gamma^{-1}\right) \in N$. However,

$$
\sigma^{-1}\left(\gamma \sigma \gamma^{-1}\right)=\tau^{-1}\left(i_{3}, i_{4}\right)\left(i_{1}, i_{2}\right)\left(i_{1}, i_{2}, i_{3}\right)\left(i_{1}, i_{2}\right)\left(i_{3}, i_{4}\right) \tau\left(i_{1}, i_{3}, i_{2}\right)=\left(i_{1}, i_{3}\right)\left(i_{2}, i_{4}\right)
$$

For ease of notation, set $\alpha:=\left(i_{1}, i_{3}\right)\left(i_{2}, i_{4}\right) \in N$.
Since $n \geq 5$, there exists $j \in X_{n} \backslash\left\{i_{1}, \ldots, i_{4}\right\}$. Set $\beta:=\left(i_{1}, i_{3}, j\right) \in A_{n}$. Then, $\alpha \beta \alpha \beta^{-1} \in N$. However,

$$
\alpha \beta \alpha \beta^{-1}=\left(i_{1}, i_{3}\right)\left(i_{2}, i_{4}\right)\left(i_{1}, i_{3}, j\right)\left(i_{1}, i_{3}\right)\left(i_{2}, i_{4}\right)\left(i_{1}, j, i_{3}\right)=\left(i_{1}, i_{3}, j\right)
$$

showing that $N$ contains a 3 -cycle. But this is a contradiction, because any 3 -cycle is the product of two 2-cycles that are not disjoint.

All possibilities for cycle decompositions that can occur have been covered by the cases above, thus $N$ must contain a 3 -cycle. It follows immediately from Steps $1,2,3$ that $A_{n}$ cannot have a proper, normal subgroup, and therefore, $A_{n}$ is a simple group.

Wednesday, February 1. We reviewed some of the properties of elements in $S_{n}$, in particular, we showed that for a $k$-cycle $\sigma \in S_{n}$ and a fixed $i \in X_{n}$ :
(i) There exists a least positive integer $s \geq 1$ such that $\sigma^{s}(i)=i$.
(ii) $\left\{\sigma^{n}(i) \mid n \in \mathbb{Z}\right\}=\left\{i, \sigma(i), \ldots, \sigma^{s-1}(i)\right\}$.

We then presented the following theorem.
Theorem. Let $\sigma \in S_{n}$. Then:
(i) $\sigma$ can be written uniquely (up to order) as a product of disjoint cycles.
(ii) $\sigma$ can be written as a product of (not necessarily disjoint) 2-cycles.

After this we then proved the theorem stating that no permutation in $S_{n}$ can be written as a product of an even number of 2-cycles on the one hand, and a product of an odd number of 2-cycles on the other hand. The proof was based upon the following. Any $\sigma \in S_{n}$ corresponds to a matrix $A_{\sigma}$ obtained by permuting the rows of the $n \times n$ identity matrix according to $\sigma$. In other words, if $R_{1}, \ldots, R_{n}$ are the rows of the identity matrix, then $R_{\sigma(1)}, \ldots, R_{\sigma(n)}$ are the rows of $A_{\sigma}$. We then showed that for any $\sigma, \tau \in S_{n}, A_{\tau} A_{\sigma}=A_{\sigma \tau}$. For the sake of completeness, here is a proof of this latter fact: We first note that the $i$ th row of $A_{\sigma}$ has all entries equal to zero, except 1 in the $\sigma(i)$ th columns. Working over $\mathbb{Q}$, say, it follows that $A_{\sigma} \cdot\left(\begin{array}{c}a_{1} \\ \vdots \\ a_{n}\end{array}\right)=\left(\begin{array}{c}a_{\sigma(1)} \\ \vdots \\ a_{\sigma(n)}\end{array}\right)$.
Now set $b_{i}:=a_{\sigma(i)}$, for all $1 \leq i \leq n$. Then for $\tau \in S_{n}$, we have

$$
A_{\tau} A_{\sigma} \cdot\left(\begin{array}{c}
a_{1} \\
\vdots \\
a_{n}
\end{array}\right)=A_{\tau} \cdot\left(\begin{array}{c}
b_{1} \\
\vdots \\
b_{n}
\end{array}\right)=\left(\begin{array}{c}
b_{\tau(1)} \\
\vdots \\
b_{\tau(n)}
\end{array}\right)=\left(\begin{array}{c}
a_{\sigma(\tau(1))} \\
\vdots \\
a_{\sigma(\tau(n))}
\end{array}\right)=A_{\sigma \tau} \cdot\left(\begin{array}{c}
a_{1} \\
\vdots \\
a_{n}
\end{array}\right)
$$

Since this holds for all vectors $\left(\begin{array}{c}a_{1} \\ \vdots \\ a_{n}\end{array}\right)$, it follows that $A_{\tau} A_{\sigma}=A_{\sigma \tau}$.
From this, one has if $\sigma \in S_{n}$ and $\sigma=\tau_{1} \cdots \tau_{r}=\gamma_{1} \cdots \gamma_{s}$, where each $\tau_{i}, \gamma_{j}$ is a 2 -cycle, then, on the one hand,

$$
\begin{aligned}
\operatorname{det}\left(A_{\sigma}\right) & =\operatorname{det}\left(A_{\tau_{1} \cdots \tau_{r}}\right) \\
& =\operatorname{det}\left(A_{\tau_{r}} \cdots A_{\tau_{1}}\right) \\
& =\operatorname{det}\left(A_{\tau_{r}}\right) \cdots \operatorname{det}\left(A_{\tau_{1}}\right) \\
& =(-1)^{r}
\end{aligned}
$$

while on the other hand,

$$
\begin{aligned}
\operatorname{det}\left(A_{\sigma}\right) & =\operatorname{det}\left(A_{\gamma_{1} \cdots \gamma_{s}}\right) \\
& =\operatorname{det}\left(A_{\gamma_{s}} \cdots A_{\gamma_{1}}\right) \\
& =\operatorname{det}\left(A_{\gamma_{s}}\right) \cdots \operatorname{det}\left(A_{\gamma_{1}}\right) \\
& =(-1)^{s} .
\end{aligned}
$$

Thus, $(-1)^{r}=(-1)^{s}$, which shows that both $r$ and $s$ are even, or both $r$ and $s$ are odd.
Monday, January 30. We spent most of the class proving the following isomorphism theorems.
First Isomorphism Theorem. Let $\phi: G_{1} \rightarrow G_{2}$ be a surjective group homomorphism with kernel $K$. Then $G_{1} / K \cong G_{2}$.

Second Isomorphism Theorem. Let $K \subseteq N \subseteq G$ be groups such that $K$ and $N$ are normal in $G$. Then $N / K$ is a normal subgroup of $G / K$ and $(G / K) /(N / K) \cong G / N$.

After proving these theorems, we began a discussion of the structure of elements in the symmetric group $S_{n}$. We started with the definition of a $k$-cycle: $\tau \in S_{n}$ is a $k$-cycle if there exists $i_{1}, \ldots, i_{k} \in X=\{1,2, \ldots, n\}$ such that $\tau\left(i_{1}\right)=i_{2}, \tau\left(i_{2}\right)=i_{3}, \ldots, \tau\left(i_{k-1}\right)=i_{k}, \tau\left(i_{k}\right)=i_{1}$, and $\tau(j)=j$, if $j \notin\left\{i_{1}, \ldots, i_{k}\right\}$. We first observed that the order of a $k$-cycle is $k$. We then showed:
(i) Disjoint cycles commute.
(ii) If $\tau$ is a $k$-cycle and $\gamma$ is an $s$-cycle, and these cycles are disjoint, then the order of $\gamma \tau$ is the least common multiple of $r$ and $s$.

We also gave a heuristic proof that any permutation in $S_{n}$ is a product of disjoint cycles.
Friday, January 27. After recalling some basic consequences of the definition of group homomorphism, we spent the rest of the lecture proving the following theorem and its corollaries.

Theorem. Let $\phi: G_{1} \rightarrow G_{2}$ be a surjective group homomorphism with kernel $K$. Then there is a one-to-one, onto correspondence between the subgroups of $G_{1}$ containing $K$ and the subgroups of $G_{2}$ given by $H \longrightarrow \phi(H)$, for $H \subseteq G_{1}$ containing $K$, and $L \longrightarrow \phi^{-1}(L)$, for $L \subseteq G_{2}$. Under this correspondence, $\phi(H)$ is normal in $G_{2}$, if $H$ is normal in $G_{1}$ and $\phi^{-1}(L)$ is normal in $G_{1}$, if $L$ is normal in $G_{2}$.
Our first corollary noted that exactly the same theorem holds for $\phi$ not necessarily surjective if we replace $G_{2}$ by $\operatorname{im}(\phi)$. We then discussed at length the following corollary.
Corollary. Let $G$ be a group and $N$ a normal subgroup. Then there is a one-to-one, onto correspondence between the subgroups of $G$ containing $N$ and the subgroups of $G / N$. Under this correspondence, the normal subgroups of $G$ containing $N$ correspond to the normal subgroups of $G / N$.

We noted that this corollary follows from the theorem by taking $\phi: G \rightarrow G / N$ to be $\phi(g)=g N$, for all $g \in G$. In particular, we saw that any subgroup of $G / N$ is necessarily of the form $H / N$, for $H$ a subgroup of $G$ containing $N$.

Wednesday, January 25. We began class with the important observation that if $N$ is a normal subgroup of the group $G$, then for all $g \in G$ and $n \in N$, there exists $n^{\prime} \in N$ such that $g n=n^{\prime} g$. In other words, in the product $g n$ we can move $g$ passed $n$, at the expense of changing $n$ to $n^{\prime}$. We then showed that $G / N$ - the set of left cosets of $N$ - forms a group under coset multiplication by showing the product $g_{1} N g_{2} N=g_{1} g_{2} N$. The resulting group is called the quotient, or factor, group of $G$ by $N$, often referred to as $G \bmod N$. We then calculated group tables for the quotient groups $\mathbb{Z} / 2 \mathbb{Z}$ and $Q_{8} / K$, for $K=\{ \pm 1\}$.

We then defined the concept of a group homomorphism: The function $\phi: G_{1} \rightarrow G_{2}$ is a group homomorphism if $\phi(a b)=\phi(a) \phi(b)$, for all $a, b \in G_{1}$. We noted that $\phi\left(e_{1}\right)=e_{2}$ and $\phi\left(g^{-1}\right)=\phi(g)^{-1}$, for all $g \in G$. This was followed by defining the kernel of $\phi$ to be the set $\operatorname{ker}(\phi):=\left\{g \in G_{1} \mid \phi(g)=e_{2}\right\}$. We followed this by proving the following:
Proposition. Let $\phi: G_{1} \rightarrow G_{2}$ be a group homomorphism, with kernel $K$. Then:
(i) $K$ is a normal subgroup of $G_{1}$.
(ii) If $H$ is a subgroup of $G_{1}$, then $\phi(H)$ is a subgroup of $G_{2}$.

We ended class with the following remarks related to the proposition above: (i) If $H$ is a normal subgroup of $G_{1}$, then $\phi(H)$ need not be a normal subgroup of $G_{2}$, though it is, if $\phi$ is surjective and (ii) Not only is the kernel of a group homomorphism a normal subgroup, any normal subgroup $N$ in a group $G$ is the kernel of a group homomorphism, namely $N$ is the kernel of the homomorphism $\phi: G \rightarrow G / N$ defined by $\phi(g)=g N$.

Monday, January 23. Continuing our discussion of subgroups, we first noted that cancellation holds in any group $G$, i.e., $g x=g y$ implies $x=y$, for $g, x, y \in G$. We then used this to show that for any subgroup $H \subseteq G$ and $g \in G$, there is a 1-1 function from $H$ to $g H$, so that $|H|=|g H|$, when $|H|<\infty$ and similarly $|H|=|H g|$. Since the distinct left (or right) costs of $H$ partition $G$ we were able to immediately deduce:
LaGrange's Theorem. Let $G$ be a finite group and $H \subseteq G$ a subgroup. Then

$$
\begin{aligned}
|G| & =|H| \cdot(\text { number of distinct left cosets of } H) \\
& =|H| \cdot(\text { number of distinct right cosets of } H) .
\end{aligned}
$$

It followed that the number of distinct left cosets of $H$ equals the number of distinct right cosets of $H$, which we defined to be the index of $H$ in $G$, denoted $[G: H]$. We followed this by a discussion of normal subgroups, noting (but not proving) the equivalent conditions from Homework 2. This was followed by a discussion of simple groups (i.e., groups with no no-trivial normal subgroups) and a bit of the history of the classification of finite simple groups.
Examples. We then gave the following list of examples of normal subgroups:
(a) Any subgroup of an abelian group.
(b) Every subgroup of the quaternion group $Q_{8}:=\{ \pm 1, \pm i, \pm, j, \pm k\}$, though $Q_{8}$ is not abelian.
(c) The subgroup of $S_{3}$ generated by the 3 -cycle $(1,2,3)$.
(d) Any subgroup of index two.
(e) For $G:=\mathrm{Gl}_{n}(\mathbb{R}), N=\mathrm{Sl}_{n}(\mathbb{R})$, i.e., the set of $n \times n$ matrices over $\mathbb{R}$ with determinant one, is a normal subgroup $G$.

Friday, January 20. We began with the definition of a subgroup $H$ of a group $G$ : A subset $H \subseteq G$ is a subgroup if : (i) $H$ is closed under the binary operation of $G$ and (ii) $h \in H$ implies $h^{-1} \in H$. It followed easily from this that $H$ is a group in its own right under the binary operation of $G$. We then gave several examples of subgroups, including: $n \mathbb{Z} \subseteq \mathbb{Z} ; \mathbb{Z} \subseteq \mathbb{Q} ; \mathrm{Sl}_{n}(\mathbb{R}) \subseteq \mathrm{Gl}_{n}(\mathbb{R})$. Most importantly, we gave the following definition:

Definition. Let $X \subseteq G$ be a subset. Then the subgroup of $G$ generated by $X$, denoted by $\langle X\rangle$, is the set of all finite expressions of the form $x_{1}^{\epsilon_{1}} \cdots x_{r}^{\epsilon_{r}}$, where each $x_{i} \in X$ and $\epsilon \in\{ \pm 1,0\}$.
We quickly noted that, in fact, $\langle X\rangle$ is a subgroup of $G$. We also stated, but did not prove, that $\langle X\rangle$ is the intersection of the subgroups of $G$ containing $X$. An important case is when $X:=\{a\}$ consists of a single element. In this case we refer to $\langle a\rangle$ as the cyclic subgroup of $G$ generated by $a$.

Given a subgroup $H \subseteq G$, we defined the left coset $g H:=\{g h \mid h \in H\}$, for any $g \in G$. Right cosets were defined similarly: $H g:=\{h g \mid h \in H\}$. We then showed that left congruence modulo $H$ gives rise to an equivalence relation on $G$ whose equivalence classes are just the left cosets of $H$ in $G$. Here, we defined $a \equiv_{l} b(\bmod H)$ to mean $b^{-1} a \in H$. It followed that for any two left (or right) cosets $g_{1} H, g_{2} H$, either $g_{1} H=g_{2} H$ or $\left(g_{1} H\right) \cap\left(g_{2} H\right)=\emptyset$. This immediately implied that the distinct left (respectively, right) cosests of $H$ in $G$ partition $G$.
We ended class by explicitly calculating several cosets in $S_{3}$. Using the notation from the previous lecture, we set $H:=\langle\sigma\rangle$ and $K:=\langle\tau\rangle$. Then we calculated $e H, \tau H$, and $\sigma \tau H$, noting that $\tau H=H \tau, \tau H=\sigma \tau H$, and that $H, \tau H$ are the distinct left cosets of $H$. We then calculated $\sigma K, \sigma^{2} K$, and $K \sigma$, noting that, not only is $\sigma K \neq K \sigma$, but also that the right coset $K \sigma$ is not equal to any of the left cosets of $K$ in $G$.
Wednesday, January 18. We began by recalling the definition of a group (and an abelian group) and gave several examples of groups, including: $\mathbb{Z}, \mathbb{Z}_{n}$ under addition, $\mathbb{Z}_{n}^{*}$ under multiplication, where these elements of $\mathbb{Z}_{n}^{*}$ are the residue classes of elements in $\mathbb{Z}$ relatively prime to $n$. We also defined $S_{n}$, the symmetric group on $n$ objects and $D_{n}$, the dihedral group with $2 n$ elements and $\mathrm{Gl}_{n}(\mathbb{R})$, the $n \times n$ invertible matrices over $\mathbb{R}$.
We defined $S_{n}$ to be the set of 1-1, onto functions from the set $X:=\{1,2, \ldots, n\}$ to itself, with composition as the group operation. $D_{n}$ was informally defined as the group of symmetries of a regular $n$-sided polygon, and we noted, but did not prove, that $D_{n}$ is generated by a reflection about an axis of symmetry and a rotation of $\frac{2 \pi}{n}$ radians clockwise about the center of the polygon, and all possible products of these two symmetries. It followed that $\left|S_{n}\right|=n$ ! and $\left|D_{n}\right|=2 n$. We then introduced the notation $\sigma=\left(\begin{array}{cccc}1 & 2 & \cdots & n \\ i_{1} & i_{2} & \cdots & i_{n}\end{array}\right)$ for elements of $S_{n}$, meaning that $\sigma(j)=i_{j}$, for $1 \leq j \leq n$, and used this notation to calculate all of the elements of $S_{3}$. In particular, taking $\sigma:=\left(\begin{array}{lll}1 & 2 & 3 \\ 2 & 3 & 1\end{array}\right)$ and $\tau:=\left(\begin{array}{lll}1 & 2 & 3 \\ 2 & 1 & 3\end{array}\right)$, we calculated the products $\sigma^{2}, \sigma \tau, \sigma^{2} \tau$, noting that $\left\{e, \sigma, \sigma^{2}, \tau, \sigma \tau, \sigma^{2} \tau\right\}$ are the six distinct permutations of $\{1,2,3\}$ and thus $S_{3}=\left\{e, \sigma, \sigma^{2}, \tau, \sigma \tau, \sigma^{2} \tau\right\}$. We ended class by noting that one can calculate all possible products of elements from $S_{3}$ by using the identities $\sigma^{3}=e=\tau^{2}$ and $\tau \sigma=\sigma^{2} \tau$.

